## Jianxun Hu Changzheng Li Leonardo C. Mihalcea Editors

Schubert Calculus and Its Applications in Combinatorics and Representation Theory
Guangzhou, China, November 2017

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## Preface

With roots in enumerative geometry and Hilbert's 15th problem, modern Schubert Calculus studies classical and quantum intersection rings on spaces with symmetries, such as flag manifolds. The presence of symmetries leads to particularly rich structures, and it connects Schubert Calculus to many branches of mathematics, including algebraic geometry, combinatorics, representation theory, and theoretical physics. For instance, the study of the quantum cohomology ring of a Grassmann manifold combines all these areas in an organic way. The current volume showcases some of the newest developments in the subject, as presented at the "International Festival in Schubert Calculus", a conference held at Sun Yat-sen University in Guangzhou, China, during November 6-10, 2017.

The event included a 1-day mini-school and a 4-day international conference entitled "Trends in Schubert Calculus". There were over 80 participants, more than one half of which were international, from countries such as Australia, France, Germany, India, Japan, Korea, Poland, Russia, United Kingdom, and the U.S.A. This event continued the tradition of conferences in Schubert Calculus with large international participation; there were three such conferences in the past decade (Toronto 2010, Osaka 2012 and Bedlewo 2015).

The current volume includes 12 papers authored by some of the speakers, covering a large array of topics, including several high-quality surveys. Each of the papers was refereed by two anonymous experts in the field. This volume could not have existed without the combined efforts of the authors and referees, and we are grateful for everyone's contribution.

Problems with roots in classical Schubert Calculus attracted significant attention. The factorial Grothendieck polynomials, which investigate polynomials representing Schubert classes in K theory, were discussed in the paper by Matsumura and Sugimoto. The related problem of finding formulas for cohomology classes of various degeneracy loci is addressed in a paper by Darondeau and Pragacz. Yet another problem with roots in Schubert's classical work, that of finding formulas for the order of contact between manifolds, is investigated in the paper by Domitrz, Mormul, and Pragacz. Finally, Duan and Zhao address and survey Schubert's
classical problem of characteristics, and find presentations for the integral cohomology rings of flag manifolds, including those of exceptional Lie types.

There are rich connections between Schubert Calculus and the combinatorics of symmetric functions and polynomials. A survey by Pechenik and Searles highlights the properties of some of the most important bases of polynomials relevant for geometry. Expanding on this, a topic of current high interest is to relate and apply Schubert Calculus methods to problems in (combinatorial and geometric) representation theory.

Three papers in the volume address such connections: Anderson and Nigro investigate the geometric Satake correspondence in relation to minuscule Schubert Calculus; McGlade, Ram, and Yang wrote a survey on the combinatorics and geometry of integrable representations of quantum affine algebras with a particular focus on level 0; this is motivated by Schubert Calculus on semi-infinite flag manifolds. Finally, Su and Zhong wrote a survey showcasing applications of Maulik and Okounkov's theory of stable envelopes on the cotangent bundle of a flag variety to various problems in geometry and representation theory. This recent direction, which one may call "Cotangent Schubert Calculus", is closely related to the study of characteristic classes of singular varieties; from this viewpoint, it is studied by Fehér, Rimányi, and Weber.

Methods and questions from Schubert Calculus can be adapted to varieties related to flag manifolds or to generalizations of the cohomology ring. A survey by Abe and Horiguchi is investigating the properties of the cohomology rings of Hessenberg varieties; these generalize the usual flag varieties, the Peterson variety, and the Springer fibre. In another direction, Hudson, Matsumura, and Perrin address the problem of defining stable Bott-Samelson classes in the algebraic cobordism; this is closely related to the outstanding problem of defining Schubert classes in more general oriented cohomology theories.

Finally, a paper by Kim, Oh, Ueda, and Yoshida gives an expository account of quasimap theory, and proves a generalization of toric residue mirror symmetry to Grassmannians.

These papers provide a broad overview of current interests in Schubert Calculus and related areas. We would like to thank again all the anonymous referees for their invaluable help, and the Springer editorial staff for the assistance with various technical parts. Finally, we are grateful to Sun Yat-sen University for generously providing funds for this conference and to Springer Nature for providing us the opportunity to publish these papers.

All papers in this volume have been refereed and are in a final form. No version of any of them will be submitted for publication elsewhere.

Guangzhou, China
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# Factorial Flagged Grothendieck Polynomials 

Tomoo Matsumura and Shogo Sugimoto


#### Abstract

The factorial flagged Grothendieck polynomials are defined by flagged set-valued tableaux of Knutson-Miller-Yong [10]. We show that they can be expressed by a Jacobi-Trudi type determinant formula, generalizing the work of Hudson-Matsumura [8]. As an application, we obtain alternative proofs of the tableau and the determinant formulas of vexillary double Grothendieck polynomials, which were originally obtained by Knutson-Miller-Yong [10] and Hudson-Matsumura [8] respectively. Furthermore, we show that each factorial flagged Grothendieck polynomial can be obtained by applying $K$-theoretic divided difference operators to a product of linear polynomials.


Keywords Factorial grothendieck polynomials • Flagged partitions • Flagged set-valued tableaux • Vexillary permutations • Jacobi-Trudi formula • Double grothendieck polynomials

## 1 Introduction

The double Grothendieck polynomials introduced by Lascoux and Schützenberger [11, 12] represent the torus-equivariant $K$-theory classes of the structure sheaves of Schubert varieties in the flag varieties. Their combinatorial formula in terms of pipe dreams or rc graphs was obtained by Fomin-Kirillov [4, 5]. By restricting to Grassmannian elements, or more generally, vexillary permutations, Knutson-MillerYong [10] expressed the associated double Grothendieck polynomials as factorial flagged Grothendieck polynomials defined in terms of flagged set-valued tableaux. This can be regarded as a unification of the work of Wachs [17] and Chen-Li-Louck [3] on flagged tableaux and the work of Buch [2] and McNamara [16] on set-valued

[^0]tableaux. On the other hand, the first author, in the joint work [7] with Hudson, Ikeda and Naruse, obtained a determinant formula of the double Grothendieck polynomials associated to Grassmannian elements (cf. [14]). Such an explicit closed formula was later generalized to the vexillary case in $[1,8]$. Motivated by these results, the first author [15] studied (non-factorial, or single) flagged Grothendieck polynomials and showed their determinant formula in general, beyond the ones given by vexillary permutations.

In this paper, we extend the results in [15] to the factorial (or double) case. Let $x=\left(x_{i}\right)_{i \in \mathbb{Z}_{>0}}$ and $b=\left(b_{i}\right)_{\mathbb{Z}_{>0}}$ be infinite sequences of variables and $\beta$ a formal variable. We denote $u \oplus v=u+v+\beta u v$ for variables $u$ and $v$. Consider a partition $\lambda=\left(\lambda_{1} \geq \cdots \geq \lambda_{r}>0\right)$ with a flagging $f=\left(0<f_{1} \leq \cdots \leq f_{r}\right)$. Let $\mathcal{T}(\lambda, f)$ be the set of flagged set-valued tableaux of the flagged partition $(\lambda, f)$. For each $T \in \mathcal{T}(\lambda, f)$, we denote

$$
[x \mid b]^{T}=\prod_{e \in T}\left(x_{v a l(e)} \oplus b_{\text {val }(e)-r(e)+c(e)}\right),
$$

where the product runs over all entries $e$ of $T, \operatorname{val}(e)$ denotes the numeric value of $e$, and $r(e)$ (resp. $c(e))$ denotes the row (resp. column) index of $e$. Following the work [9, 10] of Knutson-Miller-Yong, we define the factorial flagged Grothendieck polynomial associated to $(\lambda, f)$ by

$$
G_{\lambda, f}(x \mid b):=\sum_{T \in \mathcal{T}(\lambda, f)} \beta^{|T|-|\lambda|}[x \mid b]^{T} .
$$

The following is the main result of this paper.
Main Theorem (Theorem 3.5). For a flagged partition $(\lambda, f)$ of length $r$, we have

$$
G_{\lambda, f}(x \mid b)=\operatorname{det}\left(\sum_{s=0}^{\infty}\binom{i-j}{s} \beta^{s} \mathcal{G}_{\lambda_{i}+j-i+s}^{\left[f_{i}, f_{i}+\lambda_{i}-i\right]}\right)_{1 \leq i, j \leq r}
$$

where $\mathcal{G}_{m}^{[p, q]}=\mathcal{G}_{m}^{[p, q]}(x \mid b), m \in \mathbb{Z}, p, q \in \mathbb{Z}_{\geq 0}$, are polynomials in $x$ and $b$ with coefficients in $\mathbb{Z}[\beta]$, defined by the generating function

$$
\sum_{m \in \mathbb{Z}} \mathcal{G}_{m}^{[p, q]} u^{m}=\frac{1}{1+\beta u^{-1}} \prod_{1 \leq i \leq p} \frac{1+\beta x_{i}}{1-x_{i} u} \prod_{1 \leq i \leq q}\left(1+(u+\beta) b_{i}\right) .
$$

Our proof is a generalization of the ones in $[15,17]$ to the factorial (double) case. We prove that both the tableau and determinant expressions satisfy the same compatibility with divided difference operators, allowing us to show that they coincide by induction. The first author employed a similar argument in order to show a tableau formula for double Grothendieck polynomials associated to 321-avoiding permutations [13]. If we specialize our formula at $\beta=0$, we recover the result by

Chen-Li-Louck [3] for flagged double Schur functions. Our proof is different from theirs, which is based on the lattice-path interpretation of the tableau expression. It is also worth stressing that our proof is completely self-contained. In particular, as a consequence of our main theorem, we obtain a purely algebraic and combinatorial proof of the formula

$$
\mathcal{G}_{m}^{[p, p+m-1]}(x \mid b)=\sum_{T \in \mathcal{T}((m),(p))}[x \mid b]^{T} .
$$

To the best of the authors' knowledge, the only proof of this formula available in the literature uses a geometric argument established in [7].

As an application, we obtain an alternative proof of the tableau and determinant formulas of vexillary double Grothendieck polynomials obtained by Knutson-Miller-Yong [10] and Hudson-Matsumura [8] respectively. We also generalize it to arbitrary factorial flagged Grothendieck polynomials: we show that each of them can be obtained by consecutively applying divided difference operators to a product of linear polynomials. The corresponding results for (non-factorial) flagged Schur and Grothendieck polynomials were obtained in [15, 17] respectively.

## 2 Flagged Grothendieck Polynomials

Let $x=\left(x_{i}\right)_{\mathbb{Z}_{>0}}$ and $b=\left(b_{i}\right)_{\mathbb{Z}_{>0}}$ be sets of infinitely many variables. Let $\mathbb{Z}[\beta]$ be the polynomial ring of a formal variable $\beta$ where we set $\operatorname{deg} \beta=-1$. Let $\mathbb{Z}[\beta][x, b]$ and $\mathbb{Z}[\beta][[x, b]]$ be the rings of polynomials and of formal power series in $x$ and $b$ respectively. Throughout the paper, for each $f \in \mathbb{Z}[\beta][[x, b]]$, let $f^{\star}$ be the element obtained from $f$ by increasing by 1 the index of all the $x_{i}$ 's. We use the generalized binomial coefficients $\binom{n}{i}$ given by $(1+x)^{n}=\sum_{i \geq 0}\binom{n}{i} x^{i}$ for $n \in \mathbb{Z}$ with the convention that $\binom{n}{i}=0$ for all integers $i<0$.

### 2.1 Flagged Partitions and Their Tableaux

A partition $\lambda$ of length $r$ is a weakly decreasing sequence of positive integers $\left(\lambda_{1}, \ldots, \lambda_{r}\right)$. We identify $\lambda$ with its Young diagram $\{(i, j) \mid 1 \leq i \leq r, 1 \leq j \leq r\}$, depicted as a left-aligned array of boxes such that in the $i$ th row from the top there are exactly $\lambda_{i}$ boxes. Let $|\lambda|$ be the total number of boxes in the Young diagram of $\lambda$, i.e. $|\lambda|=\lambda_{1}+\cdots+\lambda_{r}$. A flagging $f$ of a partition of length $r$ is a weakly increasing sequence of positive integers $\left(f_{1}, \ldots, f_{r}\right)$. We call the pair $(\lambda, f)$ a flagged partition.

A set-valued tableau $T$ of shape $\lambda$ is an assignment of a finite subset of positive integers to each box of the Young diagram of $\lambda$, satisfying the requirement that

- weakly increasing in each row: $\max A \leq \min B$ if $A$ fills the box $(i, j)$ and $B$ fills the box $(i, j+1)$ for $1 \leq i \leq r$ and $1 \leq j \leq \lambda_{i}-1$.
- strictly increasing in each column: max $A<\min B$ if $A$ fills the box $(i, j)$ and $B$ fills the box $(i+1, j)$ for $1 \leq i \leq r-1$ and $1 \leq j \leq \lambda_{i+1}$.

An element $e$ of a subset assigned to a box of $T$ is called an entry of $T$ and denoted by $e \in T$. The total number of entries in $T$ is denoted by $|T|$. For each $e \in T$, let $\operatorname{val}(e)$ be its numeric value, $r(e)$ its row index, and $c(e)$ its column index. A flagged set-valued tableau of a flagged partition $(\lambda, f)$ is a set-value tableau of shape $\lambda$ additionally satisfying that the sets assigned to the boxes of the $i$-th row are subsets of $\left\{1, \ldots, f_{i}\right\}$. Let $\mathcal{T}(\lambda, f)$ be the set of all flagged set-valued tableaux of $(\lambda, f)$.

For each tableau $T \in \mathcal{T}(\lambda, f)$, we define

$$
[x \mid b]^{T}:=\prod_{e \in T}\left(x_{\text {val }(e)} \oplus b_{\text {val }(e)-r(e)+c(e)}\right)
$$

where we set $u \oplus v:=u+v+\beta u v$. Following Knutson-Miller-Yong ([9, 10]), we define the factorial flagged Grothendieck polynomial $G_{\lambda, f}(x \mid b)$ as follows.

Definition 2.1 For a flagged partition $(\lambda, f)$, we define

$$
G_{\lambda, f}=G_{\lambda, f}(x \mid b)=\sum_{T \in \mathcal{T}(\lambda, f)} \beta^{|T|-|\lambda|}[x \mid b]^{T} .
$$

If $\lambda$ is an empty partition, we set $G_{\lambda, f}=1$, and if $f_{1}=0$, we set $G_{\lambda, f}=0$.
Example 2.2 Let $\lambda=(3,1)$ and $f=(2,4)$. Then $\mathcal{T}(\lambda, f)$ contains tableaux such as

| 1 | 1 | 12 |  |
| :--- | :--- | :--- | :--- | :--- | :--- |
| 23 |  |  | 12 2 2 <br> 34  $\|$1 12 2 <br> 23   <br> 2 2 2 <br> 4   |

If we change $f$ to $f^{\prime}=(2,3)$, then $\mathcal{T}\left(\lambda, f^{\prime}\right)$ doesn't contain the second and forth tableaux

If $T$ is the first tableau in the above list, $|T|=6$ and we have

$$
[x \mid b]^{T}=\left(x_{1} \oplus b_{1}\right)\left(x_{1} \oplus b_{2}\right)\left(x_{1} \oplus b_{3}\right)\left(x_{2} \oplus b_{4}\right)\left(x_{2} \oplus b_{1}\right)\left(x_{3} \oplus b_{2}\right) .
$$

### 2.2 Divided Difference Operators and $\mathcal{G}_{m}^{[p, q]}$

Definition 2.3 For each positive integer $i$, define the $K$-theoretic divided difference operator $\pi_{i}$ by

$$
\pi_{i}(f)=\frac{\left(1+\beta x_{i+1}\right) f-\left(1+\beta x_{i}\right) s_{i}(f)}{x_{i}-x_{i+1}}
$$

for each $f \in \mathbb{Z}[\beta][[x, b]]$, where $s_{i}$ permutes $x_{i}$ and $x_{i+1}$.

The following properties of $\pi_{i}$ are easy to check (cf. [15, Sect. 2.1]).
Lemma 2.4 For each positive integer $i$ and $f, g \in \mathbb{Z}[\beta][[x, b]]$, we have
(1) $\pi_{i}(f g)=\pi_{i}(f) g+s_{i}(f) \pi_{i}(g)+\beta s_{i}(f) g$.
(2) If $f$ is symmetric in $x_{i}$ and $x_{i+1}$, then $\pi_{i}(f)=-\beta f$ and $\pi_{i}(f g)=f \pi_{i}(g)$.
(3) $\pi_{i}(f)=-\beta f$, then $f$ is symmetric in $x_{i}$ and $x_{i+1}$.

Definition 2.5 Define $\mathcal{G}_{m}^{[p, q]}=\mathcal{G}_{m}^{[p, q]}(x \mid b), m \in \mathbb{Z}, p, q \in \mathbb{Z}_{\geq 0}$, by the generating function

$$
\mathcal{G}_{u}^{[p, q]}=\sum_{m \in \mathbb{Z}} \mathcal{G}_{m}^{[p, q]} u^{m}=\frac{1}{1+\beta u^{-1}} \prod_{1 \leq i \leq p} \frac{1+\beta x_{i}}{1-x_{i} u} \prod_{1 \leq i \leq q}\left(1+(u+\beta) b_{i}\right) .
$$

If $q=0$, then we denote $\mathcal{G}_{m}^{[p]}=\mathcal{G}_{m}^{[p, q]}$.
Remark 2.6 (1) Since the degree of $u$ in $\mathcal{G}_{u}^{[0, q]}$ is at most $q$, we see that $\mathcal{G}_{m}^{[0, q]}=0$ for $m>q$.
(2) We can show that $\mathcal{G}_{m}^{[p, q]}=(-\beta)^{-m}$ for $m \leq 0$ by a direct computation. In fact, by using $1+(u+\beta) b_{i}=\frac{1-u \bar{b}_{i}}{1+\beta \bar{b}_{i}}$ where $\bar{b}_{i}=\frac{1-u b_{i}}{1+\beta b_{i}}$, we have

$$
\begin{aligned}
\mathcal{G}_{u}^{[p, q]} & =\frac{1}{1+\beta u^{-1}} \prod_{1 \leq i \leq p} \frac{1+\beta x_{i}}{1-x_{i} u} \prod_{1 \leq i \leq q} \frac{1-u \bar{b}_{i}}{1+\beta \bar{b}_{i}} \\
& =\prod_{1 \leq i \leq p}\left(1+\beta x_{i}\right) \prod_{1 \leq i \leq q} \frac{1}{1+\beta \bar{b}_{i}} \sum_{m \in \mathbb{Z}}\left(\sum_{l \geq 0} \sum_{n \geq 0} h_{m+n-l}(x) e_{l}(-\bar{b})(-\beta)^{n}\right) u^{m} .
\end{aligned}
$$

Suppose $m \leq 0$. Then the coefficient of $u^{m}$ in the summation is

$$
\prod_{1 \leq i \leq p} \frac{1}{1+\beta x_{i}} \prod_{1 \leq i \leq q}\left(1+\beta \bar{b}_{i}\right)(-\beta)^{-m}
$$

Thus we have $\mathcal{G}_{m}^{[p, q]}=(-\beta)^{-m}$ for $m \leq 0$.
(3) Similarly to (2), we can also check that $\mathcal{G}_{m}^{[1]}=x_{1}^{m}$ for $m \geq 0$.
(4) We have $\mathcal{G}_{q+r}^{[1, q]}=x_{1}^{r} \mathcal{G}_{q}^{[1, q]}$ for any integer $r \geq 0$. Indeed, if we let

$$
\prod_{i=1}^{q}\left(1+(u+\beta) b_{i}\right)=\sum_{i=0}^{q} E_{i} u^{i},
$$

then, by (3), we can compute

$$
\mathcal{G}_{q+r}^{[1, q]}=\sum_{i=0}^{d} \mathcal{G}_{i+r}^{[1]} E_{q-i}=x_{1}^{r} \sum_{i=0}^{q} \mathcal{G}_{i}^{[1]} E_{q-i}=x_{1}^{r} \mathcal{G}_{q}^{[1, q]}
$$

The following basic lemmas will be used in the next section.
Lemma 2.7 For each $m \in \mathbb{Z}, p, q \in \mathbb{Z}_{\geq 0}$, we have $\pi_{i}\left(\mathcal{G}_{m}^{[p, q]}\right)= \begin{cases}\mathcal{G}_{m-1}^{[p+1, q]} & (i=p), \\ -\beta \mathcal{G}_{m}^{[p, q]} & (i \neq p) .\end{cases}$
Proof The claim follows from applying divided difference operators to the generating function of $\mathcal{G}_{m}^{[p, q]}$ (cf. [15, Lemma 1]).

Lemma 2.8 Lett and $_{i}(i=1, \ldots, n)$ be arbitrarypositive integers and $t^{\prime}:=t+1$. We have

$$
\pi_{t}\left(\prod_{i=1}^{n}\left(x_{t} \oplus b_{t_{i}}\right)\right)=\sum_{v=0}^{n-1}\left(\prod_{i=1}^{v}\left(x_{t} \oplus b_{t_{i}}\right) \prod_{i=v+2}^{n}\left(x_{t^{\prime}} \oplus b_{t_{i}}\right)\right)+\beta \sum_{v=1}^{n-1}\left(\prod_{i=1}^{v}\left(x_{t} \oplus b_{t_{i}}\right) \prod_{i=v+1}^{n}\left(x_{t^{\prime}} \oplus b_{t_{i}}\right)\right) .
$$

In particular, $\pi_{t}\left(x_{t} \oplus b_{t_{1}}\right)=1$.
Proof We prove the claim by induction on $n$. When $n=1$, we can show that $\pi_{t}\left(x_{t} \oplus\right.$ $\left.b_{s}\right)=1$ by a direct computation. For a general $n$, we apply Lemma 2.4 (1) with $f=x_{t} \oplus b_{t_{n}}$ and $g$ the rest:

$$
\pi_{t}\left(\prod_{i=1}^{n}\left(x_{t} \oplus b_{t_{i}}\right)\right)=\prod_{i=1}^{n-1}\left(x_{t} \oplus b_{t_{i}}\right)+\left(x_{t^{\prime}} \oplus b_{t_{n}}\right) \pi_{t}\left(\prod_{i=1}^{n-1}\left(x_{t} \oplus b_{t_{i}}\right)\right)+\beta\left(x_{t^{\prime}} \oplus b_{t_{n}}\right) \prod_{i=1}^{n-1}\left(x_{t} \oplus b_{t_{i}}\right) .
$$

By the induction hypothesis, we have

$$
\begin{aligned}
\pi_{t}\left(\prod_{i=1}^{n}\left(x_{t} \oplus b_{t_{i}}\right)\right) & =\prod_{i=1}^{n-1}\left(x_{t} \oplus b_{t_{i}}\right)+\left(x_{t^{\prime}} \oplus b_{t_{n}}\right) \sum_{v=0}^{n-2}\left(\prod_{i=1}^{v}\left(x_{t} \oplus b_{t_{i}}\right) \prod_{i=v+2}^{n-1}\left(x_{t^{\prime}} \oplus b_{t_{i}}\right)\right) \\
& +\beta\left(\left(x_{t^{\prime}} \oplus b_{t_{n}}\right) \sum_{v=1}^{n-2}\left(\prod_{i=1}^{v}\left(x_{t} \oplus b_{t_{i}}\right) \prod_{i=v+1}^{n-1}\left(x_{t^{\prime}} \oplus b_{t_{i}}\right)\right)+\left(x_{t^{\prime}} \oplus b_{t_{n}}\right) \prod_{i=1}^{n-1}\left(x_{t} \oplus b_{t_{i}}\right)\right) .
\end{aligned}
$$

This is exactly the right hand side of the desired formula.
Remark 2.9 Since the left hand side of the formula in Lemma 2.8 is symmetric in the variables $b_{t_{1}}, \ldots, b_{t_{n}}$, we can conclude that the right hand side is symmetric in $x_{t}$ and $x_{t+1}$.
Lemma 2.10 ([cf. Lemma 3 [15]]) For any integers $m \in \mathbb{Z}, p \in \mathbb{Z}_{\geq 1}$ and $q \in \mathbb{Z}_{\geq 0}$, we have

$$
\frac{1}{1+x_{1} \beta}\left(\mathcal{G}_{m}^{[p, q]}-x_{1} \mathcal{G}_{m-1}^{[p, q]}\right)=\left(\mathcal{G}_{m}^{[p-1, q]}\right)^{\star}
$$

Proof It follows from the identity

$$
\sum_{m \in \mathbb{Z}} \frac{1}{1+x_{1} \beta}\left(\mathcal{G}_{m}^{[p, q]}-x_{1} \mathcal{G}_{m-1}^{[p, q]}\right) u^{m}=\left(\sum_{m \in \mathbb{Z}} \mathcal{G}_{m}^{[p-1, q]} u^{m}\right)^{\star}
$$

which can be checked by a direct computation.

## 3 Propositions and the Main Theorem

In this section, we prove the main theorem. First we show four propositions that will be used in the proof of the main theorem given at the end of this section. Throughout the section, we let $(\lambda, f)$ to be a flagged partition of length $r$.
$\underset{\widetilde{G}}{\text { Definition 3.1 }}$ We denote the following determinantal expression by $\widetilde{G}_{\lambda, f}=$ $\widetilde{G}_{\lambda, f}(x \mid b)$ :

$$
\widetilde{G}_{\lambda, f}(x \mid b)=\operatorname{det}\left(\sum_{s=0}^{\infty}\binom{i-j}{s} \beta^{s} \mathcal{G}_{\lambda_{i}-i+j+s}^{\left[f_{i}, f_{i}+\lambda_{i}-i\right]}\right)_{(1 \leq i, j \leq r)}
$$

If $\lambda=\varnothing$, we set $\widetilde{G}_{\lambda, f}=1$. If $\lambda \neq \varnothing$ and $f_{1}=0$, the first row of the determinant is identically zero by Remark 2.6 (1) so that we set $\widetilde{G}_{\lambda, f}=0$.
Proposition 3.2 For any integer $q \geq 0$, we have $\widetilde{G}_{(q),(1)}(x \mid b)=G_{(q),(1)}(x \mid b)$, or equivalently, $\mathcal{G}_{q}^{[1, q]}=\prod_{i=1}^{q}\left(x_{1} \oplus b_{i}\right)$.

Proof We prove the claim by induction on $q$. When $q=0$, it is trivial. Suppose that $q>0$. By definition of $\mathcal{G}_{m}^{[p, q]}$, we have

$$
\sum_{m \in \mathbb{Z}} \mathcal{G}_{m}^{[1, q]} u^{m}=\left(\sum_{m \in \mathbb{Z}} \mathcal{G}_{m}^{[1, q-1]} u^{m}\right)\left(1+\beta b_{q}+b_{q} u\right)
$$

so that

$$
\mathcal{G}_{q}^{[1, q]}=\mathcal{G}_{q}^{[1, q-1]}\left(1+\beta b_{q}\right)+\mathcal{G}_{q-1}^{[1, q-1]} b_{q} .
$$

Since $\mathcal{G}_{q}^{[1, q-1]}=x_{1} \mathcal{G}_{q-1}^{[1, q-1]}$ by Remark 2.6 (4), we obtain

$$
\mathcal{G}_{q}^{[1, q]}=\mathcal{G}_{q-1}^{[1, q-1]} x_{1}\left(1+\beta b_{q}\right)+\mathcal{G}_{q-1}^{[1, q-1]} b_{q}=\mathcal{G}_{q-1}^{[1, q-1]}\left(x_{1} \oplus b_{q}\right)
$$

The induction hypothesis implies the desired formula.
Proposition 3.3 If $f_{1}=1$, then we have
(1) $\widetilde{G}_{\lambda, f}(x \mid b)=\mathcal{G}_{\lambda_{1}}^{\left[1, \lambda_{1}\right]} \cdot \widetilde{G}_{\lambda^{\prime}, f^{\prime}}(x \mid b)^{\star}$,
(2) $G_{\lambda, f}(x \mid b)=\mathcal{G}_{\lambda_{1}}^{\left[1, \lambda_{1}\right]} \cdot G_{\lambda^{\prime}, f^{\prime}}(x \mid b)^{\star}$,
where $\lambda^{\prime}=\left(\lambda_{2}, \lambda_{3}, \ldots, \lambda_{r}\right)$ and $f^{\prime}=\left(f_{2}-1, f_{3}-1, \ldots, f_{r}-1\right)$.
Proof (1) We show that the left hand side coincides with the right hand side by the column operation of subtracting the $(j-1)$-st column multiplied with $x_{1}(1+$ $\left.x_{1} \beta\right)^{-1}$ from the $j$-th column for $j=2, \ldots, r$. By Remark 2.6 (4), the ( $1, j$ ) entry of $\widetilde{G}_{\lambda, f}(x \mid b)$ is

$$
\sum_{s=0}^{\infty}\binom{1-j}{s} \beta^{s} \mathcal{G}_{\lambda_{1}-1+j+s}^{\left[1, \lambda_{1}\right]}=\sum_{s=0}^{\infty}\binom{1-j}{s} \beta^{s} x_{1}^{j-1+s} \mathcal{G}_{\lambda_{1}}^{\left[1, \lambda_{1}\right]}=\left(1+x_{1} \beta\right)^{1-j} x_{1}^{j-1} \mathcal{G}_{\lambda_{1}}^{\left[1, \lambda_{1}\right]} .
$$

Thus after the above column operation, the first row of $\widetilde{G}_{\lambda, f}(x \mid b)$ becomes $\left(\mathcal{G}_{\lambda_{1}}^{\left[1, \lambda_{1}\right]}, 0, \ldots, 0\right)$. We compute the $(i, j)$ entry of the resulting determinant for $i, j \geq 2$ :

$$
\begin{aligned}
& \sum_{s=0}^{\infty}\binom{i-j}{s} \beta^{s} \mathcal{G}_{\lambda_{1}-i+j+s}^{\left[f_{i}, f_{i}+\lambda_{i}-i\right]}-\frac{x_{1}}{1+x_{1} \beta} \sum_{s=0}^{\infty}\binom{i-j+1}{s} \beta^{s} \mathcal{G}_{\lambda_{i}-i+j-1+s}^{\left[f_{i}, f_{i}+\lambda_{i}-i\right]} \\
= & \sum_{s=0}^{\infty}\binom{i-j}{s} \beta^{s} \mathcal{G}_{\lambda_{1}-i+j+s}^{\left[f_{i}, f_{i}+\lambda_{i}-i\right]}-\frac{x_{1}}{1+x_{1} \beta} \sum_{s=0}^{\infty}\left(\binom{i-j}{s}+\binom{i-j}{s-1}\right) \beta^{s} \mathcal{G}_{\lambda_{i}-i+j-1+s}^{\left[f_{i}, f_{i}+\lambda_{i}-i\right]} \\
= & \sum_{s=0}^{\infty}\binom{i-j}{s} \beta^{s}\left(\mathcal{G}_{\lambda_{1}-i+j+s}^{\left[f_{i}, f_{i}+\lambda_{i}-i\right]}-\frac{x_{1}}{1+x_{1} \beta} \mathcal{G}_{\lambda_{1}-i+j+s-1}^{\left[f_{i}, f_{i}+\lambda_{i}-i\right]}\right)-\frac{x_{1} \beta}{1+x_{1} \beta} \sum_{s^{\prime}=-1}^{\infty}\binom{i-j}{s^{\prime}} \beta^{s^{\prime}} \mathcal{G}_{\lambda_{i}-i+j+s^{\prime}}^{\left[f_{i}, f_{i}+\lambda_{i}-i\right]} \\
= & \sum_{s=0}^{\infty}\binom{i-j}{s} \beta^{s} \frac{1}{1+x_{1} \beta}\left(\mathcal{G}_{\lambda_{i}-i+j+s}^{\left[f_{i}, f_{i}+\lambda_{i}-i\right]}-x_{1} \mathcal{G}_{\lambda_{i}-i+j+s-1}^{\left[f_{i}, f_{i}+\lambda_{i}-i\right]}\right) \\
= & \sum_{s=0}^{\infty}\binom{i-j}{s} \beta^{s}\left(\mathcal{G}_{\lambda_{i}-1+j+s}^{\left[f_{i}-1+f_{i}+\lambda_{i}-i\right]}\right)^{\star} .
\end{aligned}
$$

Here we have used a well-known identity for binomial coefficients for the first equality, the fact that $\binom{i-j}{s^{\prime}}=0$ for $s^{\prime}=-1$ for the third equality, and Lemma 2.10 for the last equality. Finally the desired formula follows from the cofactor expansion with respect to the first row for the resulting determinant after the column operation.
(2) For any $T \in \mathcal{T}(\lambda, f)$, the entries on the first row of $T$ are all 1 and all other entries are greater than 1 . There is a bijection from $\mathcal{T}(\lambda, f)$ to $\mathcal{T}\left(\lambda^{\prime}, f^{\prime}\right)$ sending $T$ to $T^{\prime}$ obtained from $T$ by deleting its first row and decreasing the numeric values of the rest of the entries by 1 . Under this bijection, we have

$$
[x \mid b]^{T}=\left([x \mid b]^{T^{\prime}}\right)^{\star} \cdot \prod_{j=1}^{\lambda_{1}}\left(x_{1} \oplus b_{j}\right) .
$$

Now the claim follows from Proposition 3.2.
Proposition 3.4 If $\lambda_{1}>\lambda_{2}$ and $f_{1}<f_{2}$, then we have
(1) $\pi_{f_{1}} \widetilde{G}_{\lambda, f}(x \mid b)=\widetilde{G}_{\lambda^{\prime}, f^{\prime}}(x \mid b)$,
(2) $\pi_{f_{1}} G_{\lambda, f}(x \mid b)=G_{\lambda^{\prime}, f^{\prime}}(x \mid b)$,
where $\lambda^{\prime}=\left(\lambda_{1}-1, \lambda_{2}, \ldots, \lambda_{r}\right)$ and $f^{\prime}=\left(f_{1}+1, f_{2}, \ldots, f_{r}\right)$.
Proof (1) First observe that the entries of the determinant are symmetric in $x_{f_{1}}$ and $x_{f_{1}+1}$ except the ones on the first row, since $f_{1}<f_{2}$. We consider the cofactor expansion of $\widetilde{G}_{\lambda, f}(x \mid b)$ with respect to the first row: let $\Delta_{i, j}$ be the cofactor of the $(i, j)$ entry and we have

$$
\widetilde{G}_{\lambda, f}=\sum_{j=1}^{r}(-1)^{1+j} \Delta_{1, j} \sum_{s=0}^{\infty}\binom{1-j}{s} \beta^{s} \mathcal{G}_{\lambda_{1}-1+j+s}^{\left[f_{1}, f_{1}+\lambda_{1}-1\right]}
$$

Applying $\pi_{f_{1}}$ to this expansion, then we obtain

$$
\begin{aligned}
\pi_{f_{1}} \widetilde{G}_{\lambda, f} & =\sum_{j=1}^{r}(-1)^{1+j} \Delta_{1, j} \sum_{s=0}^{\infty}\binom{1-j}{s} \beta^{s} \pi_{f_{1}}\left(\mathcal{G}_{\lambda_{1}-1+j+s}^{\left[f_{1}, f_{1}+\lambda_{1}-1\right]}\right) \\
& =\sum_{j=1}^{r}(-1)^{1+j} \Delta_{1, j} \sum_{s=0}^{\infty}\binom{1-j}{s} \beta^{s} \mathcal{G}_{\lambda_{1}-2+j+s}^{\left[f_{1}+1, f_{1}+\lambda_{1}-1\right]},
\end{aligned}
$$

in view of Lemma 2.4 (2) and Lemma 2.7. The last expression is the cofactor expansion of $\widetilde{G}_{\lambda^{\prime}, f^{\prime}}$, and thus we obtain the desired formula.
(2) Let $t:=f_{1}$ and $t^{\prime}:=f_{1}+1$. First we define an equivalence relation in $\mathcal{T}(\lambda, f)$ as follows. Two tableaux $T_{1}, T_{2} \in \mathcal{T}(\lambda, f)$ are equivalent if the next two conditions are satisfied:
(i) the boxes containing either $t$ or $t^{\prime}$ in $T_{1}$ coincide with those in $T_{2}$;
(ii) if each box $\left(1, \lambda_{1}\right)$ in $T_{1}$ and $T_{2}$ contains $t$, then both of them contain only $t$ or both of them contain $t$ along with other entries.

Let $\mathscr{A}$ be an equivalence class for $\mathcal{T}(\lambda, f)$, then the configuration of $t$ and $t^{\prime}$ for the tableaux in $\mathscr{A}$ can be depicted as in Fig. 1. The one row rectangle $A_{1}$ on the first row consists of $m_{1}$ boxes with entries $t$. Each one-row rectangle $A_{i}(2 \leq i \leq k)$ with $*$ consists of $m_{i}$ boxes and each box contains $t$ or $t^{\prime}$ or both so that the total number of entries $t$ and $t^{\prime}$ in $A_{i}$ is $m_{i}$ or $m_{i}+1$. Each two-row rectangle $B_{j}(1 \leq j \leq k)$ consists of $r_{i}$ columns with $t$ on the first row and $t^{\prime}$ on the second. Note that $m_{i}$ and $r_{i}$ may be 0 so that the rectangles in Fig. 1 may be not connected.


Fig. 1 Configuration of $t$ and $t^{\prime}$ for $\mathscr{A}$

Let us write

$$
G_{\mathscr{A}}:=\sum_{T \in \mathscr{A}} \beta^{|T|-|\lambda|}[x \mid b]^{T}=R\left(A_{1}\right) R^{\prime}(\mathscr{A}), R^{\prime}(\mathscr{A}):=R(\mathscr{A})\left(\prod_{i=2}^{k} R\left(A_{i}\right)\right)\left(\prod_{j=1}^{k} R\left(B_{j}\right)\right),
$$

where $R(\mathscr{A})$ is the polynomial contributed from the entries other than $t$ and $t^{\prime}$ and $R\left(A_{i}\right)$ and $R\left(B_{j}\right)$ are the polynomials contributed from the entries $t$ and $t^{\prime}$ in $A_{i}$ and $B_{j}$ respectively. It is obvious that $R(\mathscr{A})$ and $R\left(B_{j}\right)(1 \leq j \leq k)$ are symmetric in $x_{t}$ and $x_{t^{\prime}}$. Moreover, in view of Remark 2.9, $R\left(A_{i}\right)(2 \leq i \leq k)$ are also symmetric in $x_{t}$ and $x_{t^{\prime}}$. Thus $R^{\prime}(\mathscr{A})$ is symmetric in $x_{t}$ and $x_{t^{\prime}}$.

Next we decompose $\mathcal{T}(\lambda, f) / \sim$ into the subsets $\mathcal{H}_{1}, \mathcal{H}_{2}, \mathcal{H}_{3}$ consisting of equivalence classes with configurations respectively satisfying conditions (1) $m_{1}=0$, (2) $m_{1}=1$ and the box $\left(1, \lambda_{1}\right)$ contains $t$ along with other entries, (3) $m_{1} \geq 1$ and the box $\left(1, \lambda_{1}\right)$ contains only $t$. There is a bijection from $\mathcal{H}_{1}$ to $\mathcal{H}_{2}$ sending $\mathscr{A}$ to $\mathscr{A}^{\prime}$ whose configuration is obtained by inserting $t$ in the box $\left(1, \lambda_{1}\right)$ of $\mathscr{A}$. This also defines a bijection from $\mathscr{A}$ to $\mathscr{A}^{\prime}$ (say, it maps $T$ to $T^{\prime}$ ). Under this bijection, we have $G_{\mathscr{A}^{\prime}}=\beta\left(x_{t} \oplus b_{t-1+\lambda_{1}}\right) G_{\mathscr{A}}$. Since $G_{\mathscr{A}}$ is symmetric in $x_{t}$ and $x_{t^{\prime}}$, by Lemma 2.4 (2) and Lemma 2.8, we have $\pi_{t}\left(G_{\mathscr{A}^{\prime}}\right)=\beta G_{\mathscr{A}}$. Thus

$$
\pi_{t}\left(\sum_{\mathscr{A} \in \mathcal{H}_{1} \sqcup \mathcal{H}_{2}} G_{\mathscr{A}}\right)=\sum_{\mathscr{A} \in \mathcal{H}_{1}}\left(\pi_{t}\left(G_{\mathscr{A}}\right)+\pi_{t}\left(G_{\mathscr{A}^{\prime}}\right)\right)=\sum_{\mathscr{A} \in \mathcal{H}_{1}}\left(-\beta G_{\mathscr{A}}+\beta G_{\mathscr{A}}\right)=0 .
$$

As a consequence, we have

$$
\begin{equation*}
\pi_{t}\left(G_{\lambda, f}\right)=\pi_{t}\left(\sum_{\mathscr{A} \in \mathcal{H}_{3}} G_{\mathscr{A}}\right) . \tag{1}
\end{equation*}
$$

Now it remains to show that the right hand side of (1) coincides with $G_{\lambda^{\prime}, f^{\prime}}$. Define an equivalence relation $\sim$ in $\mathcal{T}\left(\lambda^{\prime}, f^{\prime}\right)$ by the condition (i), that is, $T_{1}^{\prime}$ and $T_{2}^{\prime}$ in $\mathcal{T}\left(\lambda^{\prime}, f^{\prime}\right)$ are equivalent if the boxes containing either $t$ or $t^{\prime}$ in $T_{1}^{\prime}$ coincide with those in $T_{2}^{\prime}$. For an arbitrary equivalence class $\mathscr{A}^{\prime}$ for $\mathcal{T}\left(\lambda^{\prime}, f^{\prime}\right)$, we can describe its configuration of $t$ and $t^{\prime}$ as in Fig. 2. Similarly to Fig. 1, $A_{i}(2 \leq i \leq k)$ is a rectangle consisting of $m_{i}$ boxes with entries $t, t^{\prime}$ or both of them, $B_{j}(1 \leq j \leq k)$ is a two-row rectangle with $r_{j}$ columns with $t$ on the first row and $t^{\prime}$ on the second. The right-most rectangle $A_{1}^{\prime}$ has $m_{1}^{\prime}$ boxes with entries $t, t^{\prime}$ or both of them.

There is a bijection from $\mathcal{H}_{3}$ to $\mathcal{T}\left(\lambda^{\prime}, f^{\prime}\right) / \sim$ sending an equivalence class $\mathscr{A}$ to $\mathscr{A}^{\prime}$ whose configuration of $t$ and $t^{\prime}$ is obtained from the one for $\mathscr{A}$ by erasing the box $\left(1, \lambda_{1}\right)$ and replacing $t$ with $*$ in the rectangle $A_{1}$. Under this bijection, $\pi_{t}\left(G_{\mathscr{A}}\right)=G_{\mathscr{A}}$. Indeed, let $\mathscr{A} \in \mathcal{H}_{3}$ with the configuration as depicted in Fig. 1 and $\mathscr{A}^{\prime} \in \mathcal{T}\left(\lambda^{\prime}, f^{\prime}\right) / \sim$ as in Fig. 2 where $m_{1}^{\prime}=m_{1}-1 \geq 0$. We have $R\left(A_{1}\right)=$ $\prod_{i=1}^{m_{1}}\left(x_{t} \oplus b_{t-1+\lambda_{1}-m_{1}+i}\right)$. By Lemma 2.4 (2) and Lemma 2.8, we have


Fig. 2 Configuration of $t$ and $t^{\prime}$ for $\mathscr{A}^{\prime}$

$$
\begin{aligned}
\pi_{t}\left(G_{\mathscr{A}}\right)= & \left\{\sum_{v=0}^{m_{1}-1}\left(\prod_{i=1}^{v}\left(x_{t} \oplus b_{t-1+\lambda_{1}-m_{1}+i}\right) \prod_{i=v+2}^{m_{1}}\left(x_{t^{\prime}} \oplus b_{t-1+\lambda_{1}-m_{1}+i}\right)\right)\right. \\
& \left.+\beta \sum_{v=1}^{m_{1}-1}\left(\prod_{i=1}^{v}\left(x_{t} \oplus b_{t-1+\lambda_{1}-m_{1}+i}\right) \prod_{i=v+1}^{m_{1}}\left(x_{t^{\prime}} \oplus b_{t-1+\lambda_{1}-m_{1}+i}\right)\right)\right\} R^{\prime}(\mathscr{A})
\end{aligned}
$$

which is exactly $G_{\mathscr{A}^{\prime}}$. Thus we have

$$
\pi_{t} G_{\lambda, f}=\sum_{\mathscr{A} \in \mathcal{H}_{3}} \pi_{t}\left(G_{\mathscr{A}}\right)=\sum_{\mathscr{A}^{\prime} \in \mathcal{T}\left(\lambda^{\prime}, f^{\prime}\right) / \sim} G_{\mathscr{A}^{\prime}}=G_{\lambda^{\prime}, f^{\prime}}
$$

This completes the proof.
Theorem 3.5 For a flagged partition $(\lambda, f)$, we have $G_{\lambda, f}(x \mid b)=\widetilde{G}_{\lambda, f}(x \mid b)$.
Proof With the help of Proposition 3.2, 3.3, and 3.4, the proof is by induction on $|f|=f_{1}+\cdots+f_{r}$, parallel to the one in [15]. If $|f|=1$, then $\lambda=\left(\lambda_{1}\right)$ and $f=(1)$. Thus it follows from Proposition 3.2. Suppose that $|f|>1$. We prove in two cases: $f_{1}=1$ or $f_{1}>1$. If $f_{1}=1$, then we can apply Proposition 3.3 to both $\widetilde{G}_{\lambda, f}$ and $G_{\lambda, f}$. The right hand sides of the resulting formulas coincide by the induction hypothesis, and thus the claim holds. If $f_{1}>1$, then Proposition 3.4 implies that $\pi_{g_{1}} \widetilde{G}_{\mu, g}=\widetilde{G}_{\lambda, f}$ and $\pi_{g_{1}} G_{\mu, g}=G_{\lambda, f}$ where $\mu=\left(\lambda_{1}+1, \lambda_{2}, \ldots, \lambda_{r}\right)$ and $g=\left(f_{1}-1, f_{2}, \ldots, f_{r}\right)$. The left hand sides of these equalities coincide by the induction hypothesis, and thus the claim holds.

## 4 Vexillary Double Grothendieck Polynomials

In this section, we prove that the double Grothendieck polynomials associated to vexillary permutations are in fact factorial Grothendieck polynomials (Theorem 4.2), giving an alternative proof to the corresponding results in [8, 10]. Moreover we show
that any factorial Grothendieck polynomial can be obtained from a product of certain linear polynomials by applying a sequence of divided difference operators (Theorem 4.3).

The double Grothendieck polynomials were introduced by Lascoux and Schützenberger [11, 12]. For any permutation $w \in S_{n}$, we define the associated double Grothendieck polynomial $\mathfrak{G}_{w}=\mathfrak{G}_{w}(x \mid b)$ as follows. Let $w_{0}$ be the longest element of the symmetric group $S_{n}$. We set

$$
\mathfrak{G}_{w_{0}}=\prod_{i+j \leq n}\left(x_{i} \oplus b_{j}\right)
$$

For an element $w \in S_{n}$ such that $\ell(w)<\ell\left(w_{0}\right)$, we can choose a simple reflection $s_{i} \in S_{n}$ such that $\ell\left(w s_{i}\right)=\ell(w)+1$. Here $\ell(w)$ is the length of $w$. We then define

$$
\mathfrak{G}_{w}=\pi_{i}\left(\mathfrak{G}_{w s_{i}}\right)
$$

The polynomial $\mathfrak{G}_{w}$ is defined independently from the choice of such $s_{i}$, since the divided difference operators satisfy the Coxeter relations. From this point of view, we can write $\mathfrak{G}_{w}=\pi_{v} \mathfrak{G}_{w_{0}}$ with $v=w_{0} w$ and $\pi_{v}=\pi_{i_{k}} \cdots \pi_{i_{1}}$ where $v=s_{i_{1}} \cdots s_{i_{k}}$ is a reduced expression.

A permutation $w \in S_{n}$ is called vexillary if it is 2143-avoiding, i.e. there is no $a<b<c<d$ such that $w(b)<w(a)<w(d)<w(c)$. We briefly recall how to obtain a flagged partition from a vexillary permutation. We follow [6, 10] (cf. [8]). For each $w \in S_{n}$, let $r_{w}$ be the rank function of $w \in S_{n}$ defined by $r_{w}(p, q):=\sharp\{i \leq$ $p \mid w(i) \leq q\}$ for $1 \leq p, q \leq n$. The diagram $D(w)$ of $w$ is defined as

$$
D(w):=\left\{(p, q) \in\{1, \ldots, n\} \times\{1, \ldots, n\} \mid w(p)>q, \text { and } w^{-1}(q)>p\right\}
$$

The essential set $\operatorname{Ess}(w)$ is the subset of $D(w)$ given by

$$
\operatorname{Ess}(w):=\{(p, q) \mid(p+1, q),(p, q+1) \notin D(w)\}
$$

If $w$ is vexillary, we can choose a flagging set of $w$ (cf. [8]), which is a subset $\left\{\left(p_{i}, q_{i}\right), i=1, \ldots, r\right\}$ of $\{1, \ldots, n\} \times\{1, \ldots, n\}$ containing $\operatorname{Ess}(w)$ and satisfying

$$
\begin{align*}
& p_{1} \leq p_{2} \leq \cdots \leq p_{r}, \quad q_{1} \geq q_{2} \geq \cdots \geq q_{r}  \tag{2}\\
& p_{i}-r_{w}\left(p_{i}, q_{i}\right)=i, \quad \forall i=1, \ldots, r \tag{3}
\end{align*}
$$

An associated flagged partition $(\lambda(w), f(w))$ of length $r$ is given by setting $f_{i}(w):=$ $p_{i}$ and $\lambda(w)_{i}=q_{i}-p_{i}+i$ for $i=1, \ldots, r$. We remark that the set $\mathcal{T}(\lambda(w), f(w))$ depends only on $w$ but not on the choice of a flagging set.

Example 4.1 A permutation $w \in S_{n}$ is dominant if $D(w)$ is the Young diagram of a partition and the values of $r_{w}$ on $D(w)$ are zero. Such permutation is vexillary, and
in this case, $\lambda(w)$ is the partition whose Young diagram is $D(w)$ and its flagging is $f(w)=(1,2, \ldots, r)$ where $r$ is the length of $\lambda(w)$.

The following theorem was obtained in [8, 10]. We give an alternative proof.
Theorem 4.2 If $w \in S_{n}$ is vexillary, then we have $\mathfrak{G}_{w}=G_{\lambda(w), f(w)}$.
Proof We closely follow the proof of Theorem 2 in [15], which is by induction on $\left(n, \ell\left(w_{0}\right)-\ell(w)\right)$ with the lexicographic order. For the longest element $w_{0} \in S_{n}$, by definition, we have

$$
G_{\lambda\left(w_{0}\right), f\left(w_{0}\right)}=\prod_{i+j \leq n}\left(x_{i} \oplus b_{j}\right)=\mathfrak{G}_{w_{0}}
$$

Consider $w \in S_{n}$ with $w \neq w_{0}$, id. Let $d$ be the leftmost descent of $w$, i.e. $d$ is the smallest number such that $w(d)>w(d+1)$. If $d>1$ (Case 1), by Proposition 3.4, the proof is identical to the one in [15]. Suppose that $d=1$. If $m:=w(1)<n$ (Case 2), let $w^{\prime}:=s_{m} w$ and consider the dominant permutation

$$
u=(m+1, m, m+2, m+3, \ldots, n, m-1, m-2, \ldots, 1)
$$

as in [15]. The proof is identical to the one in [15], except that we have

$$
\pi_{1} \mathfrak{G}_{u}=\left(x_{1} \oplus b_{m}\right)^{-1} \mathfrak{G}_{u}, \quad G_{\lambda(w), f(w)}=\left(x_{1} \oplus b_{m}\right)^{-1} G_{\lambda\left(w^{\prime}\right), f\left(w^{\prime}\right)}
$$

The former identity holds since we have $\mathfrak{G}_{u}=\left(x_{1} \oplus b_{m}\right) R_{u}$ where $R_{u}$ is symmetric in $x_{1}$ and $x_{2}$. The latter follows from the fact that there is a bijection from $\mathcal{T}(\lambda(w), f(w))$ to $\mathcal{T}\left(\lambda\left(w^{\prime}\right), f\left(w^{\prime}\right)\right)$ sending $T$ to $T^{\prime}$ which is obtained from $T$ by adding a box in the first row with entry 1. Finally, if $w(1)=n$ (Case 3), we find a reduced expression $s_{i_{1}} \cdots s_{i_{k}}$ of $w_{0} w$ where $i_{1}, \ldots, i_{k} \geq 2$. Let $w^{\prime}=(w(2), \ldots, w(n)) \in S_{n-1}$. Then we have

$$
\mathfrak{G}_{w}=\prod_{i=1}^{n-1}\left(x_{1} \oplus b_{i}\right) \cdot \pi_{i_{k}} \pi_{i_{k-1}} \cdots \pi_{i_{1}}\left(\prod_{i+j \leq n-1}\left(x_{i+1} \oplus b_{j}\right)\right)=\prod_{i=1}^{n-1}\left(x_{1} \oplus b_{i}\right) \cdot\left(\mathfrak{G}_{w^{\prime}}\right)^{\star}
$$

By the induction hypothesis and Proposition 3.2, we have

$$
\mathfrak{G}_{w}=\mathcal{G}_{n-1}^{[1, n-1]}\left(\mathfrak{G}_{\lambda\left(w^{\prime}\right), f\left(w^{\prime}\right)}\right)^{\star}
$$

Since $f(w)_{1}=1$ and $\left(f\left(w^{\prime}\right)_{1}, \ldots, f\left(w^{\prime}\right)_{n-1}\right)=\left(f\left(w^{\prime}\right)_{2}-1, \ldots, f\left(w^{\prime}\right)_{n}-1\right)$, Proposition 3.3 implies the claim.

Theorem 4.3 Let $(\lambda, f)$ be a flagged partition of length $r$. Then we have

$$
G_{\lambda, f}=\pi_{w}\left(\prod_{i=1}^{r} \prod_{j=1}^{a_{i}}\left(x_{i} \oplus b_{j}\right)\right)
$$

where $a_{i}=\lambda_{i}+f_{i}-i$ for $1 \leq i \leq r$ and $w=\left(s_{r} s_{r+1} \cdots s_{f_{r}-1}\right) \cdots\left(s_{2} s_{3} \cdots s_{f_{2}-1}\right)$ $\left(s_{1} s_{2} \cdots s_{f_{1}-1}\right)$.

Proof The proof is by induction on $|f|$. If $|f|=1$, we see that $G_{\lambda, f}=x_{1} \oplus b_{1}$ and $w=$ id so that the claim is trivial. Suppose that $|f|>1$. If $f_{1}=1$, let $\lambda^{\prime}=$ $\left(\lambda_{2}, \ldots, \lambda_{r}\right)$ and $f^{\prime}=\left(f_{2}-1, \ldots, f_{r}-1\right)$. By Proposition 3.3, we have $G_{\lambda, f}$ $=\left(G_{\lambda^{\prime}, f^{\prime}}\right)^{\star} \prod_{j=1}^{\lambda_{1}}\left(x_{1} \oplus b_{j}\right)$. By the induction hypothesis, we can write $\left(G_{\lambda^{\prime}, f^{\prime}}\right)^{\star}=$ $\pi_{w}\left(\prod_{i=2}^{r} \prod_{j=1}^{a_{i}}\left(x_{i} \oplus b_{j}\right)\right)$. Since $s_{1}$ doesn't appear in the reduced expression of $w$, we obtain the desired formula. If $f_{1}>1$, since $\lambda^{\prime}=\left(\lambda_{1}+1, \lambda_{2}, \ldots, \lambda_{r}\right), f^{\prime}=$ ( $f_{1}-1, f_{2}, \ldots, f_{r}$ ), Proposition 3.4 and the induction hypothesis imply the claim:
$G_{\lambda, f}=\pi_{f_{1}-1} G_{\lambda^{\prime}, f^{\prime}}=\pi_{f_{1}-1} \pi_{w s_{f_{1}-1}}\left(\prod_{i=1}^{r} \prod_{j=1}^{a_{i}}\left(x_{i} \oplus b_{j}\right)\right)=\pi_{w}\left(\prod_{i=1}^{r} \prod_{j=1}^{a_{i}}\left(x_{i} \oplus b_{j}\right)\right)$.
This completes the proof.
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# Flag Bundles, Segre Polynomials, and Push-Forwards 

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#### Abstract

In this note, we give Gysin formulas for partial flag bundles for the classical groups. We then give Gysin formulas for Schubert varieties in Grassmann bundles, including isotropic ones. All these formulas are proved in a rather uniform way by using the step-by-step construction of flag bundles and the Gysin formula for a projective bundle. In this way we obtain a comprehensive list of new general formulas. The content of this paper was presented by Piotr Pragacz at the International Festival in Schubert Calculus in Guangzhou, November 6-10, 2017.


Keywords Push forward • Gysin maps • Segre polynomials • Classical flag bundles • Kempf-Laksov bundles • Schubert bundles

## 1 Introduction

Let $E \rightarrow X$ be a vector bundle of rank $n$ on a variety $X$ over an algebraically closed field. Let $\pi: \mathbf{F}(E) \rightarrow X$ be the bundle of flags of subspaces of dimensions $1,2, \ldots, n-1$ in the fibers of $E \rightarrow X$. The flag bundle $\mathbf{F}(E)$ is used, e.g., in splitting principle, a standard technique which reduces questions about vector bundles to the case of line bundles; namely the pullback bundle $\pi^{*} E$ decomposes as a direct sum of line bundles. One can construct $\mathbf{F}(E)$ inductively as a sequence of projective bundles, using the following iterative step, that decreases the rank by 1 . Let $p_{1}: \mathbf{P}(E) \rightarrow X$ denote the projective bundle of lines in $E$, and let $U_{1}:=\mathcal{O}_{\mathbf{P}(E)}(-1)$ denote the universal subbundle on $\mathbf{P}(E)$, then one has a universal exact sequence of vector bundles on $\mathbf{P}(E)$

[^1]$$
0 \rightarrow U_{1} \longrightarrow p_{1}^{*} E \longrightarrow Q_{n-1} \rightarrow 0
$$
where $Q_{n-1}$ (the universal quotient bundle on $\mathbf{P}(E)$ ) is a rank $n-1$ vector bundle. Replacing $E$ by $Q_{n-1}$, one obtains a universal subbundle on $\mathbf{P}\left(Q_{n-1}\right)$, together with a universal quotient bundle $Q_{n-2}$. Iterating this process until obtaining a quotient bundle $Q_{1}$ of rank one, one gets a sequence of projective bundles
\[

$$
\begin{equation*}
\mathbf{F}(E):=\mathbf{P}\left(Q_{2}\right) \xrightarrow{p_{n-1}} \ldots \rightarrow \mathbf{P}\left(Q_{n-1}\right) \xrightarrow{p_{2}} \mathbf{P}(E) \xrightarrow{p_{1}} X, \tag{1}
\end{equation*}
$$

\]

a flag bundle filtration

$$
0 \subsetneq\left(p_{n-1} \circ \cdots \circ p_{2}\right)^{*} U_{1} \subsetneq\left(p_{n-1} \circ \cdots \circ p_{3}\right)^{*} U_{2} \subsetneq \ldots \subsetneq U_{n-1} \subsetneq \pi^{*} E,
$$

where $U_{i} \rightarrow \mathbf{P}\left(Q_{n+1-i}\right)$ is the kernel bundle of the composition

$$
\left(p_{i} \circ \cdots \circ p_{1}\right)^{*} E \rightarrow Q_{n-1} \rightarrow \ldots \rightarrow Q_{n-i}
$$

and universal exact sequences of vector bundles on $\mathbf{P}\left(Q_{n-i+1}\right)$ :

$$
\begin{equation*}
0 \rightarrow U_{i} / p_{i}^{*} U_{i-1} \rightarrow p_{i}^{*} Q_{n-i+1} \rightarrow Q_{n-i} \rightarrow 0 \tag{2}
\end{equation*}
$$

In the Grothendieck group of $\mathbf{F}(E)$, one can write (droping the pullback notation)

$$
\pi^{*} E=U_{1}+U_{2} / U_{1}+\cdots+U_{n-1} / U_{n-2}+Q_{1}
$$

as the sum of (the pullback of) the different line bundles appearing in (1).
Now we would like to outline how to obtain a Gysin formula for the flag bundle $\pi: \mathbf{F}(E) \rightarrow X$ (cf. Example 1), and introduce some notation.

We shall work in the framework of intersection theory of [3]. Recall that a proper morphism $g: Y \rightarrow X$ of nonsingular algebraic varieties over an algebraically closed field yields an additive map $g_{*}: A^{\bullet} Y \rightarrow A^{\bullet} X$ of Chow groups induced by pushforward cycles, called the Gysin map. The theory developed in [3] allows also one to work with singular varieties, or with cohomology. In this note, $X$ will always be nonsingular.

For $E \rightarrow X$ a vector bundle, let $s(E)$ be the Segre class of $E$, that is the formal inverse of the Chern class $c(E)$ in the Chow ring of $X$. Let $\xi=c_{1}\left(\mathcal{O}_{\mathbf{P}(E)}(1)\right)$; then $A^{\bullet}(\mathbf{P}(E))$ is generated algebraically by $\xi$ over $A^{\bullet} X$-here we identify $A^{\bullet} X$ with a subring of $A^{\bullet}(\mathbf{P}(E))$-and

$$
\begin{equation*}
\left(p_{1}\right)_{*} \xi^{i}=s_{i-(n-1)}(E), \tag{3}
\end{equation*}
$$

cf. [3]. To obtain a Gysin formula for the sequence of projective bundles (1), it suffices to appropriately iterate formula (3). The intermediate formulas involve the individual Segre classes of the universal quotient bundles, that can be eliminated
using (2) and the Whitney sum formula. However, it seems rather difficult to obtain a universal formula in this way. A universal formula should hold for any polynomial in characteristic classes of universal vector bundles and depend explicitly on the Segre classes of the original bundle $E$. To obtain such a formula, we use the generating series of the Segre classes of the universal quotient bundles. A prototype is the reformulation of (3) in

$$
\begin{equation*}
\left(p_{1}\right)_{*} \xi^{i}=\left[t^{n-1}\right]\left(t^{i} s_{1 / t}(E)\right), \tag{4}
\end{equation*}
$$

where we consider the specialization in $x=1 / t$ of the Segre polynomial $s_{x}(E)=$ $\sum_{i} s_{i}(E) x^{i}$ and where for a monomial $m$ and a Laurent polynomial $P,[m](P)$ denotes the coefficient of $m$ in $P$. Formula (4) and the projection formula imply that for any polynomial $f$ in one variable with coefficients in $A^{\bullet} X$

$$
\begin{equation*}
\left(p_{1}\right)_{*} f(\xi)=\left[t^{n-1}\right]\left(f(t) s_{1 / t}(E)\right) . \tag{5}
\end{equation*}
$$

In this formula, (i) one does not need to expand $f$ into a combination of monomials; (ii) one uses the Segre polynomial that, like the total Segre class, is a group homomorphism from the Grothendieck group of $X$ to the multiplicative group of units with degree zero term $=1$ in $A^{\bullet} X$.

Iterating the Gysin formula (5) yields a closed universal Gysin formula for the flag bundle $\mathbf{F}(E) \rightarrow X$, as announced in Example 1.

It is clear that the outlined strategy of proof applies to more general step-by-step constructions than the construction (1) of the flag bundle $\mathbf{F}(E) \rightarrow X$. Considering the truncated composition $p_{k} \circ \cdots \circ p_{1}$ in (1) yields formulas for full flag bundles, i.e. bundles of flags of subspaces of dimensions $1,2, \ldots, k$ in the fibers, for $k=$ $1, \ldots, n-1$. Then, using certain commutative diagrams (see [1, (5)]), one extends these formulas to arbitrary partial flag bundles.

One other interesting generalization is to restrict to the zero locus of a section of some vector bundle at each step of the sequence of projective bundles. In other words, one can impose some geometric conditions that the subspaces of the flag have to satisfy. An illustrative example is Theorem 2.3, in the orthogonal setting, obtained by considering at each step quadric bundles of isotropic lines in projective bundles of lines.

This method of step-by-step construction of generalized flag bundles leads to uniform short proofs of the different results announced in this note.

This note is organized as follows. In Sect. 2, we shall announce universal Gysin formulas for partial flag bundles for general linear groups, symplectic groups and orthogonal groups. The proofs of the results announced there can be found in [1].

In Sect. 3 we give Gysin formulas for Kempf-Laksov flag bundles. These generalized flag bundles are used to desingularize Schubert varieties in Grassmann bundles. Theorem 3.1 is established in [2]. Theorem 3.2 is announced for the first time in the present note.

## 2 Universal Gysin Formulas for Flag Bundles

In this section, the letter $f$ denotes a polynomial in the indicated number of variables with coefficients in $A^{\bullet} X$. The appropriate symmetries that $f$ has to satisfy to be in the Chow ring of the flag bundle under consideration are always implied. Here we consider non-singular varieties $X$. Note that the theory developed in Fulton's book [3] allows one to generalize the results to singular varieties over a field and their Chow groups; moreover, for complex varieties, one can also use the cohomology rings with integral coefficients.

We shall discuss separately the cases of general linear groups, symplectic groups and orthogonal groups.

### 2.1 General Linear Groups

Let $E \rightarrow X$ be a rank $n$ vector bundle. Let $1 \leq d_{1}<\cdots<d_{m}=d \leq n-1$ be a sequence of integers. We denote by $\pi: \mathbf{F}\left(d_{1}, \ldots, d_{m}\right)(E) \rightarrow X$ the bundle of flags of subspaces of dimensions $d_{1}, \ldots, d_{m}$ in $E$. On $\mathbf{F}\left(d_{1}, \ldots, d_{m}\right)(E)$, there is a universal flag $U_{d_{1}} \subsetneq \cdots \subsetneq U_{d_{m}}$ of subbundles of $\pi^{*} E$, where $\operatorname{rk}\left(U_{d_{k}}\right)=d_{k}$ (the fiber of $U_{d_{k}}$ over the point $\left(V_{d_{1}} \subsetneq \cdots \subsetneq V_{d_{m}} \subset E(x)\right.$ ), where $x \in X$, is equal to $V_{d_{k}}$ ). For a foundational account on flag bundles, see [4].

For $i=1, \ldots, d$, set $\xi_{i}=-c_{1}\left(U_{d+1-i} / U_{d-i}\right)$.
Theorem 2.1 With the above notation, for $f\left(\xi_{1}, \ldots, \xi_{d}\right) \in A^{\bullet}\left(\mathbf{F}\left(d_{1}, \ldots, d_{m}\right)(E)\right)$, one has

$$
\pi_{*} f\left(\xi_{1}, \ldots, \xi_{d}\right)=\left[t_{1}{ }^{e_{1}} \ldots t_{d}{ }^{e_{d}}\right]\left(f\left(t_{1}, \ldots, t_{d}\right) \prod_{1 \leq i<j \leq d}\left(t_{i}-t_{j}\right) \prod_{1 \leq i \leq d} s_{1 / t_{i}}(E)\right)
$$

where for $j=d-d_{k}+i$ with $i=1, \ldots, d_{k}-d_{k-1}$, we denote $e_{j}=n-i$.

Example 1 For the complete flag bundle $\pi: \mathbf{F}(E) \rightarrow X$, one has

$$
\pi_{*} f\left(\xi_{1}, \ldots, \xi_{n-1}\right)=\left[\prod_{i=1}^{n-1} t_{i}^{n-1}\right]\left(f\left(t_{1}, \ldots, t_{n-1}\right) \prod_{1 \leq i<j \leq n-1}\left(t_{i}-t_{j}\right) \prod_{i=1}^{n-1} s_{1 / t_{i}}(E)\right) ;
$$

and for the Grassmann bundle $\pi: \mathbf{F}(d)(E) \rightarrow X$, one has

$$
\pi_{*} f\left(\xi_{1}, \ldots, \xi_{d}\right)=\left[\prod_{i=1}^{d} t_{i}^{n-i}\right]\left(f\left(t_{1}, \ldots, t_{d}\right) \prod_{1 \leq i<j \leq d}\left(t_{i}-t_{j}\right) \prod_{i=1}^{d} s_{1 / t_{i}}(E)\right) .
$$

### 2.2 Symplectic Groups

Let $E \rightarrow X$ be a rank $2 n$ vector bundle equipped with a non-degenerate symplectic form $\omega: E \otimes E \rightarrow L$ (with values in a certain line bundle $L \rightarrow X$ ). We say that a subbundle $S$ of $E$ is isotropic if $S$ is a subbundle of its symplectic complement $S^{\omega}$, where

$$
S^{\omega}:=\{w \in E \mid \forall v \in S: \omega(w, v)=0\} .
$$

Let $1 \leq d_{1}<\cdots<d_{m} \leq n$ be a sequence of integers. We denote by $\pi$ : $\mathbf{F}^{\omega}\left(d_{1}\right.$, $\left.d_{m}\right)(E) \rightarrow X$ the bundle of flags of isotropic subspaces of dimensions $d_{1}, \ldots, d_{m}$ in $E$. On $\mathbf{F}^{\omega}\left(d_{1}, \ldots, d_{m}\right)(E)$, there is a universal flag $U_{d_{1}} \subsetneq \cdots \subsetneq U_{d_{m}}$ of subbundles of $\pi^{*} E$, where $\operatorname{rk}\left(U_{d_{k}}\right)=d_{k}$.

For $i=1, \ldots, d$, set $\xi_{i}=-c_{1}\left(U_{d+1-i} / U_{d-i}\right)$.
Theorem 2.2 With the above notation, for $f\left(\xi_{1}, \ldots, \xi_{d}\right) \in A^{\bullet}\left(\mathbf{F}^{\omega}\left(d_{1}, \ldots, d_{m}\right)(E)\right)$, one has
$\pi_{*} f\left(\xi_{1}, \ldots, \xi_{d}\right)=\left[t_{1} e_{1} \ldots t_{d} e_{d}\right]\left(f\left(t_{1}, \ldots, t_{d}\right) \prod_{1 \leq i<j \leq d}\left(c_{1}(L)+t_{i}+t_{j}\right)\left(t_{i}-t_{j}\right) \prod_{1 \leq i \leq d} s_{1 / t_{i}}(E)\right)$,
where for $j=d-d_{k}+i$ with $i=1, \ldots, d_{k}-d_{k-1}$, we denote $e_{j}=2 n-i$.

Example 2 For the symplectic Grassmann bundle $\pi: \mathbf{F}^{\omega}(d)(E) \rightarrow X$, where $\omega$ has values in a trivial line bundle, one has

$$
\pi_{*} f\left(\xi_{1}, \ldots, \xi_{d}\right)=\left[\prod_{i=1}^{d} t_{i}^{2 n-i}\right]\left(f\left(t_{1}, \ldots, t_{d}\right) \prod_{1 \leq i<j \leq d}\left(t_{i}^{2}-t_{j}^{2}\right) \prod_{i=1}^{d} s_{1 / t_{i}}(E)\right) .
$$

### 2.3 Orthogonal Groups

Let $E \rightarrow X$ be a vector bundle of rank $2 n$ or $2 n+1$ equipped with a non-degenerate orthogonal form $Q: E \otimes E \rightarrow L$ (with values in a certain line bundle $L \rightarrow X$ ). We say that a subbundle $S$ of $E$ is isotropic if $S$ is a subbundle of its orthogonal complement $S^{\perp}$, where

$$
S^{\perp}:=\{w \in E \mid \forall v \in S: Q(w, v)=0\} .
$$

Let $1 \leq d_{1}<\cdots<d_{m} \leq n$ be a sequence of integers. We denote by $\pi: \mathbf{F}^{Q}\left(d_{1}, \ldots, d_{m}\right)(E) \rightarrow X$ the bundle of flags of isotropic subspaces of dimensions $d_{1}, \ldots, d_{m}$ in $E$. On $\mathbf{F}^{Q}\left(d_{1}, \ldots, d_{m}\right)(E)$, there is a universal flag $U_{d_{1}} \subsetneq \cdots \subsetneq$ $U_{d_{m}}$ of subbundles of $\pi^{*} E$, where $\operatorname{rk}\left(U_{d_{k}}\right)=d_{k}$.

For $i=1, \ldots, d$, set $\xi_{i}=-c_{1}\left(U_{d+1-i} / U_{d-i}\right)$.

Theorem 2.3 With the above notation, for $f\left(\xi_{1}, \ldots, \xi_{d}\right) \in A^{\bullet}\left(\mathbf{F}^{Q}\left(d_{1}, \ldots, d_{m}\right)(E)\right)$, one has

$$
\begin{aligned}
& \pi_{*} f\left(\xi_{1}, \ldots, \xi_{d}\right)= \\
& \quad\left[t_{1}{ }^{\left.e_{1} \cdots t_{d}{ }^{e_{d}}\right]}\right]\left(f\left(t_{1}, \ldots, t_{d}\right) \prod_{1 \leq i \leq d}\left(2 t_{i}+c_{1}(L)\right) \prod_{1 \leq i<j \leq d}\left(c_{1}(L)+t_{i}+t_{j}\right)\left(t_{i}-t_{j}\right) \prod_{1 \leq i \leq d} s_{1 / t_{i}}(E)\right),
\end{aligned}
$$

where for $j=d-d_{k}+i$ with $i=1, \ldots, d_{k}-d_{k-1}$, we denote $e_{j}=\operatorname{rk}(E)-i$.
Note that, if the rank is $2 n$ and $d=n$, we consider both of the two isomorphic connected components of the flag bundle. Thus, if one is interested in only one of the two components, the result should be divided by 2 . When $c_{1}(L)=0$, this makes appear the usual coefficient $2^{n-1}$.

## 3 Universal Gysin Formulas for Kempf-Laksov Flag Bundles

In this section, we give Gysin formulas for Kempf-Laksov flag bundles, that are desingularizations of Schubert bundles in Grassmann bundles. We also extend the results to the symplectic setting. The orthogonal cases will be treated elsewhere.

### 3.1 General Linear Groups

Let $E \rightarrow X$ be a rank $n$ vector bundle on a variety $X$ with a reference flag of bundles $E_{1} \subsetneq \cdots \subsetneq E_{n}=E$ on it, where $\operatorname{rk}\left(E_{i}\right)=i$. Let $\pi: \mathbf{G}_{d}(E)=\mathbf{F}(d)(E) \rightarrow X$ be the Grassmann bundle of subspaces of dimension $d$ in the fibers of $E$. For any partition $\lambda \subseteq(n-d)^{d}$, there is the Schubert bundle $\varpi_{\lambda}: \Omega_{\lambda}\left(E_{\bullet}\right) \rightarrow X$ in $\mathbf{G}_{d}(E)$ given over the point $x \in X$ by

$$
\begin{equation*}
\Omega_{\lambda}\left(E_{\bullet}\right)(x):=\left\{V \in \mathbf{G}_{d}(E)(x): \operatorname{dim}\left(V \cap E_{n-d-\lambda_{i}+i}(x)\right) \geq i, \text { for } i=1, \ldots, d\right\} \tag{6}
\end{equation*}
$$

We denote by

$$
\left(v_{1}, \ldots, v_{d}\right):=\left(n-d-\lambda_{d}+d, \ldots, n-d-\lambda_{1}+1\right)
$$

the dimensions of the spaces of the reference flag involved in the definition of $\Omega_{\lambda}\left(E_{0}\right)$-in reverse order-. The partition $v$ is a strict partition, and furthermore, $n-i \leq \nu_{i} \leq \nu_{1}=n-\lambda_{d} \leq n$ for any $i$. Note that the above definition of $\Omega_{\lambda}\left(E_{\boldsymbol{\bullet}}\right)$ can be restated using $v$ with the conditions

$$
\begin{equation*}
\operatorname{dim}\left(V \cap E_{v_{i}}(x)\right) \geq d+1-i, \text { for } i=1, \ldots, d \tag{7}
\end{equation*}
$$

For a strict partition $\mu \subseteq(n)^{d}$ with $d$ parts, consider the flag bundle $\vartheta_{\mu}: F_{\mu}\left(E_{\bullet}\right) \rightarrow$ $X$ given over the point $x \in X$ by

$$
\begin{equation*}
F_{\mu}\left(E_{\bullet}\right)(x):=\left\{0 \subsetneq V_{1} \subsetneq \cdots \subsetneq V_{d} \in \mathbf{F}(1, \ldots, d)(E)(x): V_{d+1-i} \subseteq E_{\mu_{i}}(x), \text { for } i=1, \ldots, d\right\} \tag{8}
\end{equation*}
$$

We will call Kempf-Laksov flag bundles such bundles $\vartheta_{\mu}$ introduced in [5].
These appear naturally as desingularizations of Schubert bundles. For a partition $\lambda \subseteq(n-d)^{d}$, defining $v$ as above, by (7) the forgetful map $\mathbf{F}(1, \ldots, d)(E) \rightarrow$ $\mathbf{G}_{d}(E)$ induces a birational morphism $F_{v}\left(E_{\mathbf{\bullet}}\right) \rightarrow \Omega_{\lambda}\left(E_{\mathbf{\bullet}}\right)$. On the Schubert cell given over the point $x \in X$ by

$$
\stackrel{\circ}{\Omega}_{\lambda}\left(E_{\bullet}\right)(x):=\left\{V \in \mathbf{G}_{d}(E)(x): \operatorname{dim}\left(V \cap E_{\nu_{i}}(x)\right)=d+1-i, \text { for } i=1, \ldots, d\right\},
$$

which is open dense in $\Omega_{\lambda}\left(E_{\bullet}\right)$, the inverse map is $V \mapsto\left(V \cap E_{v_{d}}(x), \ldots, V \cap\right.$ $E_{\nu_{1}}(x)$ ). It establishes a desingularization of $\Omega_{\lambda}\left(E_{\bullet}\right)$ (see [5]).

We construct $F_{\mu}\left(E_{\bullet}\right)$ by induction on the length of flags. For $d=1$, it is simply $\mathbf{P}\left(E_{\mu_{d}}\right)$. Assume thus that for $d>1$ we have constructed the variety $F^{\prime} \subseteq$ $\mathbf{F}(1, \ldots, d-1)(E)$ parametrizing flags

$$
\left\{0 \subsetneq V_{1} \subsetneq \cdots \subsetneq V_{d-1} \in \mathbf{F}(1, \ldots, d-1)(E)(x): V_{d+1-i} \subseteq E_{\mu_{i}}(x), \text { for } i=2, \ldots, d\right\} .
$$

Let $U_{d-1}$ be the universal subbundle of rank $d-1$ on $\mathbf{F}(1, \ldots, d-1)(E)$. Note that in restriction to $F^{\prime}$, the condition $V_{d-1} \subseteq E_{\mu_{2}}(x)$ yields: $U_{d-1} \subseteq E_{\mu_{2}} \subseteq E_{\mu_{1}}$; we can therefore consider the subvariety

$$
\mathbf{P}\left(\left(\left.E_{\mu_{1}}\right|_{F^{\prime}}\right) /\left(\left.U_{d-1}\right|_{F^{\prime}}\right)\right) \subseteq \mathbf{P}\left(\left.\left(E / U_{d-1}\right)\right|_{F^{\prime}}\right) \subseteq \mathbf{P}\left(E / U_{d-1}\right)=\mathbf{F}(1, \ldots, d)(E)
$$

Iterating this inductive step, we get a sequence of projective bundles

$$
\begin{equation*}
F_{\mu}\left(E_{\bullet}\right)=\mathbf{P}\left(E_{\mu_{1}} / U_{d-1}\right) \rightarrow \mathbf{P}\left(E_{\mu_{2}} / U_{d-2}\right) \rightarrow \ldots \rightarrow \mathbf{P}\left(E_{\mu_{d-1}} / U_{1}\right) \rightarrow \mathbf{P}\left(E_{\mu_{d}}\right) . \tag{9}
\end{equation*}
$$

Set $\xi_{i}=-c_{1}\left(U_{d-i+1} / U_{d-i}\right), i=1, \ldots, d$, the first Chern class of the universal line bundle on $\mathbf{P}\left(E_{\mu_{i}} / U_{d-i}\right)$ (we only imply the restriction of the universal subbundles to this subvariety).

Let $f$ be a polynomial in $d$ variables with coefficients in $A^{\bullet}(X)$.

Theorem 3.1 With the above notation, one has

$$
\left(\vartheta_{\mu}\right)_{*} f\left(\xi_{1}, \ldots, \xi_{d}\right)=\left[t_{1}^{\mu_{1}-1} \cdots t_{d}^{\mu_{d}-1}\right]\left(f\left(t_{1}, \ldots, t_{d}\right) \prod_{1 \leq i<j \leq d}\left(t_{i}-t_{j}\right) \prod_{1 \leq i \leq d} s_{1 / t_{i}}\left(E_{\mu_{i}}\right)\right) .
$$

A proof of this theorem is based on (9) and (5).

### 3.2 Symplectic Groups

Let $E \rightarrow X$ be a rank $2 n$ symplectic vector bundle endowed with the symplectic form $\omega: E \otimes E \rightarrow L$ with value in a line bundle $L \rightarrow X$, over a variety $X$. For $d \in\{1, \ldots, n\}$, let $\mathbf{G}_{d}^{\omega}(E)=\mathbf{F}^{\omega}(d)(E)$ be the Grassmannian bundle of isotropic $d$-planes in the fibers of $E$. Let

$$
0=E_{0} \subsetneq E_{1} \subsetneq \cdots \subsetneq E_{n}=E_{n}^{\omega} \subsetneq \cdots \subsetneq E_{0}^{\omega}=E
$$

be a reference flag of isotropic subbundles of $E$ and their symplectic complements, where $\operatorname{rk}\left(E_{i}\right)=i$. For $i=1, \ldots, n$, we set $E_{n+i}:=E_{n-i}^{\omega}$. For a partition $\lambda \subseteq(2 n-$ $d)^{d}$, there is the Schubert cell $\Omega_{\lambda}\left(E_{\bullet}\right)$ in $\mathbf{G}_{d}^{\omega}(E)$ given over the point $x \in X$ by the conditions
$\Omega_{\lambda}\left(E_{\bullet}\right)(x):=\left\{V \in \mathbf{G}_{d}^{\omega}(E)(x): \operatorname{dim}\left(V \cap E_{2 n-d+i-\lambda_{i}}(x)\right)=i\right.$, for $\left.i=1, \ldots, d\right\}$.
Denote $v_{d+1-i}:=2 n-d+i-\lambda_{i}$ the dimension of the reference space appearing in the $i$ th condition. A partition indexing the Schubert cell $\Omega_{\lambda}$ must satisfy the conditions $v_{i}+v_{j} \neq 2 n+1$ (see [6, p. 174], where this is shown for $d=n$, and for arbitrary $d$ the argument is the same). For such partitions one defines the Schubert bundle $\varpi_{\lambda}: \Omega_{\lambda} \rightarrow X$ as the Zariski-closure of $\Omega_{\lambda}$, given over a point $x \in X$ by the conditions
$\Omega_{\lambda}\left(E_{\bullet}\right)(x):=\left\{V \in \mathbf{G}_{d}^{\omega}(E)(x): \operatorname{dim}\left(V \cap E_{2 n-d+i-\lambda_{i}}(x)\right) \geq i\right.$, for $\left.i=1, \ldots, d\right\}$.
For a strict partition $\mu \subseteq(2 n)^{d}$ with $d$ parts, such that $\mu_{i}+\mu_{j} \neq 2 n+1$ for all $i, j$, we introduce the isotropic Kempf-Laksov bundle $\vartheta_{\mu}: F_{\mu}\left(E_{\bullet}\right) \rightarrow X$ given over the point $x \in X$ by

$$
F_{\mu}\left(E_{\bullet}\right)(x):=\left\{0 \subsetneq V_{1} \subsetneq \cdots \subsetneq V_{d} \in \mathbf{F}^{\omega}(1, \ldots, d)(E)(x): V_{d+1-i} \subseteq E_{\mu_{i}}(x)\right\} .
$$

Note that as in the previous section, $F_{\nu}\left(E_{\bullet}\right)$ is birational to $\Omega_{\lambda}\left(E_{\bullet}\right)$, but here it is not smooth in general.

Let $U_{i}$ stands for the restriction to $F_{\mu}\left(E_{\mathbf{\bullet}}\right)$ of the rank $i$ universal bundle on $\mathbf{F}(1, \ldots, d)(E)$. Set $\xi_{i}=-c_{1}\left(U_{d-i+1} / U_{d-i}\right)$, for $i=1, \ldots, d$.

Let $f$ be a polynomial in $d$ variables with coefficients in $A^{\bullet}(X)$.
Theorem 3.2 With the above notation, one has

```
\(\left(\vartheta_{\mu}\right)_{*} f\left(\xi_{1}, \ldots, \xi_{d}\right)=\)
        \(\left[t_{1}^{\mu_{1}-1} \ldots t_{d}^{\mu_{d}-1}\right]\left(f\left(t_{1}, \ldots, t_{d}\right) \prod_{1 \leq i<j \leq d}\left(t_{i}-t_{j}\right) \prod_{\substack{1 \leq i<j \leq d \\ \mu_{i}+\mu_{j}>2 n+1}}\left(c_{1}(L)+t_{i}+t_{j}\right) \prod_{1 \leq j \leq d} s_{1 / t_{j}}\left(E_{\mu_{j}}\right)\right)\).
```

A proof of this theorem will appear in a separate publication.

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# Order of Tangency Between Manifolds 

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#### Abstract

We study the order of tangency between two manifolds of same dimension and give that notion three quite different geometric interpretations. Related aspects of the order of tangency, e.g., regular separation exponents, are also discussed.


Keywords Order of tangency • Order of contact • Taylor polynomial • Higher jets $\cdot$ Tower of Grassmannians • Regular separation exponent • Łojasiewicz exponent

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## 1 Introduction

In the present paper we discuss the order of tangency (or that of contact) between manifolds and its relation to enumerative geometry started with classical Schubert calculus.

Two plane curves, both sufficiently smooth and nonsingular at a point $x^{0}$, are said to have a contact of order at least $k$ at $x^{0}$ if, in properly chosen regular parametriza-

[^2]tions, those two curves have identical Taylor polynomials of degree $k$ about the respective preimages of $x^{0} .{ }^{1}$

Alternatively, those curves have such contact when their minimal regular separation exponent at $x^{0}$, cf. [11], is not smaller than $k$ or is not defined.

Formulas enumerating contacts have been widely investigated. For example in [3] the authors derive a formula for the number of contacts of order $n$ between members of a specified-parameter family of plane curves and a generic plane curve of sufficiently high degree.

Contact problems of this sort have been of both old and new interests, particularly in the light of Hilbert's 15th problem to make rigorous the classical calculations of enumerative geometry, especially those undertaken by Schubert [16]. The situation regarding ordinary (i.e., first-order) contacts between families of varieties is now well understood thanks in large measure to the contact formula of Fulton, Kleiman and MacPherson [6]. The above mentioned formula in [3] generalizes that given by Schubert in [17] for the number of double contacts between a given plane curve and a specified 2-parameter family of curves. Schubert made his computations through the use of what has come to be known as "Schubert triangles". This theory has been made completely rigorous by Roberts and Speiser, see, e.g., [15], and independently by Collino and Fulton [2].

Apart from contact formulas, an important role is played by the "order of tangency". Let us discuss this notion for Thom polynomials. Among important properties of Thom polynomials we record their positivity closely related to Schubert calculus (see, e.g., [12] and also [14] for a survey). Namely, the order of tangency allows one to define the jets of Lagrangian submanifolds. The space of these jets is a fibration over the Lagrangian Grassmannian and leads to a positive decomposition of the Lagrangian Thom polynomial in the basis of Lagrangian Schubert cycles.

In this paper, we give three approaches to the order of tangency. The first one (in Sect. 2) is by the Taylor approximations of local parametrizations of manifolds. The second one (a min-max procedure in Sect.3) makes use of curves sitting in the relevant manifolds. The third approach (in Sect.4) is by Grassmann bundles. We show that these three approaches are equivalent. We basically work with manifolds over the reals (of various classes of smoothness), but the results carry over-in the holomorphic category-to complex manifolds.

In the last two sections, we discuss some issues related to the "closeness" of pairs of geometric objects: branches of algebraic sets and relations with contact geometry. In fact, in Sect. 5 discussed are the regular separation exponents of pairs of semialgebraic sets, sometimes called Łojasiewicz exponents (not to be mixed with the by now classical exponent in the renowned Łojasiewicz inequality for analytic functions). Then, in Sect. 6 we report on an unexpected application of a modification of tangency order in 3D which yields an elegant criterion for a rank-2 distribution on a 3-manifold to be contact.

[^3]These concluding sections are not less important than the preceding ones. They show that the precise measurement of closeness is sometimes more demanding-and giving more-than merely bounding below tangency orders.

In the case of singular varieties different approaches to tangency orders lead to different notions. In this respect we refer the reader to [5] where compared were two discrete symplectic invariants of singular curves: the Lagrangian tangency order and index of isotropy.

## 2 By Taylor

One situation that is frequently encountered at the crossroads of geometry and analysis deals with pairs of manifolds which are the graphs of functions of the same number of variables. Such graphs can intersect, or touch each other, at a prescribed point, with various degrees of closeness.

Our departing point is a definition of such proximity going precisely in the spirit of a benchmark reference book [9], p. 18, although not formulated expressis verbis there.
Definition Two manifolds $M$ and $\tilde{M}$ embedded in $\mathbb{R}_{\tilde{M}}^{m}$, both of class $\mathrm{C}^{r}, r \geq 1$, and the same dimension $p$, intersecting at $x^{0} \in M \cap \widetilde{M}$, for $k \leq r$, have at $x^{0}$ the order of tangency at least $k$, when there exist a neighbourhood $U \ni u^{0}$ in $\mathbb{R}^{p}$ and parametrizations ${ }^{2}$ (diffeomorphisms onto the image)

$$
q:\left(U, u^{0}\right) \rightarrow\left(M, x^{0}\right), \quad \tilde{q}:\left(U, u^{0}\right) \rightarrow\left(\tilde{M}, x^{0}\right)
$$

of class $\mathrm{C}^{r}$ such that

$$
\begin{equation*}
(\tilde{q}-q)(u)=\mathrm{o}\left(\left|u-u^{0}\right|^{k}\right) \tag{1}
\end{equation*}
$$

when $U \ni u \rightarrow u^{0}$.
(We underline the existence clause in this definition. Supposing having already such a couple of local parametrizations $q$ and $\tilde{q}$, there is an abundance of other pairs of $\mathrm{C}^{r}$ parametrizations serving the vicinities of $x^{0}$ in $M$ and $\widetilde{M}$, respectively, and not satisfying the condition (1). Note also that in this definition the order of tangency is automatically at least 0 .)

Below in $\bullet \bullet$ in Sect. 3, and also in Sect. 4 we restrict ourselves to parametrizations of very specific type-just the graphs of $\mathrm{C}^{r}$ mappings going from $p$ dimensions to $m-p$ dimensions. This appears to be possible while not violating the key condition (1).

Naturally enough, the notion of the order of tangency not smaller than ... is invariant under the local $\mathrm{C}^{r}$ diffeomorphisms of neighbourhoods in $\mathbb{R}^{m}$ of the tangency point $x^{0}$.

[^4]Attention. In the real $\mathrm{C}^{\infty}$ category it is possible for the order of tangency to be at least $k$ for all $k \in \mathbb{N}$. In other words-be infinite even though $\left\{x^{0}\right\}=M \cap \widetilde{M}$. The rest of this paper is to be read with this remark in mind.

As a matter of record, basically the same definition is evoked in Proposition on page 4 in [8]. In [8] there is also proposed the following reformulation of (1).

Proposition 1 The condition (1) is equivalent to

$$
\begin{equation*}
T_{u^{0}}^{k}(q)=T_{u^{0}}^{k}(\tilde{q}), \tag{2}
\end{equation*}
$$

where $T_{u^{0}}^{k}(\cdot)$ means the Taylor polynomial about $u^{0}$ of degree $k$.
Implication (1) $\Rightarrow(2)$.

$$
\begin{gather*}
\mathrm{o}\left(\left|u-u^{0}\right|^{k}\right)=\tilde{q}(u)-q(u)=\left(\tilde{q}(u)-T_{u^{0}}^{k}(\tilde{q})\left(u-u^{0}\right)\right) \\
+\left(T_{u^{0}}^{k}(\tilde{q})\left(u-u^{0}\right)-T_{u^{0}}^{k}(q)\left(u-u^{0}\right)\right)+\left(T_{u^{0}}^{k}(q)\left(u-u^{0}\right)-q(u)\right), \tag{3}
\end{gather*}
$$

where the first and last summands on the right hand side are o $\left(\left|u-u^{0}\right|^{k}\right)$ by Taylor. So is the middle summand

$$
T_{u^{0}}^{k}(\tilde{q})\left(u-u^{0}\right)-T_{u^{0}}^{k}(q)\left(u-u^{0}\right)=\mathrm{o}\left(\left|u-u^{0}\right|^{k}\right)
$$

and (2) follows from the following general result.
Lemma 1 Let $w \in \mathbb{R}\left[u_{1}, u_{2}, \ldots, u_{p}\right]$, $\operatorname{deg} w \leq k, w(u)=\mathrm{o}\left(|u|^{k}\right)$ when $u \rightarrow 0$ in $\mathbb{R}^{p}$. Then $w$ is identically zero.

The proof goes by induction on $k \geq 0$, with an obvious start for $k=0$. Then, assuming this for the polynomials of degrees smaller than $k \geq 1$ and taking a polynomial $w$ of degree $k$ as in the wording of the lemma, we can assume without loss of generality that $w$ is homogeneous of degree $k$ (the terms of lower degrees vanish altogether by the inductive assumption). Let $u \in \mathbb{R}^{p},|u|=1$, be otherwise arbitrary. Then

$$
t^{k} w(u)=w(t u)=\mathrm{o}\left(|t u|^{k}\right)=\mathrm{o}\left(|t|^{k}\right) \quad \text { when } t \rightarrow 0 .
$$

Hence $w(u)=0$ and the vanishing of $w$ follows.
Implication $(1) \Leftarrow(2)$.
This implication is obvious, because now the middle term on the right hand side of (3) vanishes, so that the right hand side is automatically $\mathrm{o}\left(\left|u-u^{0}\right|^{k}\right)$.

## 3 By Curves

In the discussion in this section important will be the quantity

$$
\begin{equation*}
s:=\sup \{k \in \mathbb{N}: \text { the order of tangency } \geq k\} . \tag{4}
\end{equation*}
$$

(Note that an additional restriction here on $k$ is $k \leq r$, cf. Definition above.) If the class of smoothness $r=\infty$, then, by the very definition, the condition (1) holds for all $k$ if and only if $s=\infty$.

Is it possible to ascertain something similar in the finite-order-of-tangency case?
With an answer to this question in view, we stick in the present section to the notation introduced in Sect. 2, but assume additionally that

$$
\begin{equation*}
s<r \tag{5}
\end{equation*}
$$

(Reiterating, the quantity $s$ is defined in (4) above, and $r$ is the assumed class of smoothness of the underlying manifolds, finite or infinite when the category is real. When $r=\infty$, the condition (5) simply says that $s$ is finite.)

Our second approach uses pairs of curves lying, respectively, in $M$ and $\widetilde{M}$.
We naturally assume that $T_{x^{0}} M=T_{x^{0}} \tilde{M}$. Our actual objective is to show that
Theorem 1 Under (5),

$$
\begin{equation*}
\min _{v}\left(\max _{\gamma, \tilde{\gamma}}\left(\max \left\{l \in\{0\} \cup \mathbb{N}:|\gamma(t)-\tilde{\gamma}(t)|=\mathrm{o}\left(|t|^{l}\right) \text { when } t \rightarrow 0\right\}\right)\right)=s \tag{6}
\end{equation*}
$$

The minimum is taken over all $0 \neq v \in T_{x^{0}} M=T_{x^{0}} \tilde{M}$. The outer maximum is taken over all pairs of $\mathrm{C}^{r}$ curves $\gamma \subset M, \tilde{\gamma} \subset \tilde{M}$ such that $\gamma(0)=x^{0}=\tilde{\gamma}(0)$, andboth non-zero!-velocities $\dot{\gamma}(0), \dot{\tilde{\gamma}}(0)$ are both parallel to $v$.

Attention. In this theorem the assumption (5) is essential; our proof would not work in the situation $s=r$.

Proof of Theorem $1 \bullet$ It is quick to show that the integer on the left hand side of equality (6) is at least $s$. Indeed, for every fixed vector $v$ as above, $v=d q\left(u^{0}\right) u$ (without loss of generality, $u$ is like in the proof of Lemma 1), one can take $\delta(t)=$ $q\left(u^{0}+t u\right)$ and $\tilde{\delta}(t)=\tilde{q}\left(u^{0}+t u\right)$. Then

$$
|\delta(t)-\tilde{\delta}(t)|=\mathrm{o}\left(|t u|^{s}\right)=\mathrm{o}\left(|t|^{s}\right)
$$

and so

$$
\max _{\gamma, \tilde{\gamma}}\left(\max \left\{l:|\gamma(t)-\tilde{\gamma}(t)|=\mathrm{o}\left(|t|^{l}\right) \text { when } t \rightarrow 0\right\}\right) \geq s .
$$

In view of the arbitrariness in our choice of $v$, the same remains true after taking the minimum over all admissible $v$ 's which is actually done on the left hand side of (6).
$\bullet$ To show the opposite non-sharp inequality in (6) is more involved. It is precisely in this part that the additional assumption $s \leq r-1$ is needed. We study the two manifolds in the vicinity of $x^{0}$ via an appropriate local $\mathrm{C}^{r}$ diffeomorphism of the ambient space, after which

$$
\left(M, x^{0}\right)=\left(\left\{x_{p+1}=x_{p+2}=\cdots=x_{m}=0\right\}, 0\right)
$$

and

$$
\left(\tilde{M}, x^{0}\right)=\left(\left\{x_{j}=F^{j}\left(x_{1}, x_{2}, \ldots, x_{p}\right), j=p+1, p+2, \ldots, m\right\}, 0\right)
$$

for some $\mathrm{C}^{r}$ functions $F^{j}$. Having the manifolds so neatly (graph-like) positioned, we take the most adapted parametrizations

$$
\begin{gathered}
q\left(u_{1}, u_{2}, \ldots, u_{p}\right)=\left(u_{1}, u_{2}, \ldots, u_{p}, 0,0, \ldots, 0\right), \\
\tilde{q}\left(u_{1}, u_{2}, \ldots, u_{p}\right)=\left(u_{1}, u_{2}, \ldots, u_{p}, F\left(u_{1}, u_{2}, \ldots, u_{p}\right)\right),
\end{gathered}
$$

where $F=\left(F^{p+1}, F^{p+2}, \ldots, F^{m}\right)$. This-important-necessitates some extra technical work. Firstly the initial couple of parametrizations satisfying (1) is being straightened simultaneously with manifolds $M$ and $\widetilde{M}$. Naturally enough, the resulting parametrizations keep satisfying (1), but are not yet of the above-desired form. So the parametrizations and manifolds are to be additionally slightly upgraded via another local $\mathrm{C}^{r}$ ambient diffeomorphism so as (a) to keep the simple description of manifolds and (b) to have the eventual parametrizations adapted as desired above.

Given the definition (4) of $s$, there hold

$$
T_{u^{0}}^{s}(q)=T_{u^{0}}^{s}(\tilde{q}) \quad \text { and } \quad T_{u^{0}}^{s+1}(q) \neq T_{u^{0}}^{s+1}(\tilde{q}),
$$

that is,

$$
T_{u^{0}}^{S}(F)=0 \quad \text { and } \quad T_{u^{0}}^{s+1}(F) \neq 0
$$

It follows that there exist an integer $j \in\{p+1, p+2, \ldots, m\}$ and a vector $\mathbf{w} \in \mathbb{R}^{p}$ such that

$$
\begin{equation*}
T_{u^{0}}^{s}\left(F^{j}\right)(\mathbf{w})=0 \quad \text { and } \quad T_{u^{0}}^{s+1}\left(F^{j}\right)(\mathbf{w}) \neq 0 \tag{7}
\end{equation*}
$$

Now let $u$ and $\tilde{u}$ be two $\mathrm{C}^{r}$ curves in $\mathbb{R}^{p}$ passing at $t=0$ through $u^{0}$ and such that the vectors $\dot{u}(0)$ and $\dot{\tilde{u}}(0)$ are both non-zero and parallel to $\mathbf{w}$. These curves in parameters give rise to $\mathrm{C}^{r}$ curves $\delta(t)=q(u(t))$ and $\tilde{\delta}(t)=\tilde{q}(\tilde{u}(t))$ in the manifolds, both having at $t=0$ non-zero speeds parallel to the vector $\mathbf{v}$ : $=d q\left(u^{0}\right) \mathbf{w}=d \tilde{q}\left(u^{0}\right) \mathbf{w}$. We will now estimate from above (by $s$ ) the left hand side of the equality (6) using, no wonder, $\mathbf{v}, \delta$, and $\tilde{\delta}$ :

$$
\begin{align*}
|\delta(t)-\tilde{\delta}(t)|=\sqrt{|u(t)-\tilde{u}(t)|^{2}+|F(\tilde{u}(t))|^{2}} & \\
& \geq|F(\tilde{u}(t))| \geq\left|F^{j}(\tilde{u}(t))\right| \neq \mathrm{o}\left(|t|^{s+1}\right) \tag{8}
\end{align*}
$$

where the last inequality necessitates an explanation. In fact, by (7) and for every $c \neq 0$

$$
T_{u^{0}}^{s+1}\left(F^{j}\right)(t c \mathbf{w})=(c t)^{s+1} T_{u^{0}}^{s+1}\left(F^{j}\right)(\mathbf{w}) \neq \mathrm{o}\left(|t|^{s+1}\right) \quad \text { when } t \rightarrow 0 .
$$

But $\tilde{u}(t)-\tilde{u}(0)=c t \mathbf{w}+\mathrm{o}(|t|)$ for some non-zero $c$, hence

$$
T_{u^{0}}^{s+1}\left(F^{j}\right)(\tilde{u}(t)-\tilde{u}(0)) \neq \mathrm{o}\left(|t|^{s+1}\right) \quad \text { when } t \rightarrow 0
$$

as well. Also, just by Peano in the class of smoothness $s+1 \leq r$, cf. (5),

$$
F^{j}(u)=T_{u^{0}}^{s+1}\left(F^{j}\right)\left(u-u^{0}\right)+\mathrm{o}\left(\left|u-u^{0}\right|^{s+1}\right)
$$

when $u \rightarrow u^{0}$ in $\mathbb{R}^{p}$. Therefore, $F^{j}(\tilde{u}(t)) \neq \mathrm{o}\left(|t|^{s+1}\right)$ as written in (8).
Now it is important to note that the pair of curves $\delta$ and $\tilde{\delta}$ produced by us above is completely general in the category $\mathrm{C}^{r}$ for that chosen vector $\mathbf{v}$. Hence it follows that-for this precise vector $\mathbf{v}$ !-the quantity

$$
\max _{\gamma, \tilde{\gamma}}\left(\max \left\{l \in\{0\} \cup \mathbb{N}:|\gamma(t)-\tilde{\gamma}(t)|=\mathrm{o}\left(|t|^{l}\right) \text { when } t \rightarrow 0\right\}\right)
$$

does not exceed $s$. Understandingly, so does the minimum of such quantities over all $v$ 's in $T_{x^{0}} M=T_{x^{0}} \widetilde{M}$. Theorem 1 is now proved.

## 4 By Grassmannians

Our third approach is based on the introductory pages of [8] where a natural tower of consecutive Grassmannians is being attached to every given local $\mathrm{C}^{r}$ parametrization $q$ as used by us in the preceding sections. However, to allow for a recursive definition of tower's members, a more general framework is needed.

Namely, to every $\mathrm{C}^{1}$ immersion $H: N \rightarrow N^{\prime}, N$-an $n$-dimensional manifold, $N^{\prime}$ —an $n^{\prime}$-dimensional manifold (manifolds not necessarily embedded in Euclidean spaces!), we attach the so-called image map $\mathcal{G} H: N \rightarrow G_{n}\left(N^{\prime}\right)$ of the tangent map $d H:$ for $s \in N$,

$$
\begin{equation*}
\mathcal{G} H(s)=d H(s)\left(T_{s} N\right), \tag{9}
\end{equation*}
$$

where $G_{n}\left(N^{\prime}\right)$ is the total space of the Grassmann bundle, with base $N^{\prime}$, of all $n$-planes tangent to $N^{\prime}$. That is, $G_{n}\left(N^{\prime}\right)$ is a new manifold, much bigger than $N^{\prime}$ (whenever $n^{\prime}>n$ ), of dimension $n^{\prime}+n\left(n^{\prime}-n\right)$.

We stick in the present section to the notation from Sect. 2 and invariably use the pair of parametrizations $q$ and $\tilde{q}$. So we are given the mappings

$$
\mathcal{G} q: U \longrightarrow G_{p}\left(\mathbb{R}^{m}\right), \quad \mathcal{G} \tilde{q}: U \longrightarrow G_{p}\left(\mathbb{R}^{m}\right)
$$

Upon putting $M^{(0)}=\mathbb{R}^{m}, \mathcal{G}^{(1)}=\mathcal{G}$, there emerge two sequences of recursively defined mappings. Namely, for $l \geq 1$,

$$
\mathcal{G}^{(l)} q: U \longrightarrow G_{p}\left(M^{(l-1)}\right), \quad \mathcal{G}^{(l+1)} q=\mathcal{G}\left(\mathcal{G}^{(l)} q\right)
$$

and

$$
\mathcal{G}^{(l)} \tilde{q}: U \longrightarrow G_{p}\left(M^{(l-1)}\right), \quad \mathcal{G}^{(l+1)} \tilde{q}=\mathcal{G}\left(\mathcal{G}^{(l)} \tilde{q}\right)
$$

where, naturally, $M^{(l)}=G_{p}\left(M^{(l-1)}\right)$. Now our objective is to show the following.
Theorem $2 \mathrm{C}^{r}$ manifolds $M$ and $\tilde{M}$ have at $x^{0}$ the order of tangency at least $k$ $(1 \leq k \leq r)$ iff there exist $\mathrm{C}^{r}$ parametrizations $q$ and $\tilde{q}$ of the vicinities of $x^{0}$ in, respectively, $M$ and $\widetilde{M}$, such that

$$
\begin{equation*}
\mathcal{G}^{(k)} q\left(u^{0}\right)=\mathcal{G}^{(k)} \tilde{q}\left(u^{0}\right) . \tag{10}
\end{equation*}
$$

(Observe that, in (10), there is clearly encoded that $q\left(u^{0}\right)=x^{0}=\tilde{q}\left(u^{0}\right)$. )

### 4.1 Proof of Theorem 2

In what follows, of interest for us will be the situations when $H$ in (9) above is locally (and all is local in tangency considerations!) the graph of a $\mathrm{C}^{1}$ mapping $h: \mathbb{R}^{p} \supset U \rightarrow \mathbb{R}^{t}$. That is, for $u \in U, H(u)=(u, h(u)) \in \mathbb{R}^{p+t}=\mathbb{R}^{p} \times \mathbb{R}^{t}$. Then (9) assumes by far more precise form

$$
\begin{equation*}
\mathcal{G} H(u)=(u, h(u) ; d(u, h(u))(u))=\left(u, h(u) ; \operatorname{span}\left\{\partial_{j}+h_{j}(u): j=1,2, \ldots, p\right\}\right) \tag{11}
\end{equation*}
$$

where the symbol $h_{j}$ means the partial derivative of a vector mapping $h$ with respect to the variable $u_{j}(j=1, \ldots, p)$, and $\partial_{j}+h_{j}(u)$ denotes the partial derivative of the vector mapping $(\iota, h): U \rightarrow \mathbb{R}^{p}\left(u_{1}, \ldots, u_{p}\right) \times \mathbb{R}^{t}$ with respect to $u_{j}$, where $\iota: U \hookrightarrow \mathbb{R}^{p}$ is the inclusion.

Now observe that the expression for $\mathcal{G} H(u)$ on the right hand side of (11) is still not quite useful. Yet there are charts in each newly appearing Grassmannian (see, for instance, [8] or p. 46 in [1])!

The chart in a typical fibre $G_{p}$ over a point in the base $\mathbb{R}^{p+t}$, good for (11), consists of all the entries in the bottommost rows (indexed by numbers $p+1, p+$ $2, \ldots, p+t)$ in the $(p+t) \times p$ matrices

$$
\left[v_{1}\left|v_{2}\right| \ldots \mid v_{p}\right]
$$

with non-zero upper $p \times p$ minor, after multiplying the matrix on the right by the inverse of that upper $p \times p$ submatrix. That is to say, taking as the local coordinates all the entries in rows $(p+1)$-st, $\ldots,(p+t)$ th of the matrix

$$
\left[v_{j}^{i}\right]_{\substack{1 \leq i \leq p+t \\ 1 \leq j \leq p}}\left(\left[v_{j}^{i}\right]_{\substack{1 \leq i \leq p \\ 1 \leq j \leq p}}\right)^{-1} .
$$

That is, these coordinates are all $t \times p$ entries of the matrix

$$
\left[v_{j}^{i}\right]_{\substack{p+1 \leq i \leq p+t \\ 1 \leq j \leq p}}\left(\left[v_{j}^{i}\right]_{\substack{1 \leq i \leq p \\ 1 \leq j \leq p}}\right)^{-1} .
$$

In these, extremely useful, glasses the description (11) gets stenographed to

$$
\begin{equation*}
\mathcal{G} H(u)=\left(u, h(u) ; \frac{\partial h}{\partial u}(u)\right), \tag{12}
\end{equation*}
$$

where, under the symbol $\frac{\partial h}{\partial u}(u)$ understood are all the entries of this Jacobian $(t \times p)$ matrix written in row and separated by commas. This technical simplification is central for a proof that follows.

After this, basically algebraic, preparation we come back to Theorem 2.
The order of tangency between $M$ and $\widetilde{M}$ at $x^{0}$ being at least $k$ precisely means (Proposition 1) the existence of local $\mathrm{C}^{r}$ parametrizations $q$ and $\tilde{q}$ satisfying (2). So we are just going to show that $(10) \Longleftrightarrow$ (2).

Moreover, we assume without loss of generality-much like it has been the case in the part $\bullet$ of the proof of Theorem 1-that $M$ and $\widetilde{M}$ are locally the graphs of parametrizations $q$ and $\tilde{q}$, respectively. Which, at the same time, satisfy (2). So (2) holds for $q(u)=(u, f(u)), f: U \rightarrow \mathbb{R}^{m-p}\left(y_{p+1}, \ldots, y_{m}\right)$ and for $\tilde{q}(u)=$ $(u, \tilde{f}(u)), \tilde{f}: U \rightarrow \mathbb{R}^{m-p}\left(y_{p+1}, \ldots, y_{m}\right), x^{0}=\left(u^{0}, f\left(u^{0}\right)\right)=\left(u^{0}, \tilde{f}\left(u^{0}\right)\right)$.

Implication (2) $\Rightarrow$ (10).
We will derive such expressions for $\mathcal{G}^{(k)} q(u)$ and $\mathcal{G}^{(k)} \tilde{q}(u), u \in U$, that the use of the condition (2) will just prompt by itself. An added value of this derivation will be the control over the sets of natural local coordinates in the Grassmannians in question. (With this information at hand the opposite implication $(2) \Leftarrow(10)$ will follow in no time.) Our main technical tool for the $\Rightarrow$ implication is

Lemma 2 For $1 \leq l \leq k$ there exists a local chart on the Grassmannian space $G_{p}\left(M^{(l-1)}\right)$ in which the mapping $\mathcal{G}^{(l)} q$ evaluated at $u$ assumes the form

$$
\left(u, f(u) ;\binom{l}{1} \times f_{[1]}(u),\binom{l}{2} \times f_{[2]}(u), \ldots,\binom{l}{l} \times f_{[l]}(u)\right),
$$

where $f_{[\nu]}(u)$ is a shorthand notation for the aggregate of all the partials of the $\nu t h$ order at $u$, of all the components of $f$, which are in the number $(m-p) \times p^{\nu}$, and the symbol $N \times(*)$ stands for the $N$ copies going in row and separated by commas, of an object ( $*$ ).

Attention. In this lemma we purposefully distinguish mixed derivatives taken in different orders, simply disregarding the Schwarz symmetricity discovery.

Proof $l=1$. We note that

$$
\mathcal{G}^{(1)} q(u)=\left(u, f(u) ; \operatorname{span}\left\{\partial_{j}+f_{j}(u): j=1,2, \ldots, p\right\}\right),
$$

in the relevant Jensen-Borisenko-Nikolaevskii chart, is nothing but

$$
\left(u, f(u) ; f_{[1]}(u)\right)=\left(u, f(u) ;\binom{l}{1} \times f_{[1]}(u)\right)
$$

The beginning of induction is done.
$l \Rightarrow l+1, l<k$. The mapping $\mathcal{G}^{(l)} q: U \rightarrow M^{(l)}$, evaluated at $u$, is already written down, in appropriate local chart assumed to exist in $M^{(l)}$, as

$$
\begin{equation*}
\left(u, f(u),\binom{l}{1} \times f_{[1]}(u),\binom{l}{2} \times f_{[2]}(u), \ldots,\binom{l}{l} \times f_{[l]}(u)\right) . \tag{13}
\end{equation*}
$$

We work with $\mathcal{G}^{(l+1)} q=\mathcal{G}\left(\mathcal{G}^{(l)} q\right)$. Now, (13) being clearly of the form $H(u)=$ $(u, h(u))$ in the previously introduced notation, the mapping $h$ reads

$$
h(u)=\left(f(u),\binom{l}{1} \times f_{[1]}(u),\binom{l}{2} \times f_{[2]}(u), \ldots,\binom{l}{l} \times f_{[l]}(u)\right)
$$

In order to have $\mathcal{G} H(u)$ written down, in view of (12), one ought to write in row: $u$, then $h(u)$, and then all the entries of the Jacobian matrix $\frac{\partial h}{\partial u}(u)$, also written in row and separated by commas. The latter, in our shorthand notation, are computed immediately. Namely
$\frac{\partial h}{\partial u}(u)=\left(\binom{l}{0} \times f_{[1]}(u),\binom{l}{1} \times f_{[2]}(u),\binom{l}{2} \times f_{[3]}(u), \ldots,\binom{l}{l} \times f_{[l+1]}(u)\right)$.
These entries on the right hand side are to be juxtaposed with the former entries $(u, h(u))$. For better readability, we put together the groups of same partials (a yet another permutation of Grassmann-type coordinates, cf. the wording of the lemma). In view of the elementary identities $\binom{l}{\nu-1}+\binom{l}{\nu}=\binom{l+1}{\nu}$, we get in the outcome

$$
\left(u, f(u),\binom{l+1}{1} \times f_{[1]}(u),\binom{l+1}{2} \times f_{[2]}(u), \ldots,\binom{l+1}{l} \times f_{[l]}(u),\binom{l+1}{l+1} \times f_{[l+1]}(u)\right) .
$$

Lemma 2 is now proved by induction.
We now take $l=k$ in Lemma 2 and get, for arbitrary $u \in U$, two similar visualisations of $\mathcal{G}^{(k)} q(u)$ and $\mathcal{G}^{(k)} \tilde{q}(u)$. At that, the equality (2) holds true at $u=u^{0}$. As a consequence, (10) follows.

Implication $(2) \Leftarrow(10)$.
With the information on superpositions of the mappings $\mathcal{G}$, gathered in the course of proving the implication $(2) \Rightarrow(10)$, this opposite implication is clear. Theorem 2 is now proved.

## 5 Algebraic Geometry Examples and Regular Separation Exponents

In the present section, we shall work with the regular separation exponents of pairs of sets, a notion due to Łojasiewicz [11]. We shall compute these exponents in several natural examples, and compare the results with the information (when available) about the relevant orders of tangency. These examples deal with branches of algebraic sets which often happen to be tangent one to another, with various degrees of closeness.

Example 1 In the work [4] (Fig. 2 on page 37 there) analyzed is the following algebraic set in $\mathbb{R}^{2}(x, y)$

$$
\begin{equation*}
\mathcal{C}=\left\{(x, y):\left(y-x^{2}\right)^{2}=x^{5}\right\} . \tag{14}
\end{equation*}
$$

The two branches of $C$ issuing from the point $(0,0)$,

$$
C_{-}=\left\{y=x^{2}-x^{5 / 2}, x \geq 0\right\} \text { and } C_{+}=\left\{y=x^{2}+x^{5 / 2}, x \geq 0\right\},
$$

could be naturally extended to one-dimensional manifolds $D_{-}$and $D_{+}$, both of class $\mathrm{C}^{2}$-the graphs of functions

$$
y_{-}(x)=x^{2}-|x|^{5 / 2} \quad \text { and } \quad y_{+}(x)=x^{2}+|x|^{5 / 2}
$$

respectively. The Taylor polynomials of degree 2 about $x=0$ of $y_{-}$and $y_{+}$coincide. Hence $D_{-}$and $D_{+}$have at $(0,0)$ the order of tangency at least 2 (cf. Sect. 2), and clearly not at least 3 .

This example clearly suggests that, in real algebraic geometry, it would be pertinent to use non-integer measures of closeness. For instance, for the above sets $y_{-}(x)$ and $y_{+}(x)$, we may take

$$
\sup \left\{\alpha>0: y_{+}(x)-y_{-}(x)=\mathrm{o}\left(|x|^{\alpha}\right) \text { when } x \rightarrow 0\right\} .
$$

This kind of a generalized order of tangency would be $5 / 2$ in the Colley-Kennedy example.

In the local analytic geometry there is a precise name for this notion-the minimal regular separation exponent for two (semialgebraic) branches, say $X$ and $Y$, of an algebraic set. That is, here the minimal exponent for $X=C_{-}$and $Y=C_{+}$. Some authors working after [11] used to call such a quantity the Łojasiewicz exponent of, in this case, $C_{-}, C_{+}$at $(0,0)$. And denoted it-when specialized to the present situation-by $\mathcal{L}_{(0,0)}\left(C_{-}, C_{+}\right) .{ }^{3}$
Remark 1 (i) Therefore, in Example 1, the rational number 5/2 is the minimal regular separation exponent of the semialgebraic sets $C_{-}$and $C_{+}$which touch each other at $(0,0)$.
(ii) Example 1 quickly generalizes, by means of the equation $\left(y-x^{N}\right)^{2}=x^{2 N+1}$ with $N$ arbitrarily large, to yield a pair of $\mathrm{C}^{N}$ manifolds having the order of tangency at least $N$ and not at least $N+1$, and having the minimal regular separation exponent $N+\frac{1}{2}$.

It has not been difficult in Example 1 to discern the pair of branches $C_{-}$and $C_{+}$, initially slightly hidden in a synthetic equation (14). But it can happen considerably worse in this respect. Consider, for instance

Example 2 (a) an algebraic set in the plane $\mathbb{R}^{2}(x, y)$ defined by a single equation

$$
\begin{equation*}
(x y)^{2}=\frac{1}{4}\left(x^{2}+y^{2}\right)^{3} \tag{15}
\end{equation*}
$$

This set possesses a pair (even more than one such pair) of semialgebraic branches touching each other at the point $(0,0)$. Yet it is not so immediate to ascertain their minimal regular separation exponent. Only after recognizing in (15) the classical quatrefoil $x=\cos (\varphi) \sin (2 \varphi), y=\sin (\varphi) \sin (2 \varphi)$, it becomes quick to compute the relevant minimal regular separation exponent equal to 2 .
(b) It is even more interestingly with another algebraic set in 2D given by the equation

$$
\begin{equation*}
\left(x^{2}+y^{2}-\frac{1}{2} x\right)^{2}=\frac{1}{4}\left(x^{2}+y^{2}\right) . \tag{16}
\end{equation*}
$$

This set possesses as well a pair of semialgebraic branches $\{y \leq 0\}$ and $\{y \geq 0\}$ touching each other at the point $(0,0)$. Yet it takes some time to find their minimal regular separation exponent. In fact, after discovering in (16) the classical cardioid $r=\frac{1}{2}(1+\cos \varphi)$, that exponent turns out to be-one more time-a non-integer ( $3 / 2$, in the occurrence).

[^5](c) It is worthy of note that for both explicitly defined algebraic sets (15) and (16) one can apply general type bounds above, for the minimal separation exponent, produced in [10]. Yet the estimations got in that way are unrealistically (by factors of thousands) high.
Returning to the notion of the order of tangency, in the realm of algebraic geometry the distinction between that order and minimal regular separation exponent sometimes happens to be fairly clear, with both discussed quantities effectively computable. An instructive instance of such a situation occurs in [18] (Example 3.5 there).

Example 3 The author of [18] deals there with a pair of one-dimensional algebraic manifolds $N$ and $Z$ in $\mathbb{R}^{2}(x, y)$ intersecting at $(0,0)$. The manifold $N=\{y=0\}$ is already utmostly simplified, whereas $Z=\left\{y^{d}+y x^{d-1}+x^{s}=0\right\}$ depends on two integer parameters $d$ and $s, 1<d<s, d$ odd. What we are going to discuss here is a kind of reworking of Tworzewski's original approach, see also Attention below.

These manifolds have at $(0,0)$ the order of tangency at least $s-d$, and not at least $s-d+1$, while their minimal regular separation exponent at $(0,0)$ is $s-d+1$.

Indeed-to justify this one tries to present $Z$ as the graph of a function $y=y(x)$. Clearly, $y(0)=0$ and a function $y(x)$ could not be divisible by, for instance, $x^{s+1}$. So, with no loss of generality,

$$
y(x)=x^{k} z(x)-x^{s-d+1}
$$

for certain integer $k \geq 1$ and another function $z(x)$ such that $z(0) \neq 0$.

- The possibility $k<s-d+1$ boils rather quickly down to $k=1$, and then to the relation $\left(z-x^{s-d}\right)^{d}+z=0$, impossible at $x=0$, for $(z(0))^{d}+z(0) \neq 0, d$ being odd.
$\bullet$ So $k \geq s-d+1$ and now

$$
y(x)=x^{s-d+1} z(x)-x^{s-d+1}
$$

for a certain function $z(x)$. Upon substituting this $y(x)$ to the defining equation of $Z$ and simplifying, $(z-1)^{d} x^{(d-1)(s-d)}+z=0$. Hence $z(0)=0$. The Implicit Function Theorem is applicable here around $(0,0)$, because

$$
\left.\frac{\partial}{\partial z}\left((z-1)^{d} x^{(d-1)(s-d)}+z\right)\right|_{(0,0)}=1
$$

One gets a locally unique $\mathrm{C}^{\infty}$ function $z(x), z(0)=0$, hence also a locally unique function $y(x)=x^{s-d+1} z(x)-x^{s-d+1}$ whose graph is $Z$. Because the function $z$ vanishes at 0 , the minuend in this expression for $y$ is an ' $o$ 'of the subtrahend when $x \rightarrow 0$. So the statements about the order of tangency and minimal regular separation exponent follow immediately.

Attention. When $2 \mid d$, the above-found resolving function $y(x)$ is not the only solution to the defining equation of $Z$. Namely, the necessary equality $(z(0))^{d}+z(0)=0$, $z(0) \neq 0$, is then possible with $z(0)=-1$ and

$$
\left.\frac{\partial}{\partial z}\left(\left(z-x^{s-d}\right)^{d}+z\right)\right|_{(0,-1)}=d(-1)^{d-1}+1=1-d \neq 0
$$

Hence the Implicit Function Theorem gives this time a locally unique (for that $k=1$ ) $\mathrm{C}^{\infty}$ function $\tilde{z}(x), \tilde{z}(0)=-1$. Then the graph of

$$
\tilde{y}(x)=x \tilde{z}(x)-x^{s-d+1}=-x+(\text { higher powers of } x)
$$

is a second branch of $Z$ passing through $(0,0) \in \mathbb{R}^{2}$, transversal to $N$, in a stark distinction to the previously found, tangent to $N$, branch.

Remark 2 More generally, one could not hope to get a precise information regarding the minimal regular separation exponent for the pair of manifolds $M, \widetilde{M}$ on the sole basis of the assumptions in Theorem 1. That is, basically, under (5). Despite the inequality (8), that exponent need not necessarily be $s+1$. Following Example 1 earlier in this section, one could just take the curves $C_{-}$and $C_{+}$as the curves $\delta$ and $\tilde{\delta}$, respectively, in the proof of Theorem 1 . That is, to take $M=C_{-}=\delta$ and $\tilde{M}=C_{+}=\tilde{\delta}$. Then, as the reader has lately seen, $r=2$ and the quantity $s$ defined in (4) is also 2 , while the minimal regular separation exponent is but $s+\frac{1}{2}\left(=\frac{5}{2}\right)$.

Even restricting oneself to a benchmark setting $\operatorname{dim} M=\operatorname{dim} \tilde{M}$, an additional enormous complication could come from the fact that the intersection $M \cap \widetilde{M}$ might be a topologically highly nontrivial set (think about the $\mathrm{C}^{\infty}$ category). And it is precisely $M \cap \widetilde{M}$ which enters the definition of regular separation exponents for the pair $M, \tilde{M}$.

## 6 Relation with Contact Topology

Unsurprisingly, the notion of order of contact proves useful not only in algebraic geometry (cf. Introduction), but also in geometry tout court. One not so obvious application in the real category deals with the real contact structures in three dimensions. Our summarising it here follows closely Sect. 1.6 in [7]. The author considers there a couple $\Sigma \subset M$, where $M$ is a contact 3 -dimensional manifold and $\Sigma$-a fixed embedded surface in it. Contact means $M$ being endowed with a contact structure, say $\xi$, in $T M$.

When one approaches a given point $p \in \Sigma$ by points $q$ staying within $\Sigma$, a natural question is about the order of smallness of the angle $\angle\left(T_{q} \Sigma, \xi_{q}\right)$. If that angle is an 'O' of the distance of $q$ to $p$ taken to power $k$ (the distance measured in any chosen, and hence every, set of smooth local coordinates about $p$ ), then it is said that $\xi$ has the order of contact at least $k$ with $\Sigma$ at $p$. (Therefore, what is discussed in this section
differs a little from the notion of closeness of a pair of manifolds investigated in the preceding sections. Yet the added value is substantial.)

That is, to say that the new order of contact is at least 1 at a given point $p$ is tantamount to saying that $\xi_{p}=T_{p} \Sigma$. And it is exactly 0 at $p$ whenever $\xi_{p} \neq T_{p} \Sigma$.

So it comes as a not small surprise that this elementary notion allows one to characterise the contact structures as such! Namely, a theorem proved in [7] asserts that a rank-2 tangent distribution $\xi$ on a 3-dimensional $M$ is contact iff $\xi$ has the new order of contact at most 1 with every surface $\Sigma$ embedded in $M$, and this at every point of $\Sigma$.

The next natural question in this direction is whether it is possible to similarly characterise contact structures on $(2 n+1)$-dimensional manifolds, $n \geq 2$. The author of [7] says nothing in this respect.

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# On Schubert's Problem of Characteristics 

Haibao Duan and Xuezhi Zhao


#### Abstract

The Schubert varieties on a flag manifold $G / P$ give rise to a cell decomposition on $G / P$ whose Kronecker duals, known as the Schubert classes on $G / P$, form an additive base of the integral cohomology $H^{*}(G / P)$. The Schubert's problem of characteristics asks to express a monomial in the Schubert classes as a linear combination in the Schubert basis. We present a unified formula expressing the characteristics of a flag manifold $G / P$ as polynomials in the Cartan numbers of the group $G$. As application we develop a direct approach to our recent works on the Schubert presentation of the cohomology rings of flag manifolds $G / P$.


Keywords Lie group $\cdot$ Flag manifold $\cdot$ Schubert variety $\cdot$ Cohomology
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## 1 The Problem

Schubert considered what he called the problem of characteristics to be the main theoretical problem of enumerative geometry. —Kleimann [40, 1987]

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[^6]The existence of a finite basis for the homologies in every closed manifold implies furthermore the solvability of Schubert's "characteristics problems" in general. - Van der Waerden $[59,1930]^{1}$
Schubert calculus is the intersection theory of the 19th century, together with applications to enumerative geometry. Justifying this calculus was a major topic of the 20 century algebraic geometry, and was also the content of Hilbert's 15th problem "Rigorous foundation of Schubert's enumerative calculus" [33, 39, 50]. Thanks to the pioneer works [26,59] of Van der Waerden and Ehresmann, the problem of characteristics [53], considered by Schubert as the fundamental problem of enumerative geometry, has had a concise statement by the 1950's.

Let $G$ be compact connected Lie group with a maximal torus $T$. For an one parameter subgroup $\alpha: \mathbb{R} \rightarrow G$ its centralizer $P$ is a parabolic subgroup on $G$, while the homogeneous space $G / P$ is a projective variety, called a flag manifold of $G$. Let $W(P, G)$ be the set of left cosets of the Weyl group $W$ of $G$ by the Weyl group $W(P)$ of $P$ with associated length function $l: W(P, G) \rightarrow \mathbb{Z}$. The following result was discovered by Ehresmann [26, 1934] for the Grassmannians $G_{n, k}$ of $k$-dimensional linear subspaces on $\mathbb{C}^{n}$, announced by Chevalley $[13,1958]$ for the complete flag manifolds $G / T$, and extended to all flag manifolds $G / P$ by Bernstein-Gel'fandGel'fand [5, 1973].

Theorem 1.1 The flag manifold $G / P$ admits a decomposition into the cells indexed by the elements of $W(P, G)$,

$$
\begin{equation*}
G / P=\underset{w \in W(P, G)}{\cup} X_{w}, \quad \operatorname{dim} X_{w}=2 l(w), \tag{1.1}
\end{equation*}
$$

with each cell $X_{w}$ the closure of an algebraic affine space, called the Schubert variety on $G / P$ associated to $w$.

Since only even dimensional cells are involved in the decomposition (1.1), the set $\left\{\left[X_{w}\right], w \in W(P, G)\right\}$ of fundamental classes forms an additive basis of the integral homology $H_{*}(G / P)$. The cocycle classes $s_{w} \in H^{*}(G / P)$ Kronecker dual to the basis (i.e. $\left.\left\langle s_{w},\left[X_{u}\right]\right\rangle=\delta_{w, u}, w, u \in W(P, G)\right)$ gives rise to the Schubert class associated to $w \in W(P, G)$. Theorem 1.1 implies the following result, well known as the basis theorem of Schubert calculus [59, Sect. 8].
Theorem 1.2 The set $\left\{s_{w}, w \in W(P, G)\right\}$ of Schubert classes forms an additive basis of the integral cohomology $H^{*}(G / P)$.

An immediate consequence is that any monomial $s_{w_{1}} \cdots s_{w_{k}}$ in the Schubert classes on $G / P$ can be expressed as a linear combination of the basis elements

$$
\begin{equation*}
s_{w_{1}} \cdots s_{w_{k}}=\sum_{w \in W(P, G), l(w)=l\left(w_{1}\right)+\cdots+l\left(w_{k}\right),} a_{w_{1}, \ldots, w_{k}}^{w} \cdot s_{w}, a_{w_{1}, \ldots, w_{k}}^{w} \in \mathbb{Z}, \tag{1.2}
\end{equation*}
$$

[^7]where the coefficients $a_{w_{1}, \ldots, w_{k}}^{w}$ are called characteristics by Schubert [26, 53, 59].
The problem of characteristics. Given a monomial $s_{w_{1}} \cdots s_{w_{k}}$ in the Schubert classes, determine the characteristics $a_{w_{1}, \ldots, w_{k}}^{w}$ for all $w \in W(P, G)$ with $l(w)=l\left(w_{1}\right)+\cdots+l\left(w_{r}\right)$.

The characteristics are of particular importance in geometry, algebra and topology. They provide solutions to the problems of enumerative geometry [51-54]; were seen by Hilbert as "the degree of the final equations and the multiplicity of their solutions" of a system [33]; and are requested by describing the cohomology ring $H^{*}(G / P)$ in the Schubert basis [60, p. 331]. Notably, the degree of a Schubert variety [54] and the multiplicative rule of two Schubert classes [53] are two special cases of the problem which have received considerable attentions in literatures, see [15, 16] for accounts on the earlier relevant works.

This paper summaries and simplifies our series works [15, 16, 19-21] devoted to describe the integral cohomologies of flag manifolds by a minimal system of generators and relations in the Schubert classes. Precisely, based on a formula of the characteristics $a_{w_{1}, \ldots, w_{k}}^{w}$ stated in Sect. 2 and established in Sect. 3, we address in Sects. 4 and 5 a more direct approach to the Schubert presentations [20, 21] of the cohomology rings of flag manifolds $G / P$.

## 2 The Formula of the Characteristics $a_{w_{1}}^{w}, \ldots, w_{k}$

To investigate the topology of a flag manifold $G / P$ we may assume that the Lie group $G$ is 1-connected and simple. Resorting to the geometry of the Stiefel diagram of the Lie group $G$ we present in Theorem 2.4 a formula that boils down the characteristics $a_{w_{1}, \ldots, w_{k}}^{w}$ to the Cartan matrix of the group $G$.

Fix a maximal torus $T$ on $G$ and set $n=\operatorname{dim} T$. Equip the Lie algebra $L(G)$ with an inner product (, ) so that the adjoint representation acts as isometries of $L(G)$. The Cartan subalgebra of $G$ is the linear subspace $L(T) \subset L(G)$.

The restriction of the exponential map $\exp : L(G) \rightarrow G$ on $L(T)$ defines a set $S(G)$ of $\frac{1}{2}(\operatorname{dim} G-n)$ hyperplanes on $L(T)$, namely, the set of singular hyperplanes through the origin in $L(T)$ [11, p. 226]. These planes divide $L(T)$ into finitely many convex regions, each one is called a Weyl chamber of $G$. The reflections on $L(T)$ in these planes generate the Weyl group $W$ of $G$ [35, p. 49].

The map exp on $L(T)$ carries the normal line $l$ (through the origin) of a hyperplane $L \in S(G)$ to a circle subgroup on $T$. Let $\pm \alpha \in l$ be the non-zero vectors with minimal length so that $\exp ( \pm \alpha)=e$ (the group unit). The set $\Phi(G)$ consisting of all those vectors $\pm \alpha$ is called the root system of $G$. Fixing a regular point $x_{0} \in L(T)-$ $\cup_{L \in D(G)} L$ the set of simple roots relative to $x_{0}$ is

$$
\Delta\left(x_{0}\right)=\left\{\beta \in \Phi(G) \mid\left(\beta, x_{0}\right)>0\right\}[35, \text { p. 47]. }
$$

In addition, for a simple root $\beta \in \Delta$ the simple reflection relative to $\beta$ is the reflection $\sigma_{\beta}$ in the plane $L_{\beta} \in S(G)$ perpendicular to $\beta$. If $\beta, \beta^{\prime} \in \Delta$ the Cartan number

$$
\beta \circ \beta^{\prime}:=2\left(\beta, \beta^{\prime}\right) /\left(\beta^{\prime}, \beta^{\prime}\right)
$$

is always an integer, and only $0, \pm 1, \pm 2, \pm 3$ can occur [35, p. 55].
Since the set of simple reflections $\left\{\sigma_{\beta} \mid \beta \in \Delta\right\}$ generates $W$ [11, p. 193] every $w \in W$ admits a factorization of the form

$$
\begin{equation*}
w=\sigma_{\beta_{1}} \circ \cdots \circ \sigma_{\beta_{m}}, \beta_{i} \in \Delta \tag{2.1}
\end{equation*}
$$

Definition 2.1 The length $l(w)$ of an element $w \in W$ is the least number of factors in all decompositions of $w$ in the form (2.1). The decomposition (2.1) is called reduced if $m=l(w)$.

For a reduced decomposition (2.1) of $w$ the $m \times m$ (strictly upper triangular) matrix $A_{w}=\left(a_{i, j}\right)$ with $a_{i, j}=0$ if $i \geq j$ and $-\beta_{j} \circ \beta_{i}$ if $i<j$ is called the Cartan matrix of $w$ relative to the decomposition (2.1).

Example 2.2 In [57] Stembridge asked for an approach to find a reduced decomposition (2.1) for each $w \in W$. Resorting to the geometry of the Cartan subalgebra $L(T)$ this task can be implemented by the following method.

Picture $W$ as the $W$-orbit $\left\{w\left(x_{0}\right) \in L(T) \mid w \in W\right\}$ through the regular point $x_{0}$. For a $w \in W$ let $C_{w}$ be a line segment on $L(T)$ from the Weyl chamber containing $x_{0}$ to $w\left(x_{0}\right)$, that crosses the planes in $S(G)$ once at a time. Assume that they are met in the order $L_{\alpha_{1}}, \ldots, L_{\alpha_{k}}, \alpha_{i} \in \Phi(G)$. Then $l(w)=k$ and $w=\sigma_{\alpha_{k}} \circ \cdots \circ \sigma_{\alpha_{1}}$. Set

$$
\beta_{1}=\alpha_{1}, \quad \beta_{2}=\sigma_{\alpha_{1}}\left(\alpha_{2}\right), \ldots, \quad \beta_{k}=\sigma_{\alpha_{1}} \circ \cdots \circ \sigma_{\alpha_{k-1}}\left(\alpha_{k}\right)
$$

Then, from $\beta_{i} \in \Delta$ and $\sigma_{\beta_{i}}=\sigma_{\alpha_{1}} \circ \cdots \circ \sigma_{\alpha_{i-1}} \circ \sigma_{\alpha_{i}} \circ \sigma_{\alpha_{i-1}} \circ \cdots \circ \sigma_{\alpha_{1}}$ one sees that $w=\sigma_{\beta_{1}} \circ \cdots \circ \sigma_{\beta_{k}}$, which is reduced because of $l(w)=k$.

Let $\mathbb{Z}\left[x_{1}, \ldots, x_{m}\right]=\oplus_{n \geq 0} \mathbb{Z}\left[x_{1}, \ldots, x_{m}\right]^{(n)}$ be the ring of integral polynomials in $x_{1}, \ldots, x_{m}$, graded by $\operatorname{deg} x_{i}=1$.

Definition 2.3 For a $m \times m$ strictly upper triangular integer matrix $A=\left(a_{i, j}\right)$ the triangular operator $T_{A}$ associated to $A$ is the additive homomorphism $T_{A}$ : $\mathbb{Z}\left[x_{1}, \ldots, x_{m}\right]^{(m)} \rightarrow \mathbb{Z}$ defined recursively by the following elimination rules:
(i) If $m=1$ (consequently $A=(0)$ ) then $T_{A}\left(x_{1}\right)=1$;
(ii) If $h \in \mathbb{Z}\left[x_{1}, \ldots, x_{m-1}\right]^{(m)}$ then $T_{A}(h)=0$;
(iii) For any $h \in \mathbb{Z}\left[x_{1}, \ldots, x_{m-1}\right]^{(m-r)}$ with $r \geq 1$,

$$
T_{A}\left(h \cdot x_{m}^{r}\right)=T_{A^{\prime}}\left(h \cdot\left(a_{1, m} x_{1}+\cdots+\bar{a}_{m-1, m} x_{m-1}\right)^{r-1}\right)
$$

where $A^{\prime}$ is the $(m-1) \times(m-1)$ strictly upper triangular matrix obtained from $A$ by deleting both of the $m$ th column and row.

By additivity, $T_{A}$ is defined for every $h \in \mathbb{Z}\left[x_{1}, \ldots, x_{m}\right]^{(m)}$ using the unique expansion $h=\sum_{0 \leq r \leq m} h_{r} \cdot x_{m}^{r}$ with $h_{r} \in \mathbb{Z}\left[x_{1}, \ldots, x_{m-1}\right]^{(m-r)}$.

For a parabolic subgroup $P$ of $G$ the set $W(P, G)$ of left cosets of $W$ by $W(P)$ can be identified with the subset of $W$

$$
\left.W(P, G)=\left\{w \in W \mid l(w) \leq l\left(w w^{\prime}\right), w^{\prime} \in W(P)\right\} \text { (by }[5,5.1]\right),
$$

where $l$ is the length function on $W$. Assume that $w=\sigma_{\beta_{1}} \circ \cdots \circ \sigma_{\beta_{m}}, \beta_{i} \in \Delta$, is a reduced decomposition of an element $w \in W(P, G)$ with associated Cartan matrix $A_{w}=\left(a_{i, j}\right)_{m \times m}$. For a multi-index $I=\left\{i_{1}, \ldots, i_{k}\right\} \subseteq\{1, \ldots, m\}$ we put $|I|:=k$ and set

$$
\sigma_{I}:=\sigma_{\beta_{i_{1}}} \circ \cdots \circ \sigma_{\beta_{i_{k}}} \in W, x_{I}:=x_{i_{1}} \cdots x_{i_{k}} \in \mathbb{Z}\left[x_{1}, \ldots, x_{m}\right] .
$$

Our promised formula for the characteristics is:
Theorem 2.4 For every monomial $s_{w_{1}} \cdots s_{w_{k}}$ in the Schubert classes on $G / P$ with $l(w)=l\left(w_{1}\right)+\cdots+l\left(w_{k}\right)$ we have

$$
\begin{equation*}
a_{w_{1}, \ldots, w_{k}}^{w}=T_{A_{w}}\left(\prod_{i=1, \ldots, k}\left(\sum_{\sigma_{I}=w_{i},|I|=l\left(w_{i}\right), I \subseteq\{1, \ldots, m\}} x_{I}\right)\right) . \tag{2.2}
\end{equation*}
$$

Remark 2.5 Formula (2.2) reduces the characteristic $a_{w_{1}, \ldots, w_{k}}^{w}$ to a polynomial in the Cartan numbers of the group $G$, hence applies uniformly to all flag manifolds $G / P$.

If $k=2$ the characteristic $a_{w_{1}, w_{2}}^{w}$ is well known as a Littlewood-Richardson coefficient, and the formula (2.2) has been obtained by Duan in [16]. In [62] Willems generalizes the formula of $a_{w_{1}, w_{2}}^{w}$ to the more general context of flag varieties associated to Kac-Moody groups, and for the equivariant cohomologies. Recently, Bernstein and Richmond [10] obtained also a formula expressing $a_{w_{1}, w_{2}}^{w}$ in the Cartan numbers of $G$.

## 3 Proof of the Characteristics Formula (2.2)

In this paper the homologies and cohomologies are over the ring $\mathbb{Z}$ of integers. If $f: X \rightarrow Y$ is a continuous map between two topological spaces $X$ and $Y, f_{*}$ (resp. $f^{*}$ ) denote the homology (resp. cohomology) homomorphism induced by $f$. For an oriented closed manifold $M$ (resp. a connected projective variety) the notion $[M] \in H_{\operatorname{dim} M}(M)$ stands for the orientation class. In addition, the Kronecker pairing between cohomology and homology of a space $X$ is written as $\langle\rangle:, H^{*}(X) \times H_{*}(X) \rightarrow \mathbb{Z}$. The proof of Theorem 2.4 makes use of the celebrated $K$-cycles on $G / T$ constructed by Bott and Samelson in [8, 9]. We begin by recalling the construction of these cycles, as well as their basic properties developed in [8, 9 , 15, 16].

For a simple Lie group $G$ fix a regular point $x_{0} \in L(T)$ and let $\Delta$ be the set of simple roots relative to $x_{0}$. For a $\beta \in \Delta$ let $K_{\beta}$ be the centralizer of the subspace $\exp \left(L_{\beta}\right)$ on $G$, where $L_{\beta} \in S(G)$ is the plane perpendicular to $\beta$. Then $T \subset K_{\beta}$ and the quotient $K_{\beta} / T$ is diffeomorphic to the 2 -sphere [ $9, \mathrm{p} .996$ ].

The 2 -sphere $K_{\beta} / T$ carries a natural orientation specified as follows. The Cartan decomposition of the Lie algebra $L\left(K_{\beta}\right)$ relative to the maximal torus $T \subset K_{\beta}$ takes
the form $L\left(K_{\beta}\right)=L(T) \oplus \vartheta_{\beta}$, where $\vartheta_{\beta} \subset L(G)$ is the root space belonging to the root $\beta$ [35, p. 35]. Taking a non-zero vector $v \in \vartheta_{\beta}$ and letting $v^{\prime} \in \vartheta_{\beta}$ be such that $\left[v, v^{\prime}\right]=\beta$, where $\left[\right.$, ] is the Lie bracket on $L(G)$, then the ordered base $\left\{v, v^{\prime}\right\}$ furnishes $\vartheta_{\beta}$ with an orientation that is irrelevant to the choices of $v$. The tangent map of the quotient $\pi_{\beta}: K_{\beta} \rightarrow K_{\beta} / T$ at the group unit $e \in K_{\beta}$ maps the 2-plane $\vartheta_{\beta}$ isomorphically onto the tangent space to the sphere $K_{\beta} / T$ at the point $\pi_{\beta}(e)$. In this manner the orientation $\left\{v, v^{\prime}\right\}$ on $\vartheta_{\beta}$ furnishes the sphere $K_{\beta} / T$ with the orientation $\omega_{\beta}=\left\{\pi_{\beta}(v), \pi_{\beta}\left(v^{\prime}\right)\right\}$.

For an ordered sequence $\beta_{1}, \ldots, \beta_{m} \in \Delta$ of $m$ simple roots (repetitions like $\beta_{i}=$ $\beta_{j}$ may occur) let $K\left(\beta_{1}, \ldots, \beta_{m}\right)$ be the product group $K_{\beta_{1}} \times \cdots \times K_{\beta_{m}}$. With $T \subset$ $K_{\beta_{i}}$ the product $T \times \cdots \times T$ ( $m$-copies) acts on $K\left(\beta_{1}, \ldots, \beta_{m}\right)$ by

$$
\left(g_{1}, \ldots, g_{m}\right)\left(t_{1}, \ldots, t_{m}\right)=\left(g_{1} t_{1}, t_{1}^{-1} g_{2} t_{2}, \ldots, t_{m-1}^{-1} g_{m} t_{m}\right)
$$

Let $\Gamma\left(\beta_{1}, \ldots, \beta_{m}\right)$ be the base manifold of this principal action that is oriented by the $\omega_{\beta_{i}}, 1 \leq i \leq m$. The point on $\Gamma\left(\beta_{1}, \ldots, \beta_{m}\right)$ corresponding to the point $\left(g_{1}, \ldots, g_{m}\right) \in K\left(\beta_{1}, \ldots, \beta_{m}\right)$ is called $\left[g_{1}, \ldots, g_{m}\right]$.

The integral cohomology of the oriented manifold $\Gamma\left(\beta_{1}, \ldots, \beta_{m}\right)$ has been determined by Bott and Samelson in [8]. Let $\varphi_{i}: K_{\beta_{i}} / T \rightarrow \Gamma\left(\beta_{1}, \ldots, \beta_{m}\right)$ be the embedding induced by the inclusion $K_{\beta_{i}} \rightarrow K\left(\beta_{1}, \ldots, \beta_{m}\right)$ onto the $i$ th factor group, and put

$$
y_{i}=\varphi_{i_{*}}\left(\omega_{\beta_{i}}\right) \in H_{2}\left(\Gamma\left(\beta_{1}, \ldots, \beta_{m}\right)\right), 1 \leq i \leq m .
$$

Form the $m \times m$ strictly upper triangular matrix $A=\left(a_{i, j}\right)_{m \times m}$ by setting

$$
a_{i, j}=0 \text { if } i \geq j, \text { but } a_{i, j}=-2\left(\beta_{j}, \beta_{i}\right) /\left(\beta_{i}, \beta_{i}^{\prime}\right) \text { if } i<j
$$

It is easy to see from the construction that the set $\left\{y_{1}, \ldots, y_{m}\right\}$ forms a basis of the second homology group $H_{2}\left(\Gamma\left(\beta_{1}, \ldots, \beta_{m}\right)\right)$.
Lemma 3.1 ([9]) Let $x_{1}, \ldots, x_{m} \in H^{2}\left(\Gamma\left(\beta_{1}, \ldots, \beta_{m}\right)\right)$ be the Kronecker duals of the cycle classes $y_{1}, \ldots, y_{m}$ on $\Gamma\left(\beta_{1}, \ldots, \beta_{m}\right)$. Then

$$
\begin{equation*}
H^{*}\left(\Gamma\left(\beta_{1}, \ldots, \beta_{m}\right)\right)=\mathbb{Z}\left[x_{1}, \ldots, x_{m}\right] / J \tag{3.1}
\end{equation*}
$$

where $J$ is the ideal generated by $x_{j}^{2}-\sum_{i<j} a_{i, j} x_{i} x_{j}, 1 \leq j \leq m$.
In view of (3.1) the map $p_{\Gamma\left(\beta_{1}, \ldots, \beta_{m}\right)}$ from the polynomial ring $\mathbb{Z}\left[x_{1}, \ldots, x_{m}\right]$ onto its quotient $H^{*}\left(\Gamma\left(\beta_{1}, \ldots, \beta_{m}\right)\right)$ gives rise to the additive map

$$
\int_{\Gamma\left(\beta_{1}, \ldots, \beta_{m}\right)}: \mathbb{Z}\left[x_{1}, \ldots, x_{m}\right]^{(m)} \rightarrow \mathbb{Z}
$$

evaluated by $\int_{\Gamma\left(\beta_{1}, \ldots, \beta_{m}\right)} h=\left\langle p_{\Gamma\left(\beta_{1}, \ldots, \beta_{m}\right)}(h),\left[\Gamma\left(\beta_{1}, \ldots, \beta_{m}\right)\right]\right\rangle$. The geometric implication of the triangular operator $T_{A}$ in Definition 2.3 is shown by the following result.

Lemma 3.2 ([15, Proposition 2]) We have

$$
\int_{\Gamma\left(\beta_{1}, \ldots, \beta_{m}\right)}=T_{A}: \mathbb{Z}\left[x_{1}, \ldots, x_{m}\right]^{(m)} \rightarrow \mathbb{Z}
$$

In particular, $\int_{\Gamma\left(\beta_{1}, \ldots, \beta_{m}\right)} x_{1} \cdots x_{m}=1$.
For a parabolic $P$ on $G$ we can assume, without loss of the generalities, that $T \subseteq$ $P \subset G$. For a sequence $\beta_{1}, \ldots, \beta_{m}$ of simple roots the associated Bott-Samelson's $K$-cycle on $G / P$ is the map

$$
\varphi_{\beta_{1}, \ldots, \beta_{m} ; P}: \Gamma\left(\beta_{1}, \ldots, \beta_{m}\right) \rightarrow G / P
$$

defined by $\varphi_{\beta_{1}, \ldots, \beta_{m} ; P}\left(\left[g_{1}, \ldots, g_{m}\right]\right)=g_{1} \cdots g_{m} P$. If $P=T$ Hansen [32] has shown that certain $K$-cycles are desingularizations of the Schubert varieties on $G / T$. The following more general result allows one to translate the calculation with Schubert classes on $G / P$ to computing with monomials in the much simpler ring $H^{*}\left(\Gamma\left(\beta_{1}, \ldots, \beta_{m}\right)\right)$.

Lemma 3.3 With respect the Schubert basis on $H^{*}(G / P)$ the induced map of $\varphi_{\beta_{1}, \ldots, \beta_{m} ; P}$ on the cohomologies is given by

$$
\begin{equation*}
\varphi_{\beta_{1}, \ldots, \beta_{m} ; P}^{*}\left(s_{w}\right)=(-1)^{l(w)} \sum_{\sigma_{I}=w, I l \mid l(w), I \subseteq\lfloor 1, \ldots, m]} x_{I}, w \in W(P ; G) . \tag{3.2}
\end{equation*}
$$

Proof With $T \subseteq P \subset G$ the map $\varphi_{\beta_{1}, \ldots, \beta_{m} ; P}$ factors through $\varphi_{\beta_{1}, \ldots, \beta_{m} ; T}$ in the fashion

$$
\begin{aligned}
& \Gamma\left(\beta_{1}, \ldots, \beta_{m}\right) \xrightarrow{\varphi_{\beta_{1}}, \ldots \beta_{m} ; T} G / T, \\
& \varphi_{\beta_{1}, \ldots, \beta_{m} ; P} \quad G / P
\end{aligned}
$$

where the map $\pi$ is the fibration with fiber $P / T$. By [16, Lemma 5.1] formula (3.2) holds for the case $P=T$. According to [5, Sect. 5] the induced map $\pi^{*}$ : $H^{*}(G / P) \rightarrow H^{*}(G / T)$ is given by $\pi^{*}\left(s_{w}\right)=s_{w}, w \in W(P ; G)$, showing formula (3.2) for the general case $T \subset P$.

Proof of Theorem 2.4 For a monomial $s_{w_{1}} \cdots s_{w_{k}}$ in the Schubert classes of $G / P$ assume as in (1.2) that

$$
\begin{equation*}
s_{w_{1}} \cdots s_{w_{k}}=\sum_{w \in W(P ; G), l(w)=m} a_{w_{1}, \ldots, w_{k}}^{w} \cdot s_{w}, a_{w_{1}, \ldots, w_{k}}^{w} \in \mathbb{Z} \tag{3.3}
\end{equation*}
$$

where $m=l\left(w_{1}\right)+\cdots+l\left(w_{k}\right)$. For an element $w_{0} \in W(P ; G)$ with a reduced decomposition $w_{0}=\sigma_{\beta_{1}} \circ \cdots \circ \sigma_{\beta_{m}}, \beta_{i} \in \Delta$, let $A_{w_{0}}=\left(a_{i, j}\right)_{m \times m}$ be the relative Cartan matrix. Applying the ring map $\varphi_{\beta_{1}, \ldots, \beta_{m} ; P}^{*}$ to the Eq. (3.3) on $H^{*}(G / P)$ we obtain by (3.2) the equality on the group $H^{2 m}\left(\Gamma\left(\beta_{1}, \ldots, \beta_{m}\right)\right)$

$$
(-1)^{l\left(w_{1}\right)+\cdots+l\left(w_{k}\right)} \prod_{1 \leq i \leq k}\left(\sum_{\sigma_{I}=w_{i},|l|=l\left(w_{i}\right), I \subseteq\{1, \ldots, m\}} x_{I}\right)=(-1)^{m} a_{w_{1}, \ldots, w_{k}}^{w_{0}} \cdot x_{1} \cdots x_{m}
$$

Applying $\int_{\Gamma\left(\beta_{1}, \ldots, \beta_{m}\right)}$ to both sides we get by Lemma 3.2 that

$$
(-1)^{l\left(w_{1}\right)+\cdots+l\left(w_{k}\right)} \cdot T_{A_{w}}\left(\prod_{1 \leq i \leq k}\left(\sum_{\sigma_{I}=w_{i}, I l \mid=l\left(w_{i}\right), I \leq\{1, \ldots, m\}} x_{I}\right)\right)=(-1)^{m} \cdot a_{w_{1}, \ldots, w_{k}}^{w_{0}}
$$

This is identical to (2.2) because of $m=l\left(w_{1}\right)+\cdots+l\left(w_{k}\right)$.

## 4 The Cohomology of Flag Manifolds $\boldsymbol{G} / \boldsymbol{P}$

The classical Schubert calculus amounts to the determination of the intersection rings on Grassmann varieties and on the so called "flag manifolds" of projective geometry. -Weil [60, p. 331]

A classical problem of topology is to express the integral cohomology ring $H^{*}(G / H)$ of a homogeneous space $G / H$ by a minimal system of explicit generators and relations. The traditional approach due to H. Cartan, Borel, Baum, Toda utilize various spectral sequence techniques [2,3,37,58, 64], and the calculation encounters the same difficulties when applied to a Lie group $G$ with torsion elements in its integral cohomology, in particular, when $G$ is one of the five exceptional Lie groups [38, 58, 61].

However, if $P \subset G$ is parabolic, Schubert calculus makes the structure of the ring $H^{*}(G / P)$ appearing in a new light. Given a set $\left\{y_{1}, \ldots, y_{k}\right\}$ of $k$ elements let $\mathbb{Z}\left[y_{1}, \ldots, y_{k}\right]$ be the ring of polynomials in $y_{1}, \ldots, y_{k}$ with integer coefficients. For a subset $\left\{r_{1}, \ldots, r_{m}\right\} \subset \mathbb{Z}\left[y_{1}, \ldots, y_{k}\right]$ of homogenous polynomials denote by $\left\langle r_{1}, \ldots, r_{m}\right\rangle$ the ideal generated by $r_{1}, \ldots, r_{m}$.

Theorem 4.1 For each flag manifold $G / P$ there exist a set $\left\{y_{1}, \ldots, y_{k}\right\}$ of Schubert classes on $G / P$, and a set $\left\{r_{1}, \ldots, r_{m}\right\} \subset \mathbb{Z}\left[y_{1}, \ldots, y_{k}\right]$ of polynomials, so that the inclusion $\left\{y_{1}, \ldots, y_{k}\right\} \subset H^{*}(G / P)$ induces a ring isomorphism

$$
\begin{equation*}
H^{*}(G / P)=\mathbb{Z}\left[y_{1}, \ldots, y_{k}\right] /\left\langle r_{1}, \ldots, r_{m}\right\rangle, \tag{4.1}
\end{equation*}
$$

where both the numbers $k$ and $m$ are minimal subject to this presentation.
Proof Let $D\left(H^{*}(G / P)\right) \subset H^{*}(G / P)$ be the ideal of the decomposable elements. Since the ring $H^{*}(G / P)$ is torsion free and has a basis consisting of Schubert classes, there is a set $\left\{y_{1}, \ldots, y_{k}\right\}$ of Schubert classes on $G / P$ that corresponds to a basis of the quotient group $H^{*}(G / P) / D\left(H^{*}(G / P)\right)$. In particular, the inclusion $\left\{y_{1}, \ldots, y_{k}\right\} \subset$ $H^{*}(G / P)$ induces a surjective ring map

$$
f: \mathbb{Z}\left[y_{1}, \ldots, y_{k}\right] \rightarrow H^{*}(G / P) .
$$

Since ker $f$ is an ideal the Hilbert basis theorem implies that there exists a finite subset $\left\{r_{1}, \ldots, r_{m}\right\} \subset \mathbb{Z}\left[y_{1}, \ldots, y_{n}\right]$ so that $\operatorname{ker} f=\left\langle r_{1}, \ldots, r_{m}\right\rangle$. We can of course assume that the number $m$ is minimal subject to this constraint.

As the cardinality of a basis of the quotient group $H^{*}(G / P) / D\left(H^{*}(G / P)\right)$ the number $k$ is an invariant of $G / P$. In addition, if one changes the generators $y_{1}, \ldots, y_{k}$ to $y_{1}^{\prime}, \ldots, y_{k}^{\prime}$, then each old generator $y_{i}$ can be expressed as a polynomial $g_{i}$ in the new ones $y_{1}^{\prime}, \ldots, y_{k}^{\prime}$, and the invariance of the number $m$ is shown by the presentation

$$
H^{*}(G / P)=\mathbb{Z}\left[y_{1}^{\prime}, \ldots, y_{k}^{\prime}\right] /\left\langle r_{1}^{\prime}, \ldots, r_{m}^{\prime}\right\rangle,
$$

where $r_{j}^{\prime}$ is obtained from $r_{j}$ by substituting $g_{i}$ for $y_{i}, 1 \leq j \leq m$.
A presentation of the ring $H^{*}(G / P)$ in the form of (4.1) will be called a Schubert presentation of the cohomology of $G / P$, while the set $\left\{y_{1}, \ldots, y_{k}\right\}$ of generators will be called a set of special Schubert classes on $G / P$. Based on the characteristic formula (2.2) we develop in this section algebraic and computational machineries implementing Schubert presentation of the ring $H^{*}(G / P)$. To be precise the following conventions will be adopted throughout the remaining part of this section.
(i) $G$ is a 1 -connected simple Lie group with Weyl group $W$, and a fixed maximal torus $T$;
(ii) A set $\Delta=\left\{\beta_{1}, \ldots, \beta_{n}\right\}$ of simple roots of $G$ is given and ordered as the vertex of the Dykin diagram of $G$ pictured on [32, p. 58];
(iii) For each simple root $\beta_{i} \in \Delta$ write $\sigma_{i}$ instead of $\sigma_{\beta_{i}} \in W ; \omega_{i}$ in place of the Schubert class $s_{\sigma_{\beta_{i}}} \in H^{2}(G / T)$.

Note that Theorem 1.2 implies that the set $\left\{\omega_{1}, \ldots, \omega_{n}\right\}$ is the Schubert basis of the second cohomology $H^{2}(G / T)$, whose elements is identical to the fundamental dominant weights of $G$ in the context of Borel and Hirzebruch [4, 17].

### 4.1 Decomposition

By convention (iii) each $w \in W$ admits a factorization of the form

$$
w=\sigma_{i_{1}} \circ \cdots \circ \sigma_{i_{k}}, 1 \leq i_{1}, \ldots, i_{k} \leq n, l(w)=k .
$$

hence can be written as $w=\sigma[I]$ with $I=\left(i_{1}, \ldots, i_{k}\right)$. Such expressions of $w$ may not be unique, but the ambiguity can be dispelled by employing the following notion. Furnish the set of all reduced decompositions of $w$

$$
D(w):=\left\{w=\sigma[I] \mid I=\left(i_{1}, \ldots, i_{k}\right), l(w)=k\right\} .
$$

with the order $\leq$ given by the lexicographical order on the multi-indexes $I$. Call a decomposition $w=\sigma[I]$ minimized if $I$ is the minimal one with respect to the order. As result every $w \in W$ possesses a unique minimized decomposition.

For a subset $K \subset\{1, \ldots, n\}$ let $P_{K} \subset G$ be the centralizer of the 1-parameter subgroup $\{\exp (t b) \in G \mid t \in \mathbb{R}\}$ on $G$, where $b \in L(T)$ is a vector that satisfies

$$
\left(\beta_{i}, b\right)>0 \text { if } i \in K ;\left(\beta_{i}, b\right)=0 \text { if } i \notin K
$$

Then every parabolic subgroup $P$ is conjugate in $G$ to some $P_{K}$ with $K \subset\{1, \ldots, n\}$, while the Weyl group $W(P) \subset W$ is generated by the simple reflections $\sigma_{j}$ with $j \notin K$. Resorting to the length function $l$ on $W$ we embed the set $W(P ; G)$ as the subset of $W$ (as in Sect. 2)

$$
W(P ; G)=\left\{w \in W \mid l\left(w_{1}\right) \geq l(w), w_{1} \in w W(P)\right\}
$$

and put $W^{r}(P ; G):=\{w \in W(P ; G) \mid l(w)=r\}$. Since every $w \in W^{r}(P ; G)$ has a unique minimized decomposition as $w=\sigma[I]$, the set $W^{r}(P ; G)$ is also ordered by the lexicographical order on the multi-index $I$ 's, hence can be expressed as

$$
\begin{equation*}
W^{r}(P ; G)=\left\{w_{r, i} \mid 1 \leq i \leq \beta(r)\right\}, \beta(r):=\left|W^{r}(P ; G)\right| \tag{4.2}
\end{equation*}
$$

where $w_{r, i}$ is the $i$ th element in the ordered set $W^{r}(P ; G)$. In [19] a program entitled "Decomposition" is composed, whose function is summarized below.

Algorithm 4.2 Decomposition.
Input: The Cartan matrix $A=\left(a_{i j}\right)_{n \times n}$ of $G$, and a subset $K \subset\{1, \ldots, n\}$.
Output: The set $W\left(P_{K} ; G\right)$ being presented by the minimized decompositions of its elements, together with the index system (4.2) imposed by the order $\leq$.

For examples of the results coming from Decomposition we refer to [22, 1.1-7.1].

### 4.2 Factorization of the Ring $H^{*}(G / T)$ Using Fibration

The cardinality of the Schubert basis of $G / T$ agrees with the order of the Weyl group $W$, which in general is very large. To reduce the computation costs we may take a proper subset $K \subset\{1, \ldots, n\}$ and let $P:=P_{K}$ be the corresponding parabolic subgroup. The inclusion $T \subset P \subset G$ then induces the fibration

$$
\begin{equation*}
P / T \stackrel{i}{\hookrightarrow} G / T \xrightarrow{\pi} G / P, \tag{4.3}
\end{equation*}
$$

where the induced maps $\pi^{*}$ and $i^{*}$ behave well with respect to the Schubert bases of the three flag manifolds $P / T, G / P$ and $G / T$ in the following sense:
(i) With respect to the inclusion $W(P) \subset W$ the map $i^{*}$ carries the subset $\left\{s_{w}\right\}_{w \in W(P) \subset W}$ of the Schubert basis of $H^{*}(G / T)$ onto the Schubert basis of $H^{*}(P / T)$.
(ii) With respect to the inclusion $W(P ; G) \subset W$ the map $\pi^{*}$ identifies the Schubert basis $\left\{s_{w}\right\}_{w \in W(P ; G)}$ of $H^{*}(G / P)$ with a subset of the Schubert basis of $H^{*}(G / T)$.

For these reasons we can make no difference in notation between an element in $H^{*}(G / P)$ and its $\pi^{*}$ image in $H^{*}(G / T)$, and between a Schubert class on $P / T$ and its $i^{*}$ pre-image on $G / T$.

Assume now that $\left\{y_{1}, \ldots, y_{n_{1}}\right\}$ and $\left\{x_{1}, \ldots, x_{n_{2}}\right\}$ are respectively special Schubert classes on $P / T$ and $G / P$, and with respect to them one has the Schubert presentations

$$
\begin{equation*}
H^{*}(P / T)=\frac{\mathbb{Z}\left[y_{i}\right]_{1 \leq i \leq n_{1}}}{\left\langle h_{s}\right\rangle_{1 \leq s \leq m_{1}}} ; H^{*}(G / P)=\frac{\mathbb{Z}\left[x_{j}\right]_{1 \leq j \leq n_{2}}}{\left\langle r_{t}\right\rangle_{1 \leq t \leq m_{2}}}, \tag{4.4}
\end{equation*}
$$

where $h_{s} \in \mathbb{Z}\left[y_{i}\right]_{1 \leq i \leq n_{1}}, r_{t} \in \mathbb{Z}\left[x_{j}\right]_{1 \leq j \leq n_{2}}$. The following result allows one to formulate the ring $H^{*}(G / T)$ by the simpler ones $H^{*}(P / T)$ and $H^{*}(G / P)$.

Theorem 4.3 The inclusions $y_{i}, x_{j} \in H^{*}(G / T)$ induces a surjective ring map

$$
\varphi: \mathbb{Z}\left[y_{i}, x_{j}\right]_{1 \leq i \leq n_{1}, 1 \leq j \leq n_{2}} \rightarrow H^{*}(G / T) .
$$

Furthermore, if $\left\{\rho_{s}\right\}_{1 \leq s \leq m_{1}} \subset \mathbb{Z}\left[y_{i}, x_{j}\right]$ is a system satisfying

$$
\begin{equation*}
\rho_{s} \in \operatorname{ker} \varphi \text { and }\left.\rho_{s}\right|_{x_{j}=0}=h_{s}, \tag{4.5}
\end{equation*}
$$

then $\varphi$ induces a ring isomorphism

$$
\begin{equation*}
H^{*}(G / T)=\mathbb{Z}\left[y_{i}, x_{i}\right]_{1 \leq i \leq n_{1}, 1 \leq j \leq n_{2}} /\left\langle\rho_{s}, r_{t}\right\rangle_{1 \leq s \leq m_{1}, 1 \leq t \leq m_{2}} . \tag{4.6}
\end{equation*}
$$

Proof By the property (i) above the bundle (4.3) has the Leray-Hirsch property. That is, the cohomology $H^{*}(G / T)$ is a free module over the ring $H^{*}(G / P)$ with the basis $\left\{1, s_{w}\right\}_{w \in W(P)}$ :

$$
\begin{equation*}
H^{*}(G / T)=H^{*}(G / P)\left\{1, s_{w}\right\}_{w \in W(P)}([36, p .231]), \tag{4.7}
\end{equation*}
$$

implying that $\varphi$ surjects. It remains to show that for any $g \in \operatorname{ker} \varphi$ one has

$$
g \in\left\langle\rho_{s}, r_{t}\right\rangle_{1 \leq s \leq m_{1}, 1 \leq t \leq m_{2}} .
$$

To this end we notice by (4.5) and (4.7) that

$$
g \equiv \sum_{w \in W\left(P_{K}\right)} g_{w} \cdot s_{w} \bmod \left\langle\rho_{s}\right\rangle_{1 \leq s \leq m_{1}} \text { with } g_{w} \in \mathbb{Z}\left[x_{j}\right]_{1 \leq j \leq n_{2}}
$$

Thus $\varphi(g)=0$ implies $\varphi\left(g_{w}\right)=0$, showing $g_{w} \in\left\langle r_{t}\right\rangle_{1 \leq t \leq m_{2}}$ by (4.4).

### 4.3 The Generalized Grassmannians

For a topological space $X$ we set

$$
H^{\text {even }}(X):=\oplus_{r \geq 0} H^{2 r}(X), \quad H^{\text {odd }}(X):=\oplus_{r \geq 0} H^{2 r+1}(X)
$$

Then $H^{\text {even }}(X)$ is a subring of $H^{*}(X)$, while $H^{\text {odd }}(X)$ is a module over the ring $H^{\text {even }}(X)$.

If $P$ is a parabolic subgroup that corresponds to a singleton $K=\{i\}$, the flag manifold $G / P$ is called generalized Grassmannians of $G$ corresponding to the weight $\omega_{i}$ [20]. With $W^{1}(P ; G)=\left\{\sigma_{i}\right\}$ consisting of a single element the basis theorem implies that $H^{2}(G / P)=\mathbb{Z}$ is generated by $\omega_{i}$. Furthermore, letting $P^{s}$ be the semi-simple part of $P$, then the projection $p: G / P^{s} \rightarrow G / P$ is an oriented circle bundle on $G / P$ with Euler class $\omega_{i}$. With $H^{\text {odd }}(G / P)=0$ by the basis theorem the Gysin sequence [48, p. 143]

$$
\cdots \rightarrow H^{r}(G / P) \xrightarrow{p^{*}} H^{r}\left(G / P^{s}\right) \xrightarrow{\beta} H^{r-1}(G / P) \xrightarrow{\omega \cup} H^{r+1}(G / P) \xrightarrow{p^{*}} \cdots .
$$

of $p$ breaks into the short exact sequences

$$
\begin{equation*}
0 \rightarrow \omega_{i} \cup H^{2 r-2}(G / P) \rightarrow H^{2 r}(G / P) \xrightarrow{p^{*}} H^{2 r}\left(G / P^{s}\right) \rightarrow 0 \tag{4.8}
\end{equation*}
$$

as well as the isomorphisms

$$
\begin{equation*}
\beta: H^{2 r-1}\left(G / P^{s}\right) \cong \operatorname{ker}\left\{H^{2 r-2}(G / P) \xrightarrow{\omega_{i} \cup} H^{2 r}(G / P)\right\}, \tag{4.9}
\end{equation*}
$$

where $\omega_{i} \cup$ means taking cup product with $\omega_{i}$. In particular, formula (4.8) implies that

Lemma 4.4 If $S=\left\{y_{1}, \ldots, y_{m}\right\} \subset H^{*}(G / P)$ is a subset so that $p^{*} S=\left\{p^{*}\left(y_{1}\right), \ldots, p^{*}\left(y_{m}\right)\right\}$ is a minimal set of generators of the ring $H^{\text {even }}\left(G / P^{s}\right)$, then $S^{\prime}=\left\{\omega_{i}, y_{1}, \ldots, y_{m}\right\}$ is a minimal set of generators of $H^{*}(G / P)$.

By Lemma 4.4 the inclusions $\left\{\omega_{i}\right\} \cup S \subset H^{*}(G / P), p^{*} S \subset H^{*}\left(G / P^{s}\right)$ extend to the surjective maps $\pi$ and $\bar{\pi}$ that fit in the commutative diagram

$$
\begin{align*}
\mathbb{Z}\left[\omega_{i}, y_{1}, \ldots, y_{m}\right]^{(2 r)} & \xrightarrow{\varphi} \mathbb{Z}\left[y_{1}, \ldots, y_{m}\right]^{(2 r)} \\
\pi \downarrow & \bar{\pi} \downarrow  \tag{4.10}\\
H^{2 r-2}(G / P) \xrightarrow{\omega_{i} \cup} H^{2 r}(G / P) \quad & \xrightarrow{p^{*}} H^{2 r}\left(G / P^{s}\right) \quad \rightarrow 0
\end{align*}
$$

where $\mathbb{Z}\left[\omega_{i}, y_{1}, \ldots, y_{m}\right]$ is graded by $\operatorname{deg} \omega_{i}=2, \operatorname{deg} y_{i}$, and where

$$
\varphi\left(\omega_{i}\right)=0, \varphi\left(y_{i}\right)=y_{i} ; \bar{\pi}\left(y_{i}\right)=p^{*}\left(y_{i}\right) .
$$

The next result showing in [20, Lemma 8] enables us to formulate a presentation of the ring $H^{*}(G / P)$ in term of $H^{*}\left(G / P^{s}\right)$.

Theorem 4.5 Assume that $\left\{h_{1}, \ldots, h_{d}\right\} \subset \mathbb{Z}\left[y_{1}, \ldots, y_{m}\right]$ is a subset so that

$$
\begin{equation*}
H^{\text {even }}\left(G / P^{s}\right)=\mathbb{Z}\left[p^{*}\left(y_{1}\right), \ldots, p^{*}\left(y_{m}\right)\right] /\left\langle p^{*}\left(h_{1}\right), \ldots, p^{*}\left(h_{d}\right)\right\rangle, \tag{4.11}
\end{equation*}
$$

and that $\left\{d_{1}, \ldots, d_{t}\right\}$ is a basis of the module $H^{\text {odd }}\left(G / P^{s}\right)$ over $H^{\text {even }}\left(G / P^{s}\right)$. Then

$$
\begin{equation*}
H^{*}(G / P)=\mathbb{Z}\left[\omega_{i}, y_{1}, \ldots, y_{m}\right] /\left\langle r_{1}, \ldots, r_{d} ; \omega_{i} g_{1}, \ldots, \omega_{i} g_{t}\right\rangle \tag{4.12}
\end{equation*}
$$

where $\left\{r_{1}, \ldots, r_{d}\right\},\left\{g_{1}, \ldots, g_{t}\right\} \subset \mathbb{Z}\left[\omega_{i}, y_{1}, \ldots, y_{m}\right]$ are two sets of polynomials that satisfy respectively the following "initial constraints"
(i) $r_{k} \in \operatorname{ker} \pi$ with $\left.r_{k}\right|_{\omega_{i}=0}=h_{k}, \quad 1 \leq k \leq d$;
(ii) $\pi\left(g_{j}\right)=\beta\left(d_{j}\right), 1 \leq j \leq t$.

### 4.4 The Characteristics

To make Theorems 4.3 and 4.5 applicable in practical computation we develop in this section a series of three algorithms, entitled Characteristics, Null-space, Giambelli polynomials, all of them are based on the characteristic formula (2.2).

### 4.4.1 The Characteristics

For a $w \in W(P ; G)$ with the minimized decomposition $w=\sigma_{i_{1}} \circ \cdots \circ \sigma_{i_{m}}, 1 \leq$ $i_{1}, \ldots, i_{m} \leq n, l(w)=m$, we observe in formula (2.2) that
(i) The Cartan matrix $A_{w}$ of $w$ can be read directly from the Cartan matrix [35, p. 59] of the Lie group $G$;
(ii) For a $u \in W(P ; G)$ with $l(u)=r<m$ the solutions in the multi-indices $I=$ $\left\{j_{1}, \ldots, j_{r}\right\} \subseteq\left\{i_{1}, \ldots, i_{m}\right\}$ to the equation $\sigma_{I}=u$ in $W(P ; G)$ agree with the solutions to the linear system $\sigma_{I}\left(x_{0}\right)=u\left(x_{0}\right)$ on the vector space $L(T)$, where $x_{0}$ is the fixed regular point;
(iii) The evaluation the operator $T_{A_{w}}$ on a polynomial have been programmed using different methods in [19, 65].

Summarizing, granted with Decomposition, formula (2.2) indicates an effective algorithm to implement a parallel program whose function is briefed below.

Algorithm 4.6:Characteristics.
Input: The Cartan matrix $A=\left(a_{i j}\right)_{n \times n}$ of $G$, and a monomial $s_{w_{1}} \cdots s_{w_{k}}$ in Schubert classes on $G / P$.
Output: The expansion (1.2) of $s_{w_{1}} \cdots s_{w_{k}}$ in the Schubert basis.

### 4.4.2 The Null-Space

Let $\mathbb{Z}\left[y_{1}, \ldots, y_{k}\right]$ be the ring of polynomials in $y_{1}, \ldots, y_{k}$ graded by deg $y_{i}>0$, and let $\mathbb{Z}\left[y_{1}, \ldots, y_{k}\right]^{(m)}$ be the $\mathbb{Z}$-module consisting of all the homogeneous polynomials with degree $m$. Denote by $\mathbb{N}^{k}$ the set of all $k$-tuples $\alpha=\left(b_{1}, \ldots, b_{k}\right)$ of non-negative integers. Then the set of monomials basis of $\mathbb{Z}\left[y_{1}, \ldots, y_{k}\right]^{(m)}$ is

$$
\begin{equation*}
B(m)=\left\{y^{\alpha}=y_{1}^{b_{1}} \cdots y_{k}^{b_{k}} \mid \alpha=\left(b_{1}, \ldots, b_{k}\right) \in \mathbb{N}^{k}, \operatorname{deg} y^{\alpha}=m\right\} . \tag{4.13}
\end{equation*}
$$

It will be considered as an ordered set with respect to the lexicographical order on $\mathbb{N}^{k}$, whose cardinality is called $b(m)$.

Let $S=\left\{y_{1}, \ldots, y_{k}\right\}$ be a set of Schubert classes on $G / P$ that generates the ring $H^{*}(G / P)$ multiplicatively. Then the inclusion $S \subset H^{*}(G / P)$ induces a surjective ring map $f: \mathbb{Z}\left[y_{1}, \ldots, y_{k}\right] \rightarrow H^{*}(G / P)$ whose restriction to degree $2 m$ is

$$
f_{m}: \mathbb{Z}\left[y_{1}, \ldots, y_{k}\right]^{(2 m)} \rightarrow H^{2 m}(G / P) .
$$

Combining the Characteristics with the function "Null-space" in Mathematica, a basis of ker $f_{m}$ can be explicitly exhibited.

Let $s_{m, i}$ be the Schubert class corresponding to the element $w_{m, i} \in W(P ; G)$. With respect to the Schubert basis $\left\{s_{m, i} \mid 1 \leq i \leq \beta(m)\right\}$ on $H^{2 m}(G / P)$ every monomial $y^{\alpha} \in B(2 m)$ has the unique expansion

$$
\pi_{m}\left(y^{\alpha}\right)=c_{\alpha, 1} \cdot s_{m, 1}+\cdots+c_{\alpha, \beta(m)} \cdot s_{m, \beta(m)}, c_{\alpha, i} \in \mathbb{Z}
$$

where the coefficients $c_{\alpha, i}$ can be evaluated by the Characteristics. The matrix $M\left(f_{m}\right)=\left(c_{\alpha, i}\right)_{b(2 m) \times \beta(m)}$ so obtained is called the structure matrix of $f_{m}$. The builtin function Null-space in Mathematica transforms $M\left(f_{m}\right)$ to another matrix $N\left(f_{m}\right)$ in the fashion

```
In:=Null-space \(\left[M\left(f_{m}\right)\right]\)
Out:= a matrix \(N\left(f_{m}\right)=\left(b_{j, \alpha}\right)_{(b(2 m)-\beta(m)) \times b(2 m)}\),
```

whose significance is shown by the following fact.
Lemma 4.7 The set $\kappa_{i}=\sum_{y^{\alpha} \in B(2 m)} b_{i, \alpha} \cdot y^{\alpha}, 1 \leq i \leq b(2 m)-\beta(m)$, of polynomials is a basis of the $\mathbb{Z}$ module ker $f_{m}$.

### 4.4.3 The Giambelli polynomials (i.e. The Schubert Polynomials [6])

For the unitary group $G=U(n)$ of rank $n$ with parabolic subgroup $P=U(k) \times$ $U(n-k)$ the flag manifold $G_{n, k}=G / P$ is the Grassmannian of $k$-planes through the origin on $\mathbb{C}^{n}$. Let $1+c_{1}+\cdots+c_{k}$ be the total Chern class of the canonical $k$-bundle on $G_{n, k}$. Then the $c_{i}$ 's can be identified with appropriate Schubert classes
on $G_{n, k}$ (i.e. the special Schubert class on $G_{n, k}$ ), and one has the classical Schubert presentation

$$
H^{*}\left(G_{n, k}\right)=\mathbb{Z}\left[c_{1}, \ldots, c_{k}\right] /\left\langle r_{n-k+1}, \ldots, r_{n}\right\rangle,
$$

where $r_{j}$ is the component of the formal inverse of $1+c_{1}+\cdots+c_{k}$ in degree $j$. It follows that every Schubert class $s_{w}$ on $G_{n, k}$ can be written as a polynomial $\mathcal{G}_{w}\left(c_{1}, \ldots, c_{k}\right)$ in the special ones, and such an expression is afforded by the classical Giambelli formula [34, p. 112].

In general, assume that $G / P$ is a flag variety, and that a Schubert presentation (4.1) of the ring $H^{*}(G / P)$ has been specified. Then each Schubert class $s_{w}$ of $G / P$ can be expressed as a polynomial $\mathcal{G}_{w}\left(y_{1}, \ldots, y_{k}\right)$ in these special ones, and such an expression will be called a Giambelli polynomial of the class $s_{w}$. Based on the Characteristics a program implementing $\mathcal{G}_{w}\left(y_{1}, \ldots, y_{k}\right)$ has been compiled, whose function is summarized below.

```
Algorithm 4.8 : Giambelli polynomials
Input: A set {\mp@subsup{y}{1}{},\ldots,\mp@subsup{y}{k}{}}\mathrm{ of special Schubert classes on G/P.}
Output: Giambelli polynomials \mathcal{G}}(\mp@subsup{y}{1}{},\ldots,\mp@subsup{y}{k}{})\mathrm{ for all }w\inW(P;G)
```

We clarify the details in this program. By (4.13) we can write the ordered monomial basis $B(2 m)$ of $\mathbb{Z}\left[y_{1}, \ldots, y_{k}\right]^{(2 m)}$ as $\left\{y^{\alpha_{1}}, \ldots, y^{\alpha_{b(2 m)}}\right\}$. The corresponding structure matrix $M\left(f_{m}\right)$ in degree $2 m$ then satisfies

$$
\left(\begin{array}{c}
y^{\alpha_{1}} \\
\vdots \\
y^{\alpha_{b(2 m)}}
\end{array}\right)=M\left(f_{m}\right)\left(\begin{array}{c}
s_{m, 1} \\
\vdots \\
s_{m, \beta(m)}
\end{array}\right) .
$$

Since $f_{m}$ surjects the matrix $M\left(f_{m}\right)$ has a $\beta(m) \times \beta(m)$ minor equal to $\pm 1$. The standard integral row and column operation diagonalizing $M\left(f_{m}\right)$ [55, p. 162-164] provides us with two unique invertible matrices $P=P_{b(2 m) \times b(2 m)}$ and $Q=Q_{\beta(m) \times \beta(m)}$ that satisfy

$$
\begin{equation*}
P \cdot M\left(f_{m}\right) \cdot Q=\binom{I_{\beta(m)}}{C}_{b(2 m) \times \beta(m)}, \tag{4.14}
\end{equation*}
$$

where $I_{\beta(m)}$ is the identity matrix of $\operatorname{rank} \beta(m)$. The Giambelli polynomials is realized by the procedure below.

Step 1. Compute $M\left(f_{m}\right)$ using the Characteristics;
Step 2. Diagonalize $M\left(f_{m}\right)$ to get the matrices $P$ and $Q$;
Step 3. $\operatorname{Set}\left(\begin{array}{c}\mathcal{G}_{m, 1} \\ \vdots \\ \mathcal{G}_{m, \beta(m)}\end{array}\right)=Q \cdot[P]\left(\begin{array}{c}y^{\alpha_{1}} \\ \vdots \\ y^{\alpha_{b(2 m)}}\end{array}\right)$,
where $[P]$ is formed by the first $\beta(m)$ rows of $P$. Obviously, the polynomial $\mathcal{G}_{m, j}$ so obtained depends only on the special Schubert classes $\left\{y_{1}, \ldots, y_{k}\right\}$ on $G / P$, and is a Giambelli polynomial of $s_{m, j}, 1 \leq j \leq \beta(m)$.

## 5 Application to the Flag Manifolds $G / T$

A calculus, or science of calculation, is one which has organized processes by which passage is made, mechanically, from one result to another. -De Morgan.

Among all the flag manifolds $G / P$ associated to a Lie group $G$ it is the complete flag manifold $G / T$ that is of crucial importance, since the inclusion $T \subseteq P \subset G$ of subgroups induces the fibration

$$
P / T \hookrightarrow G / T \xrightarrow{\pi} G / P
$$

in which the induced map $\pi^{*}$ embeds the ring $H^{*}(G / P)$ as a subring of $H^{*}(G / T)$, see the proof of Theorem4.3. Further, according to E. Cartan [63, p. 674] all the 1-connected simple Lie groups consist of the three infinite families $\operatorname{SU}(n)$, $\operatorname{Sp}(n), \operatorname{Spin}(n)$ of the classical groups, as well as the five exceptional ones: $G_{2}, F_{4}, E_{6}, E_{7}, E_{8}$, while for any compact connected Lie group $G$ with a maximal torus $T$ one has a diffeomorphism

$$
G / T=G_{1} / T_{1} \times \cdots \times G_{k} / T_{k}
$$

with each $G_{i}$ an 1-connected simple Lie group and $T_{i} \subset G_{i}$ a maximal torus. Thus, the problem of finding Schubert presentations of flag manifolds may be reduced to the special cases $G / T$ where $G$ is 1-connected and simple.

In this section we determine the Schubert presentation of the ring $H^{*}(G / T)$ in accordance to $G$ is classical or exceptional. Recall that for a simple Lie group $G$ with rank $n$ the fundamental dominant weights $\left\{\omega_{1}, \ldots, \omega_{n}\right\}$ of $G[4]$ is precisely the Schubert basis of $H^{2}(G / T)$ [17, Lemma 2.4].

### 5.1 The Ring $H^{*}(G / T)$ Classical Lie Group

If $G=S U(n+1)$ or $S p(n)$ Borel [3] has shown that

$$
H^{*}(G / T)=\mathbb{Z}\left[\omega_{1}, \ldots, \omega_{n}\right] /\left\langle\mathbb{Z}\left[\omega_{1}, \ldots, \omega_{n}\right]^{+, W}\right\rangle
$$

where $\mathbb{Z}\left[\omega_{1}, \ldots, \omega_{n}\right]^{+, W}$ is the ring of the integral Weyl invariants of $G$ in positive degrees. It follows that if we let $c_{r}(G) \in H^{2 r}(G / T)$ be respectively the $r$ th element symmetric polynomial in the sets

$$
\left\{\omega_{1}, \omega_{k}-\omega_{k-1},-\omega_{n} \mid 2 \leq k \leq n\right\} \text { or }\left\{ \pm \omega_{1}, \pm\left(\omega_{k}-\omega_{k-1}\right) \mid 2 \leq k \leq n\right\},
$$

then we have
Theorem 5.1 For $G=S U(n)$ or $S p(n)$ Schubert presentation of $H^{*}(G / T)$ is

$$
\begin{align*}
H^{*}(S U(n) / T) & =\mathbb{Z}\left[\omega_{1}, \ldots, \omega_{n-1}\right] /\left\langle c_{2}, \ldots, c_{n}\right\rangle, c_{r}=c_{r}(S U(n))  \tag{5.1}\\
H^{*}(S p(n) / T) & =\mathbb{Z}\left[\omega_{1}, \ldots, \omega_{n}\right] /\left\langle c_{2}, \ldots, c_{2 n}\right\rangle, c_{2 r}=c_{2 r}(\operatorname{Sp}(n)) \tag{5.2}
\end{align*}
$$

Turning to the group $G=\operatorname{Spin}(2 n)$ let $y_{k}$ be the Schubert class on $\operatorname{Spin}(2 n) / T$ associated to the Weyl group element

$$
w_{k}=\sigma[n-k, \ldots, n-2, n-1], 2 \leq k \leq n-1
$$

(in the notation of Sect.4.1). According to Marlin [46, Proposition 3]

$$
\begin{equation*}
H^{*}(\operatorname{Spin}(2 n) / T)=\mathbb{Z}\left[\omega_{1}, \ldots, \omega_{n}, y_{2}, \ldots, y_{n-1}\right] /\left\langle\delta_{i}, \xi_{j}, \mu_{k}\right\rangle \tag{5.3}
\end{equation*}
$$

where

$$
\begin{aligned}
& \delta_{i}:=2 y_{i}-c_{i}\left(\omega_{1}, \ldots, \omega_{n}\right), 1 \leq i \leq n-1, \\
& \xi_{j}:=y_{2 j}+(-1)^{j} y_{j}^{2}+2 \sum_{1 \leq r \leq j-1}(-1)^{r} y_{r} y_{2 j-r}, 1 \leq j \leq\left[\frac{n-1}{2}\right], \\
& \mu_{k}:=(-1)^{k} y_{k}^{2}+2 \sum_{2 k-n+1 \leq r \leq k-1}(-1)^{r} y_{r} y_{2 k-r},\left[\frac{n}{2}\right] \leq k \leq n-1,
\end{aligned}
$$

and where $c_{i}\left(\omega_{1}, \ldots, \omega_{n}\right)$ is the $i$ th elementary symmetric function on set

$$
\left\{\omega_{n}, \omega_{i}-\omega_{i-1}, \omega_{n-1}+\omega_{n}-\omega_{n-2}, \omega_{n-1}-\omega_{n} \mid 2 \leq i \leq n-2\right\} .
$$

Since each Schubert class $y_{2 j}$ with $1 \leq j \leq\left[\frac{n-1}{2}\right]$ can be expressed as a polynomial in the $y_{2 i+1}$ 's by the relations of the type $\xi_{k}$, we obtain that
Theorem 5.2 The Schubert presentation of $H^{*}(\operatorname{Spin}(2 n) / T)$ is

$$
\begin{equation*}
H^{*}(\operatorname{Spin}(2 n) / T)=\mathbb{Z}\left[\omega_{1}, \ldots, \omega_{n}, y_{3}, y_{5}, \ldots, y_{2\left[\frac{n-1}{2}\right]-1}\right] /\left\langle r_{i}, h_{k}\right\rangle \tag{5.4}
\end{equation*}
$$

where $r_{i}$ and $h_{k}$ are the polynomials obtained respectively from $\delta_{i}$ and $\mu_{k}$ by replacing the classes $y_{2 r}$ with the polynomials (by the relation $\xi_{r}$ )

$$
(-1)^{r-1} y_{r}^{2}+2 \sum_{1 \leq k \leq r-1}(-1)^{k-1} y_{k} y_{2 r-k} .
$$

Similarly, if $G=\operatorname{Spin}(2 n+1)$ one can deduce a Schubert presentation of the ring $H^{*}(G / T)$ from Marlin's formula [46, Proposition 2].

Remark 5.3 For the classical Lie groups $G$ the Giambelli polynomials (i.e. Schubert polynomials) of the Schubert classes on $G / T$ have been determined by Billey and Haiman [6].

In comparison with Marlin's formula (5.3) the presentation (5.4) is more concise for involving fewer generators and relations. A basic requirement of topology is to present the cohomology of a space $X$ by a minimal system of generators and relations, so that the (rational) minimal model and $\kappa$-invariants of the Postnikov tower of $X$ [30] can be formulated accordingly.

### 5.2 The Ring $H^{*}(G / T)$ for an Exceptional Lie Group

Having clarified the Schubert presentation of the ring $H^{*}(G / T)$ for the classical $G$ we proceed to the exceptional cases $G=F_{4}, E_{6}$ or $E_{7}$ (the result for the case $E_{8}$ comes from the same calculation, only the presentation [21, Theorem 5.1] is slightly lengthy). In these cases the dimension $s=\operatorname{dim} G / T$ and the number $t$ of the Schubert classes on $G / T$ are

$$
(s, t)=(48,1152),(72,51840) \text { or }(126,2903040),
$$

respectively. Instead of describing the ring $H^{*}(G / T)$ using the totality of $t^{3}$ Littlewood-Richardson coefficients $c_{u, v}^{w}$ (with $c_{u, v}^{w}=0$ for $l(w) \neq l(u)+l(v)$ being understood) the idea of Schubert presentation brings us the following concise and explicit formulae of the ring $H^{*}(G / T)$.

Theorem 5.4 For $G=F_{4}, E_{6}$ and $E_{7}$ the Schubert presentations of the cohomologies $H^{*}(G / T)$ are

$$
\begin{align*}
& H^{*}\left(F_{4} / T\right)=\mathbb{Z}\left[\omega_{1}, \omega_{2}, \omega_{3}, \omega_{4}, y_{3}, y_{4}\right] /\left\langle\rho_{2}, \rho_{4}, r_{3}, r_{6}, r_{8}, r_{12}\right\rangle, \text { where }  \tag{5.5}\\
& \rho_{2}=c_{2}-4 \omega_{1}^{2} ; \\
& \rho_{4}=3 y_{4}+2 \omega_{1} y_{3}-c_{4} ; \\
& r_{3}=2 y_{3}-\omega_{1}^{3} ; \\
& r_{6}=y_{3}^{2}+2 c_{6}-3 \omega_{1}^{2} y_{4} ; \\
& r_{8}=3 y_{4}^{2}-\omega_{1}^{2} c_{6} ; \\
& r_{12}=y_{4}^{3}-c_{6}^{2} \\
& H^{*}\left(E_{6} / T\right)=\mathbb{Z}\left[\omega_{1}, \ldots, \omega_{6}, y_{3}, y_{4}\right] /\left\langle\rho_{2}, \rho_{3}, \rho_{4}, \rho_{5}, r_{6}, r_{8}, r_{9}, r_{12}\right\rangle, \text { where }  \tag{5.6}\\
& \rho_{2}=4 \omega_{2}^{2}-c_{2} ; \\
& \rho_{3}=2 y_{3}+2 \omega_{2}^{3}-c_{3} ;
\end{align*}
$$

```
\(\rho_{4}=3 y_{4}+\omega_{2}^{4}-c_{4}\);
\(\rho_{5}=2 \omega_{2}^{2} y_{3}-\omega_{2} c_{4}+c_{5}\);
\(r_{6}=y_{3}^{2}-\omega_{2} c_{5}+2 c_{6}\);
\(r_{8}=3 y_{4}^{2}-2 c_{5} y_{3}-\omega_{2}^{2} c_{6}+\omega_{2}^{3} c_{5}\);
\(r_{9}=2 y_{3} c_{6}-\omega_{2}^{3} c_{6}\);
\(r_{12}=y_{4}^{3}-c_{6}^{2}\).
```

$H^{*}\left(E_{7} / T\right)=\mathbb{Z}\left[\omega_{1}, \ldots, \omega_{7}, y_{3}, y_{4}, y_{5}, y_{9}\right] /\left\langle\rho_{i}, r_{j}\right\rangle$, where
$\rho_{2}=4 \omega_{2}^{2}-c_{2}$;
$\rho_{3}=2 y_{3}+2 \omega_{2}^{3}-c_{3}$;
$\rho_{4}=3 y_{4}+\omega_{2}^{4}-c_{4}$;
$\rho_{5}=2 y_{5}-2 \omega_{2}^{2} y_{3}+\omega_{2} c_{4}-c_{5}$;
$r_{6}=y_{3}^{2}-\omega_{2} c_{5}+2 c_{6}$;
$r_{8}=3 y_{4}^{2}+2 y_{3} y_{5}-2 y_{3} c_{5}+2 \omega_{2} c_{7}-\omega_{2}^{2} c_{6}+\omega_{2}^{3} c_{5}$;
$r_{9}=2 y_{9}+2 y_{4} y_{5}-2 y_{3} c_{6}-\omega_{2}^{2} c_{7}+\omega_{2}^{3} c_{6}$;
$r_{10}=y_{5}^{2}-2 y_{3} c_{7}+\omega_{2}^{3} c_{7}$;
$r_{12}=y_{4}^{3}-4 y_{5} c_{7}-c_{6}^{2}-2 y_{3} y_{9}-2 y_{3} y_{4} y_{5}+2 \omega_{2} y_{5} c_{6}+3 \omega_{2} y_{4} c_{7}+c_{5} c_{7}$;
$r_{14}=c_{7}^{2}-2 y_{5} y_{9}+2 y_{3} y_{4} c_{7}-\omega_{2}^{3} y_{4} c_{7}$;
$r_{18}=y_{9}^{2}+2 y_{5} c_{6} c_{7}-y_{4} c_{7}^{2}-2 y_{4} y_{5} y_{9}+2 y_{3} y_{5}^{3}-5 \omega_{2} y_{5}^{2} c_{7}$,
where the $c_{r}$ 's are the polynomials $c_{r}(P)$ in the weights $\omega_{1}, \ldots, \omega_{n}$ defined in (5.17); and where the $y_{i}$ 's are the Schubert classes on $G / T$ associated to the Weyl group elements tabulated below:

| $y_{i}$ | $F_{4} / T\left(F / P_{\{1\}}\right)$ | $E_{6} / T\left(E_{6} / P_{\{2\}}\right)$ | $E_{7} / T\left(E_{7} / P_{\{2\}}\right)$ |
| :--- | :--- | :--- | :--- |
| $y_{3}$ | $\sigma_{[3,2,1]}$ | $\sigma_{[5,4,2]}$ | $\sigma_{[5,4,2]}$ |
| $y_{4}$ | $\sigma_{[4,3,2,1]}$ | $\sigma_{[6,5,4,2]}$ | $\sigma_{[6,5,4,2]}$ |
| $y_{5}$ |  |  | $\sigma_{[7,6,5,4,2]}$ |
| $y_{6}$ | $\sigma_{[3,2,4,3,2,1]}$ | $\sigma_{[1,3,6,5,4,2]}$ | $\sigma_{[1,3,6,5,4,2]}$ |
| $y_{7}$ |  |  | $\sigma_{[1,3,7,6,5,4,2]}$ |
| $y_{9}$ |  |  | $\sigma_{[1,5,4,3,7,6,5,4,2]}$. |

To reduce the computational complexity of deriving Theorem 5.4 we choose for each $G=F_{4}, E_{6}$ or $E_{7}$ a parabolic subgroup $P$ associated to a singleton $K=\{i\}$, where the index $i$, as well as the isomorphism types of $P$ and its simple part $P^{s}$, is stated in the table below

| $G$ | $F_{4}$ | $E_{6}$ | $E_{7}$ |
| :--- | :--- | :--- | :--- |
| $i$ | 1 | 2 | 2 |
| $P ; P^{s}$ | $S p(3) \cdot S^{1} ; S p(3)$ | $S U(6) \cdot S^{1} ; S U(6)$ | $S U(7) \cdot S^{1} ; S U(7)$ |

In view of the circle bundle associated to $P$ (see in Sect.4.3)

$$
\begin{equation*}
S^{1} \hookrightarrow G / P^{s} \xrightarrow{p} G / P \tag{5.10}
\end{equation*}
$$

the calculation will be divided into three steps, in accordance to cohomologies of the three homogeneous spaces $G / P^{s}, G / P$ and $G / T$.

Step 1. The cohomologies $H^{*}\left(G / P^{s}\right)$. By the formulae (4.8) and (4.9) the additive structure of $H^{*}\left(G / P^{s}\right)$ is determined by the homomorphisms

$$
H^{2 r-2}(G / P) \xrightarrow{\cup \omega} H^{2 r}(G / P) .
$$

Explicitly, with respect to the Schubert basis $\left\{s_{r, 1}, \ldots, s_{r, \beta(r)}\right\}$ of $H^{2 r}(G / P), \beta(r)=$ $\left|W^{r}(P ; G)\right|$, the expansions

$$
\omega \cup s_{r-1, i}=\sum a_{i, j} \cdot s_{r, j}
$$

give rise to a $\beta(r-1) \times \beta(r)$ matrix $A_{r}$ that satisfies the linear system

$$
\left(\begin{array}{l}
\omega \cup s_{r-1,1} \\
\omega \cup s_{r-1,2} \\
\vdots \\
\omega \cup s_{r-1, \beta(r-1)}
\end{array}\right)=A_{r}\left(\begin{array}{l}
s_{r, 1} \\
s_{r, 2} \\
\vdots \\
s_{r, \beta(r)}
\end{array}\right)
$$

Since $\omega \cup s_{r-1, i}$ is a monomial in Schubert classes, the Characteristics is applicable to evaluate the entries of $A_{r}$. Diagonalizing $A_{r}$ by the integral row and column reductions [50, p. 162-166] one obtains the non-trivial groups $H^{r}\left(G / P^{s}\right)$, together with their basis, as that tabulated below, where
(1) $\bar{y}_{i}:=p^{*}\left(y_{i}\right)$ with $y_{i}$ the Schubert classes in table (5.8);
(2) For simplicity the non-trivial groups $H^{r}\left(G / P^{s}\right)$ are printed only up to the stage where all the generators and relations of the ring $H^{\text {even }}\left(G / P^{s}\right)$ emerge.

Step 2. The cohomologies $H^{*}(G / P)$. Summarizing the contents of Table 1 we find that

$$
H^{\text {even }}\left(F_{4} / P^{s}\right)=\mathbb{Z}\left[\bar{y}_{3}, \bar{y}_{4}, \bar{y}_{6}\right] /\left\langle p^{*}\left(h_{3}\right), p^{*}\left(h_{6}\right), p^{*}\left(h_{8}\right), p^{*}\left(h_{12}\right)\right\rangle,
$$

where

$$
h_{3}=2 y_{3}, h_{6}=2 y_{6}+y_{3}^{2}, h_{8}=3 y_{4}^{2}, h_{12}=y_{6}^{2}-y_{4}^{3},
$$

and that $H^{\text {odd }}\left(F_{4} / P^{s}\right)$ has the $H^{\text {even }}\left(F_{4} / P^{s}\right)$-module basis $\left\{d_{23}\right\}$. By Theorem 4.3 we obtain the partial presentation

$$
H^{*}\left(F_{4} / P\right)=\mathbb{Z}\left[\omega_{1}, y_{3}, y_{4}, y_{6}\right] /\left\langle r_{3}, r_{6}, r_{8}, r_{12}, r_{12}^{\prime}\right\rangle,
$$

Table 1 Non-trivial cohomologies of $F_{4} / P^{s}$

| Nontrivial $H^{k}$ | Basis elements | Relations |
| :--- | :--- | :--- |
| $H^{6} \cong \mathbb{Z}_{2}$ | $\bar{s}_{3,1}\left(=\bar{y}_{3}\right)$ | $2 \bar{y}_{3}=0$ |
| $H^{8} \cong \mathbb{Z}$ | $\bar{s}_{4,2}\left(=\bar{y}_{4}\right)$ |  |
| $H^{12} \cong \mathbb{Z}_{4}$ | $\bar{s}_{6,2}\left(=\bar{y}_{6}\right)$ | $-2 \bar{y}_{6}=\bar{y}_{3}^{2}$ |
| $H^{14} \cong \mathbb{Z}_{2}$ | $\bar{y}_{3} \bar{y}_{4}$ |  |
| $H^{16} \cong \mathbb{Z}_{3}$ | $\bar{y}_{4}^{2}$ | $3 \bar{y}_{4}^{2}=0$ |
| $H^{18} \cong \mathbb{Z}_{2}$ | $\bar{y}_{3} \bar{y}_{6}$ |  |
| $H^{20} \cong \mathbb{Z}_{4}$ | $\bar{y}_{4} \bar{y}_{6}$ |  |
| $H^{26} \cong \mathbb{Z}_{2}$ | $\bar{y}_{3} \bar{y}_{4} \bar{y}_{6}$ |  |
| $H^{23} \cong \mathbb{Z}$ | $d_{23}=\beta^{-1}\left(2 s_{11,1}-s_{11,2}\right)$ |  |

Table 2 Non-trivial cohomologies of $E_{6} / P^{s}$

| Nontrivial $H^{k}$ | Basis elements | Relations |
| :--- | :--- | :--- |
| $H^{6} \cong \mathbb{Z}$ | $\bar{s}_{3,2}\left(=\bar{y}_{3}\right)$ |  |
| $H^{8} \cong \mathbb{Z}$ | $\bar{s}_{4,3}\left(=\bar{y}_{4}\right)$ |  |
| $H^{12} \cong \mathbb{Z}$ | $\bar{s}_{6,1}\left(=\bar{y}_{6}\right)$ | $-2 \bar{y}_{6}=\bar{y}_{3}^{2}$ |
| $H^{14} \cong \mathbb{Z}$ | $\bar{y}_{3} \bar{y}_{4}$ |  |
| $H^{16} \cong \mathbb{Z}_{3}$ | $\bar{y}_{4}^{2}$ | $3 \bar{y}_{4}^{2}=0$ |
| $H^{18} \cong \mathbb{Z}_{2}$ | $\bar{y}_{3} \bar{y}_{6}$ | $2 \bar{y}_{3} \bar{y}_{6}=0$ |
| $H^{20} \cong \mathbb{Z}$ | $\bar{y}_{4} \bar{y}_{6}$ |  |
| $H^{22} \cong \mathbb{Z} 3$ | $\bar{y}_{3} \bar{y}_{4}^{2}$ |  |
| $H^{26} \cong \mathbb{Z}_{2}$ | $\bar{y}_{3} \bar{y}_{4} \bar{y}_{6}$ |  |
| $H^{28} \cong \mathbb{Z}_{3}$ | $\bar{y}_{4}^{2} \bar{y}_{6}$ |  |
| $H^{23} \cong \mathbb{Z}$ | $d_{23}=\beta^{-1}\left(s_{11,1}-s_{11,2}-s_{11,3}+s_{11,4}-s_{11,5}+s_{11,6}\right)$ |  |
| $H^{29} \cong \mathbb{Z}$ | $d_{29}=\beta^{-1}\left(-s_{14,1}+s_{14,2}+s_{14,4}-s_{14,5}\right)$ | $2 d_{29}= \pm \bar{y}_{3} d_{23}$ |

indicating that the inclusion $\left\{\omega_{1}, y_{3}, y_{4}, y_{6}\right\} \subset H^{*}\left(F_{4} / P\right)$ induces the surjective ring $\operatorname{map} f: \mathbb{Z}\left[\omega_{1}, y_{3}, y_{4}, y_{6}\right] \rightarrow H^{*}\left(F_{4} / P\right)$. Further, according to Lemma4.7, computing with the Null-space $N\left(f_{m}\right)$ in the order $m=3,6,8$ and 12 suffices to decide the generators of the ideal ker $f$ to yields the Schubert presentation

$$
\begin{equation*}
H^{*}\left(F_{4} / P\right)=\mathbb{Z}\left[\omega_{1}, y_{3}, y_{4}, y_{6}\right] /\left\langle r_{3}, r_{6}, r_{8}, r_{12}\right\rangle \text {, where } \tag{5.11}
\end{equation*}
$$

$$
\begin{aligned}
& r_{3}=2 y_{3}-\omega_{1}^{3} ; \\
& r_{6}=2 y_{6}+y_{3}^{2}-3 \omega_{1}^{2} y_{4} ; \\
& r_{8}=3 y_{4}^{2}-\omega_{1}^{2} y_{6} ; \\
& r_{12}=y_{6}^{2}-y_{4}^{3} .
\end{aligned}
$$

Similarly, combining the Null-space with the contents of Tables 2 and 3 one gets the Schubert presentations of the ring $H^{*}(G / P)$ for $G=E_{6}$ and $E_{7}$ as

Table 3 Non-trivial cohomologies of $E_{7} / P^{s}$

| Nontrivial $H^{k}$ | Basis elements |
| :---: | :---: |
| $H^{6} \cong \mathbb{Z}$ | $s_{3,2}=\bar{y}_{3}$ |
| $H^{8} \cong \mathbb{Z}$ | $s_{4,3}=\bar{y}_{4}$ |
| $H^{10} \cong \mathbb{Z}$ | $s_{5,4}=\bar{y}_{5}$ |
| $H^{12} \cong \mathbb{Z}$ | $s_{6,5}=\bar{y}_{3}^{2}+\bar{y}_{6}$ |
| $H^{14} \cong \mathbb{Z} \oplus \mathbb{Z}$ | $s_{7,6}=-\bar{y}_{7} ; \bar{y}_{3} \bar{y}_{4}$ |
| $H^{16} \cong \mathbb{Z}$ | $\bar{y}_{3} \bar{y}_{5}-2 \bar{y}_{4}^{2}$ |
| $H^{18} \cong \mathbb{Z}_{2} \oplus \mathbb{Z} \oplus \mathbb{Z}$ | $\bar{y}_{3}^{3}+\bar{y}_{3} \bar{y}_{6}+\bar{y}_{4} \bar{y}_{5}+\bar{y}_{9} ;-\bar{y}_{4} \bar{y}_{5} ; \bar{y}_{3}^{3}+\bar{y}_{9}$ |
| $H^{20} \cong \mathbb{Z} \oplus \mathbb{Z}$ | $\bar{y}_{3} \bar{y}_{7}-\bar{y}_{5}^{2} ; \bar{y}_{3} \bar{y}_{7}+\bar{y}_{4} \bar{y}_{6}$ |
| $H^{22} \cong \mathbb{Z} \oplus \mathbb{Z}$ | $\bar{y}_{3}^{2} \bar{y}_{5}-\bar{y}_{3} \bar{y}_{4}^{2}+\bar{y}_{4} \bar{y}_{7} ;-2 \bar{y}_{4} \bar{y}_{7}+\bar{y}_{5} \bar{y}_{6}$ |
| $H^{24} \cong \mathbb{Z}_{2} \oplus \mathbb{Z} \oplus \mathbb{Z}$ | $\begin{aligned} & \bar{y}_{3}^{4}+\bar{y}_{3}^{2} \bar{y}_{6}+\bar{y}_{3} \bar{y}_{4} \bar{y}_{5}+\bar{y}_{3} \bar{y}_{9} \\ & \bar{y}_{3}^{2} \bar{y}_{6}+\bar{y}_{3} \bar{y}_{4} \bar{y}_{5}-\bar{y}_{4}^{3} ;-\bar{y}_{4}^{3}+\bar{y}_{5} \bar{y}_{7}+\bar{y}_{6}^{2} \end{aligned}$ |
| $H^{26} \cong \mathbb{Z}_{2} \oplus \mathbb{Z} \oplus \mathbb{Z}$ | $\begin{aligned} & \bar{y}_{3}^{3} \bar{y}_{4}+\bar{y}_{3} \bar{y}_{4} \bar{y}_{6}+\bar{y}_{4}^{2} \bar{y}_{5}+\bar{y}_{4} \bar{y}_{9} \\ & \bar{y}_{3} \bar{y}_{4} \bar{y}_{6}+\bar{y}_{3} \bar{y}_{5}^{2}-3 \bar{y}_{4}^{2} \bar{y}_{5}+\bar{y}_{6} \bar{y}_{7} \\ & \bar{y}_{3}^{3} \bar{y}_{4}+\bar{y}_{3}^{2} \bar{y}_{7}+2 \bar{y}_{4}^{2} \bar{y}_{5}+\bar{y}_{4} \bar{y}_{9} \end{aligned}$ |
| $H^{28} \cong \mathbb{Z}_{2} \oplus \mathbb{Z} \oplus \mathbb{Z}$ | $\begin{aligned} & \bar{y}_{3}^{3} \bar{y}_{5}+\bar{y}_{3} \bar{y}_{5} \bar{y}_{6}+\bar{y}_{4} \bar{y}_{5}^{2}+\bar{y}_{5} \bar{y}_{9} \\ & 3 \bar{y}_{3}^{2} \bar{y}_{4}^{2}+5 \bar{y}_{4} \bar{y}_{5}^{2}+\bar{y}_{5} \bar{y}_{9} \\ & \bar{y}_{3}^{3} \bar{y}_{5}+2 \bar{y}_{3}^{2} \bar{y}_{4}^{2}+\bar{y}_{3} \bar{y}_{4} \bar{y}_{7}+\bar{y}_{3} \bar{y}_{5} \bar{y}_{6}+4 \bar{y}_{4} \bar{y}_{5}^{2}+ \\ & \bar{y}_{5} \bar{y}_{9} \end{aligned}$ |
| $H^{30} \cong \mathbb{Z}_{2} \oplus \mathbb{Z} \oplus \mathbb{Z}$ | $\begin{aligned} & -\bar{y}_{3}^{5}-\bar{y}_{3} \bar{y}_{4}^{3}+\bar{y}_{4} \bar{y}_{5} \bar{y}_{6}+\bar{y}_{5}^{3}+\bar{y}_{6} \bar{y}_{9} \\ & -\bar{y}_{3}^{5}+\bar{y}_{3} \bar{y}_{4} \bar{y}_{5}-3 \bar{y}_{3} \bar{y}_{4}^{3}+\bar{y}_{3} \bar{y}_{5} \bar{y}_{7}-\bar{y}_{4}^{2} \bar{y}_{7}+\bar{y}_{5}^{3} \\ & \bar{y}_{3}^{2} \bar{y}_{4} \bar{y}_{5}+\bar{y}_{3} \bar{y}_{5} \bar{y}_{7}+\bar{y}_{4}^{2} \bar{y}_{7}+\bar{y}_{4} \bar{y}_{5} \bar{y}_{6} \end{aligned}$ |
| $H^{32} \cong \mathbb{Z}_{2} \oplus \mathbb{Z}_{2} \oplus \mathbb{Z} \oplus \mathbb{Z}$ | $\begin{aligned} & \bar{y}_{3}^{4} \bar{y}_{4}+\bar{y}_{3}^{2} \bar{y}_{4} \bar{y}_{6}+\bar{y}_{3} \bar{y}_{4}^{2} \bar{y}_{5}+\bar{y}_{3} \bar{y}_{4} \bar{y}_{9} \\ & \bar{y}_{3}^{2} \bar{y}_{5}^{2}+\bar{y}_{3} \bar{y}_{6} \bar{y}_{7}+\bar{y}_{4} \bar{y}_{5} \bar{y}_{7}+\bar{y}_{5}^{2} \bar{y}_{6}+\bar{y}_{7} \bar{y}_{9} \\ & \bar{y}_{3}^{4} \bar{y}_{4}+\bar{y}_{3} \bar{y}_{4}^{2} \bar{y}_{5}+\bar{y}_{3} \bar{y}_{4} \bar{y}_{9}+\bar{y}_{3} \bar{y}_{6} \bar{y}_{7}+\bar{y}_{4}^{4} \\ & \bar{y}_{3}^{2} \bar{y}_{4} \bar{y}_{6}+\bar{y}_{3} \bar{y}_{4}^{2} \bar{y}_{5}-\bar{y}_{4}^{4}+\bar{y}_{4} \bar{y}_{5} \bar{y}_{7}+\bar{y}_{5}^{2} \bar{y}_{6} \end{aligned}$ |
| $H^{34} \cong \mathbb{Z}_{38} \oplus \mathbb{Z}$ | $\begin{aligned} & \bar{y}_{3}^{4} \bar{y}_{5}+\bar{y}_{3} \bar{y}_{4} \bar{y}_{5}^{2}+\bar{y}_{3} \bar{y}_{5} \bar{y}_{9}+2 \bar{y}_{4}^{3} \bar{y}_{5} \\ & \bar{y}_{3}^{4} \bar{y}_{5}+2 \bar{y}_{3}^{3} \bar{y}_{4}^{2}+\bar{y}_{3}^{2} \bar{y}_{4} \bar{y}_{7}+9 \bar{y}_{4}^{3} \bar{y}_{5}+\bar{y}_{4} \bar{y}_{6} \bar{y}_{7} \end{aligned}$ |
| $H^{36} \cong \mathbb{Z}_{2} \oplus \mathbb{Z}_{2} \oplus \mathbb{Z}_{20} \oplus \mathbb{Z}$ | $\begin{aligned} & \bar{y}_{3}^{3} \bar{y}_{4} \bar{y}_{5}+\bar{y}_{3} \bar{y}_{4} \bar{y}_{5} \bar{y}_{6}+\bar{y}_{4}^{2} \bar{y}_{5}^{2}+\bar{y}_{4} \bar{y}_{5} \bar{y}_{9} \\ & 13 \bar{y}_{3}^{6}+30 \bar{y}_{3}^{2} \bar{y}_{4}^{3}+\bar{y}_{3} \bar{y}_{6} \bar{y}_{9}+\bar{y}_{4} \bar{y}_{5}^{2}+\bar{y}_{4} \bar{y}_{5} \bar{y}_{9}+ \\ & \bar{y}_{6}^{3} \\ & -4 \bar{y}_{3}^{6}+\bar{y}_{3}^{3} \bar{y}_{4} \bar{y}_{5}-11 \bar{y}_{3}^{2} \bar{y}_{4}^{3}+\bar{y}_{3} \bar{y}_{6} \bar{y}_{9}+\bar{y}_{5} \bar{y}_{6} \bar{y}_{7} \\ & 13 \bar{y}_{3}^{6}+\bar{y}_{3}^{3} \bar{y}_{4} \bar{y}_{5}+28 \bar{y}_{3}^{2} \bar{y}_{4}^{3}-\bar{y}_{4}^{2} \bar{y}_{5}^{2}+\bar{y}_{4} \bar{y}_{5} \bar{y}_{9} \\ & \hline \end{aligned}$ |

$$
\begin{equation*}
H^{*}\left(E_{6} / P\right)=\mathbb{Z}\left[\omega_{2}, y_{3}, y_{4}, y_{6}\right] /\left\langle r_{6}, r_{8}, r_{9}, r_{12}\right\rangle \text {, where } \tag{5.12}
\end{equation*}
$$

$$
\begin{aligned}
& r_{6}=2 y_{6}+y_{3}^{2}-3 \omega_{2}^{2} y_{4}+2 \omega_{2}^{3} y_{3}-\omega_{2}^{6} ; \\
& r_{8}=3 y_{4}^{2}-6 \omega_{2} y_{3} y_{4}+\omega_{2}^{2} y_{6}+5 \omega_{2}^{2} y_{3}^{2}-2 \omega_{2}^{5} y_{3} ; \\
& r_{9}=2 y_{3} y_{6}-\omega_{2}^{3} y_{6} ; \\
& r_{12}=y_{6}^{2}-y_{4}^{3} .
\end{aligned}
$$

$$
\begin{align*}
& H^{*}\left(E_{7} / P\right)=\mathbb{Z}\left[\omega_{2}, y_{3}, y_{4}, y_{5}, y_{6}, y_{7}, y_{9}\right] / \quad\left\langle r_{j}\right\rangle_{j \in \Lambda} \text { where } \Lambda=\{6,8,9,10,12,14,18\}, \\
& \quad r_{6}=2 y_{6}+y_{3}^{2}+2 \omega_{2} y_{5}-3 \omega_{2}^{2} y_{4}+2 \omega_{2}^{3} y_{3}-\omega_{2}^{6} ; \\
& r_{8}=3 y_{4}^{2}-2 y_{3} y_{5}+2 \omega_{2} y_{7}-6 \omega_{2} y_{3} y_{4}+\omega_{2}^{2} y_{6}+5 \omega_{2}^{2} y_{3}^{2}+2 \omega_{2}^{3} y_{5}-2 \omega_{2}^{5} y_{3} ; \\
& r_{9}=2 y_{9}+2 y_{4} y_{5}-2 y_{3} y_{6}-4 \omega_{2} y_{3} y_{5}-\omega_{2}^{2} y_{7}+\omega_{2}^{3} y_{6}+2 \omega_{2}^{4} y_{5} ; \\
& r_{10}=y_{5}^{2}-2 y_{3} y_{7}+\omega_{2}^{3} y_{7} ; \\
& r_{12}=y_{6}^{2}+2 y_{5} y_{7}-y_{4}^{3}+2 y_{3} y_{9}+2 y_{3} y_{4} y_{5}+2 \omega_{2} y_{5} y_{6}-6 \omega_{2} y_{4} y_{7}+\omega_{2}^{2} y_{5}^{2} ; \\
& r_{14}=y_{7}^{2}-2 y_{5} y_{9}+y_{4} y_{5}^{2} ; \\
& r_{18}=y_{9}^{2}+2 y_{5} y_{6} y_{7}-y_{4} y_{7}^{2}-2 y_{4} y_{5} y_{9}+2 y_{3} y_{5}^{3}-\omega_{2} y_{5}^{2} y_{7} .
\end{align*}
$$

Step 3. Computing with the Weyl invariants. In addition to (5.10) the parabolic subgroup $P$ on $G$ specified by table (5.9) induces also the fibration

$$
\begin{equation*}
P / T \stackrel{i}{\hookrightarrow} G / T \xrightarrow{\pi} G / P(\text { i.e.(4.3)) } \tag{5.14}
\end{equation*}
$$

where Schubert presentation of the cohomology of the base space $G / P$ has been decided by (5.11), (5.12) and (5.13). On the other hand, with

$$
P / T=S p(3) / T^{3}, S U(6) / T^{5} \text { or } S U(7) / T^{6} \text { for } G=F_{4}, E_{6} \text { or } E_{7}
$$

the cohomology of the fiber space $P / T$ is given by Theorem 5.1 as

$$
H^{*}(P / T)=\left\{\begin{array}{l}
\frac{\mathbb{Z}\left[\omega_{2}, \omega_{3}, \omega_{4}\right]}{\left.\mathbb{( c _ { 2 } , c _ { 4 } , c _ { 6 } \rangle}\right]} \text { if } G=F_{4}  \tag{5.15}\\
\frac{\left.\mathbb{L} \omega_{1}, \omega_{3}, \ldots, \omega_{2}\right]}{\left\langle c_{r}, 2 \leq r \leq n\right\rangle} \text { if } G=E_{n} \text { with } n=6,7 .
\end{array}\right.
$$

Thus, Theorem 4.3 is applicable to fashion the ring $H^{*}(G / T)$ in question from the known ones $H^{*}(P / T)$ and $H^{*}(G / P)$. To this end we need only to specify a system $\left\{\rho_{r}\right\}$ satisfying the constraints (4.5). The invariant theory of Weyl groups serves this purpose.

Recall that the Weyl group $W$ of $G$ can be identified with the subgroup of $\operatorname{Aut}\left(H^{2}(G / T)\right)$ generated by the automorphisms $\sigma_{i}, 1 \leq i \leq n$, whose action on the Schubert basis $\left\{\omega_{1}, \ldots, \omega_{n}\right\}$ of $H^{2}(G / T)$ is given by the Cartan matrix $\left(a_{i j}\right)_{n \times n}$ of $G$ as

$$
\sigma_{i}\left(\omega_{k}\right)=\left\{\begin{array}{l}
\omega_{i} \text { if } k \neq i ;  \tag{5.16}\\
\omega_{i}-\sum_{1 \leq j \leq n} a_{i j} \omega_{j} \text { if } k=i .
\end{array}\right.
$$

Introduce for each $G=F_{4}, E_{6}$ and $E_{7}$ the polynomials $c_{r}(P)$ in $\omega_{1}, \ldots, \omega_{n}$ by the formula

$$
c_{r}(P):=\left\{\begin{array}{l}
e_{r}\left(o\left(\omega_{4}, W(P)\right)\right), 1 \leq r \leq 4 \text { if } G=F_{4}  \tag{5.17}\\
e_{r}\left(o\left(\omega_{n}, W(P)\right), 1 \leq r \leq n \text { if } G=E_{n}, n=6,7,\right.
\end{array}\right.
$$

where $o(\omega, W(P)) \subset H^{2}(G / T)$ denotes the $W(P)$-orbit through $\omega \in H^{2}(G / T)$, and where $e_{r}(o(\omega, W(P))) \in H^{2 r}(G / T)$ is the $r$ th elementary symmetric function on the set $o(\omega, W(P))$. For instance if $G=F_{4}, E_{6}$ we have by (5.16) that

$$
\begin{gathered}
o\left(\omega_{4}, W(P)\right)=\left\{\omega_{4}, \omega_{3}-\omega_{4}, \omega_{2}-\omega_{3}, \omega_{1}-\omega_{2}+\omega_{3}, \omega_{1}-\omega_{3}+\omega_{4}, \omega_{1}-\omega_{4}\right\}, \\
o\left(\omega_{6}, W(P)\right)=\left\{\omega_{6}, \omega_{5}-\omega_{6}, \omega_{4}-\omega_{5}, \omega_{2}+\omega_{3}-\omega_{4}, \omega_{1}+\omega_{2}-\omega_{3}, \omega_{2}-\omega_{1}\right\} .
\end{gathered}
$$

On the other hand, according to Bernstein-Gel'fand-Gel'fand [5, Proposition 5.1] the induced map $\pi^{*}$ in (4.3) injects, and satisfies the relation

$$
\operatorname{Im} \pi^{*}=H^{*}(G / T)^{W(P)}=H^{*}(G / P),
$$

implying $c_{r}(P) \in H^{*}(G / P)$. Since $c_{r}(P)$ is an explicit polynomial in the Schubert classes $\omega_{i}$ the Giambelli polynomials is functional to express it as a polynomial $g_{r}$ in the special Schubert classes on $H^{*}(G / P)$ given by table (5.8):

| $G$ | $F_{4}$ | $E_{6}$ | $E_{7}$ |
| :--- | :--- | :--- | :--- |
| $g_{2}$ | $4 \omega_{1}^{2}$ | $4 \omega_{2}^{2}$ | $4 \omega_{2}^{2}$ |
| $g_{3}$ | $y_{2}$ | $2 y_{3}+2 \omega_{2}^{3}$ | $2 y_{3}+2 \omega_{2}^{3}$ |
| $g_{4}$ | $3 y_{4}+2 \omega_{1} y_{3}$ | $3 y_{4}+\omega_{2}^{4}$ | $3 y_{4}+\omega_{2}^{4}$ |
| $g_{5}$ | $3 \omega_{2} y_{4}-2 \omega_{2}^{2} y_{3}+\omega_{2}^{5}$ | $2 y_{5}+3 \omega_{2} y_{4}-2 \omega_{2}^{2} y_{3}+\omega_{2}^{5}$ |  |
| $g_{6}$ | $y_{6}$ | $y_{6}$ | $y_{6}+2 \omega_{2} y_{5}$ |
| $g_{7}$ |  | $y_{7}$ |  |

Up to now we have accumulated sufficient information to show Theorem 5.4.
Proof of Theorem 5.4 For each $G=F_{4}, E_{6}$ or $E_{7}$ Schubert presentations for the cohomologies of the base $G / P$ and of the fiber $P / T$ have been determined by (5.11)-(5.13) and (5.15), respectively, while a system $\left\{\rho_{r}\right\}$ satisfying the relation (4.5) is seen to be $\rho_{r}:=c_{r}(P)-g_{r}$. Therefore, Theorem 4.3 is directly applicable to formulate a presentation of the ring $H^{*}(G / T)$. The results can be further simplified to yield the desired formulae (5.5)-(5.7) by the following observations:
(a) Certain Schubert classes $y_{k}$ from the base space $G / P$ can be eliminated against appropriate relations of the type $\rho_{k}$, e.g. if $G=E_{7}$ the generators $y_{6}, y_{7}$ and the relations $\rho_{6}, \rho_{7}$ can be excluded by the formulae of $g_{6}$ and $g_{7}$, which implies that $y_{6}=c_{6}-2 \omega_{2} y_{5}$ and $y_{7}=c_{7}$, respectively;
(b) Without altering the ideal, higher degree relations of the type $r_{i}$ may be simplified modulo the lower degree ones by the following fact. For two ordered sequences $\left\{f_{i}\right\}_{1 \leq i \leq n}$ and $\left\{h_{i}\right\}_{1 \leq i \leq n}$ of a graded polynomial ring with

$$
\operatorname{deg} f_{1}<\cdots<\operatorname{deg} f_{n} \text { and } \operatorname{deg} h_{1}<\cdots<\operatorname{deg} h_{n}
$$

write $\left\{h_{i}\right\}_{1 \leq i \leq n} \sim\left\{f_{i}\right\}_{1 \leq i \leq n}$ to denote the statements that $\operatorname{deg} h_{i}=\operatorname{deg} f_{i}$ and that $\left(f_{i}-h_{i}\right) \in\left\langle f_{j}\right\rangle_{1 \leq j<i}$. Then $\left\{f_{i}\right\}_{1 \leq i \leq n} \sim\left\{h_{i}\right\}_{1 \leq i \leq n}$ implies that $\left\langle h_{1}, \ldots, h_{n}\right\rangle$ $=\left\langle f_{1}, \ldots, f_{n}\right\rangle$.

### 5.3 A Type Free Characterization of the Ring $H^{*}(G / T)$

For an 1-connected simple Lie group $G$ with rank $n$ denote by $D(G) \subset H^{*}(G / T)$ the ideal of decomposable elements. Let $h(G)$ be the cardinality of a basis of the quotient group $H^{*}(G / T) / D(G)$ and set $m=h(G)-n-1$. The results of Theorems 5.1, 5.2 and 5.4 can be summarized into one formula, without referring to the types of the group $G$ (see [21, Theorems 1.2 and 1.3]).

Theorem 5.5 For each simple Lie group $G$ there exist a set $\left\{y_{1}, \ldots, y_{m}\right\}$ of $m$ Schubert classes on $G / T$ with $2<\operatorname{deg} y_{1}<\cdots<\operatorname{deg} y_{m}$, so that the inclusion $\left\{\omega_{1}, \ldots, \omega_{n}, y_{1}, \ldots, y_{m}\right\} \in H^{*}(G / T)$ induces the Schubert presentation

$$
\begin{equation*}
H^{*}(G / T)=\mathbb{Z}\left[\omega_{1}, \ldots, \omega_{n}, y_{1}, \ldots, y_{m}\right] /\left\langle e_{i}, f_{j}, g_{j}\right\rangle_{1 \leq i \leq k ; 1 \leq j \leq m} \tag{5.19}
\end{equation*}
$$

where
(i) $k=n-m$ for all $G \neq E_{8}$ but $k=n-m+2$ for $G=E_{8}$;
(ii) $e_{i} \in\left\langle\omega_{1}, \ldots, \omega_{n}\right\rangle, 1 \leq i \leq k$;
(iii) the pair $\left(f_{j}, g_{j}\right)$ of polynomials is related to the Schubert class $y_{j}$ in the fashion

$$
f_{j}=p_{j} \cdot y_{j}+\alpha_{j}, g_{j}=y_{j}^{k_{j}}+\beta_{j}, 1 \leq j \leq m
$$

where $p_{j} \in\{2,3,5\}$ and $\alpha_{j}, \beta_{j} \in\left\langle\omega_{1}, \ldots, \omega_{n}\right\rangle$;
(iv) ignoring the ordering, the sequence $\left\{\operatorname{deg} e_{i}, \operatorname{deg} g_{j}\right\}_{1 \leq i \leq k ; 1 \leq j \leq m}$ of integers agrees with the degree sequence of the basic Weyl invariants of the group $G$ (over the field of rationals).

Concerning assertion (iv) we remark that for each simple Lie group $G$ the degree sequence, as well as explicit formulae, of the basic Weyl invariants $P_{1}, \ldots, P_{n}$ of $G$ has been determined by Chevalley and Mehta [12, 47].

## 6 Further Remarks on the Characteristics

6.1. Certain parameter spaces of the geometric figures concerned by Schubert [51, Chap. 4] may fail to be flag manifolds, but can be constructed by performing finite number of blow-ups on flag manifolds along the centers again in flag manifolds, see the examples in Fulton [27, Example 14.7.12], or in [18] for the construction of the parameter spaces of the complete conics and quadrics on the 3 -space $\mathbb{P}^{3}$. As results
the relevant characteristics can be computed from those of flag manifolds via strict transformations (e.g. [18, Examples 5.11; 5.12]).
6.2. As the intersection multiplicities of certain Schubert varieties on $G / P$, the characteristics $a_{w_{1}, \ldots, w_{k}}^{w}$ are always non-negative by Van der Waerden [59]. Due to the importance of these numbers in geometry their effective computability (rather than positivity) had been the top priority in the classical approach [51, 52, 54], see also Fulton [27, 14.7].

Motivated by the Littlewood-Richardson rule [45] for the structure constants of the Grassmannian $G_{n, k}$ a remarkable development of Schubert calculus has taken place in algebraic combinatorics since 1970's, where the main concern is to find combinatorial descriptions of characteristics by which the positivity become transparent. This idea has inspired beautiful results on the enumerations of Yong tableaux, Mondrian tableaux, Chains in the Bruhat order, and puzzles by Buch, Graham, Coskun, Knutson and Tao [7, 14, 28, 43, 44], greatly enriched the classical Schubert calculus.
6.3. According to Van der Waerden [59] and Weil [60, p. 331] Hilbert's 15 th problem has been solved satisfactorily. In particular, in the context of modern intersection theory (e.g. [27, 29]) rigorous treatment of the major enumerative results of Schubert [51] had been completed independently by many authors (e.g. [1, 41, 42, 49]) ${ }^{2}$; granted with the basis theorem the characteristics of flag manifolds can be evaluated uniformly by the formula (2.2), while the Schubert presentations of the cohomology rings of flag manifolds have also been available (e.g. [3, 6, 20, 21, 46]).

However, Schubert calculus remains a vital and powerful tool in constructing the cohomologies of much broad spaces, such as the homogeneous spaces $G / H$ associated to Lie groups $G$. In contrast to the basis theorem (i.e. Theorem 1.2) the cohomologies of such spaces may be nontrivial in odd degrees, and may contain torsion elements. Nevertheless, inputting the formula (5.19) into the Koszul complex

$$
E_{2}^{*, *}(G)=H^{*}(G / T) \otimes H^{*}(T)
$$

associated to the fibration $G \rightarrow G / T$ a unified construction of the integral cohomology rings of all the 1-connected simple Lie groups $G$ has been carried out by Duan and Zhao in [23]. In addition, the formula (2.2) of the characteristics has been extended to evaluate the Steenrod operators on the $\bmod p$ cohomologies of flag manifolds [24], and of the exceptional Lie groups [25].
6.4. As is of today Schubert calculus has entered the intersection of several rapidly developing fields of mathematics, and has been generalized to the studies of other generalized cohomology theories, such as equivariant, quantum cohomology, K-theory, and cobordism, all of them are different deformations of the ordinary cohomology [31]. In this regard the present paper is by no means a comprehensive survey on the contemporary Schubert calculus. It illustrates a passage from the Cartan matrices of

[^8]Lie groups to the cohomology of homogeneous spaces, where Schubert's characteristics play a central role.

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# Asymmetric Function Theory 

Oliver Pechenik and Dominic Searles


#### Abstract

The classical theory of symmetric functions has a central position in algebraic combinatorics, bridging aspects of representation theory, combinatorics, and enumerative geometry. More recently, this theory has been fruitfully extended to the larger ring of quasisymmetric functions, with corresponding applications. Here, we survey recent work extending this theory further to general asymmetric polynomials.


Keywords Symmetric function • Quasisymmetric function • Schubert
polynomial • Slide polynomial • Demazure character • Demazure atom • Quasikey polynomial

## 1 The Three Worlds: Symmetric, Quasisymmetric, and General Polynomials

One of the gems of 20th-century mathematics is theory of symmetric functions and symmetric polynomials, as expounded in the classic textbooks [58, 59, 78]. In addition to its intrinsic beauty, this theory has important applications in representation theory, algebraic geometry, and combinatorics.

First, we will review aspects of this symmetric function theory. Then, we discuss the more general theory of quasisymmetric functions and polynomials, a very active area of contemporary research. Finally, we turn to the combinatorial theory of general asymmetric polynomials. While this seems naively like a very simple object, the polynomial ring turns out to have a rich and beautiful combinatorial structure analogous to that of the symmetric and quasisymmetric worlds, but far less explored.

[^9]While the theory of symmetric polynomials is a textbook subject and that of quasisymmetric polynomials has recently been nicely expounded in [61], we believe this article is the first survey of corresponding developments in the asymmetric world.

In each world, we will consider a variety of additive bases for the algebra in question. The power of the combinatorial theory comes from the bases having the following three characteristics:
(1) positive combinatorial rules for change of basis,
(2) positive combinatorial multiplication rules in various bases, and
(3) algebraic/geometric interpretations of the basis elements.

We will attempt to elucidate all three properties in each case, to the extent that such properties are known to hold. In some cases, our descriptions of geometric interpretations require background that is beyond the scope of this article to recall; readers without such background knowledge may safely skip such episodes.

## 2 The Symmetric World

Consider the $\mathbb{Z}$-algebra Poly $_{n}:=\mathbb{Z}\left[x_{1}, \ldots, x_{n}\right]$ of integral multivariate polynomials. It carries a natural action of the symmetric group $\mathcal{S}_{n}$ on $n$ letters, where the simple transposition $(i, i+1)$ acts on $f \in$ Poly $_{n}$ by swapping the variables $x_{i}$ and $x_{i+1}$. Let $\operatorname{Sym}_{n}:=\operatorname{Poly}_{n}^{\mathcal{S}_{n}}$, the $\mathcal{S}_{n}$-invariants. It is easy to see that $\operatorname{Sym}_{n}$ is a subring of Poly ${ }_{n}$; we call it the ring of symmetric polynomials in $n$ variables. $\operatorname{Sym}_{n}$ is naturally a graded ring, inheriting the grading by degree from Poly $_{n}$. We denote the degree $m$ homogeneous piece of a graded ring $R$ by $R^{(m)}$.

For $m \leq n$, we can map $\operatorname{Sym}_{n}$ onto $\operatorname{Sym}_{m}$ by setting the last $n-m$ variables equal to 0 . The inverse limit of the $\left\{\mathrm{Sym}_{n}\right\}$ with respect to these restriction maps is called the ring of symmetric functions Sym, although its elements are not functions, but rather formal power series in infinitely-many variables. Classically, one generally prefers to study Sym; however, we will usually prefer the essentially equivalent theory of $\mathrm{Sym}_{n}$, as it extends more naturally to the asymmetric setting that is our focus.

We will consider four of the most important additive bases of $\mathrm{Sym}_{n}$ : the monomial, elementary, homogeneous, and Schur bases. Given bases $A, B$ of a free $\mathbb{Z}$-module, we say that $A$ refines $B$ if every element of $B$ can be written as a sum of elements of $A$ with nonnegative coefficients. For a much deeper exposition of symmetric function theory than we provide, see any of the textbooks [58, 59, 78], which also provide proofs that we omit. The relations among these four bases are illustrated in Fig. 1.

For a weak composition $a$ (i.e., an infinite sequence of nonnegative integers with finite sum), define a monomial

$$
\mathbf{x}^{a}:=x_{1}^{a_{1}} x_{2}^{a_{2}} \cdots
$$

For a partition $\lambda$ (i.e., a weakly decreasing weak composition), let

$$
m_{\lambda}:=\sum_{a} \mathbf{x}^{a}
$$

where the sum is over all distinct weak compositions that can be obtained by rearranging the parts of $\lambda$. If $\mathbf{x}^{a}$ is a monomial of the symmetric function $f$, then necessarily $\mathbf{x}^{b}$ is also a monomial of $f$ for every $b$ that can be obtained by rearranging the parts of $a$. Thus $f$ can be written uniquely as a finite sum of the monomial symmetric functions $m_{\lambda}$. Therefore $\left\{m_{\lambda}\right\}$ is a $\mathbb{Z}$-linear basis of Sym and the dimension of $\operatorname{Sym}^{(m)}$ is the number of partitions of $m$.

We may consider a second action of $\mathcal{S}_{n}$ on Poly $y_{n}$ where a permutation acts by permuting variables and then multiplying by the sign of the permutation. Note that this merely amounts to twisting the original action by tensoring with the 1 -dimensional sign representation of $\mathcal{S}_{n}$. The invariants of this twisted action are the alternating polynomials in $n$ variables, $v_{n} \mathrm{Sym}_{n}$. These are precisely the polynomials where setting any two variables equal yields 0 . The sum of two alternating polynomials is alternating, but the product is generally not. Hence $v_{n} \operatorname{Sym}_{n}$ is not a subring of Poly ${ }_{n}$, although it is a module over $\operatorname{Sym}_{n}$. As with $\operatorname{Sym}_{n}, v_{n} \operatorname{Sym}_{n}$ is graded by degree, although for technical reasons one might prefer to shift the degree by $\binom{n}{2}$.

Let $v_{n}:=\prod_{1 \leq i<j \leq n}\left(x_{i}-x_{j}\right)$ be the Vandermonde determinant. This is an alternating polynomial and moreover divides every other alternating polynomial. The quotients are necessarily symmetric. Hence every alternating polynomial can be written as $v_{n}$ times a symmetric polynomial. (This fact justifies the notation $v_{n} \operatorname{Sym}_{n}$ for the module of alternating polynomials.)

For a weak composition $a$ of length $n$, define

$$
\tilde{j}_{a}:=\sum_{\sigma \in S_{n}} \operatorname{sgn}(\sigma) x^{\sigma(a)}
$$

Note that if $\mathbf{x}^{a}$ is a term of the alternating polynomial $f$, then so is every other term of $\tilde{j}_{a}$. Moreover if $a$ has any repeated parts, then clearly $\tilde{j}_{a}=0$. Hence $v_{n} \operatorname{Sym}_{n}$ has a natural basis of polynomials $\tilde{j}_{\theta}$, for $\theta$ ranging over strict partitions, that is partitions with distinct parts.

Every strict partition of length $n$ may be written uniquely as $\delta+\lambda$, where $\delta=$ $(n-1, n-2, \ldots, 0), \lambda$ is a partition, and the sum is componentwise. We write $j_{\lambda}:=\tilde{j}_{\delta+\lambda}$, to obtain a basis of $v_{n}$ Sym indexed by partitions. That is, the dimension of the space of alternating polynomials of degree $m+\binom{n}{2}$ equals the dimension of the space of symmetric polynomials of degree $m$. Indeed, we can even identify the isomorphism; it is just multiplication by $v_{n}=\tilde{j}_{\delta}=j_{(0)}$. If we shifted the grading of $v_{n} \operatorname{Sym}_{n}$ as suggested above (so that $v_{n}$ is in degree 0 ), then multiplication by $v_{n}$ is an isomorphism $\operatorname{Sym}_{n} \rightarrow v_{n} \mathrm{Sym}_{n}$ of graded Sym-modules.

The basis of $\operatorname{Sym}_{n}$ obtained by pulling back the $j_{\lambda}$ basis of $v_{n} \operatorname{Sym}_{n}$ is not the basis of monomial symmetric polynomials, but rather something more interesting. These important objects

$$
s_{\lambda}:=\frac{j_{\lambda}}{v_{n}}
$$

are called the Schur polynomials.
Although the Schur polynomials are clearly symmetric and hence can be expanded in the monomial basis, it is remarkable that the non-zero expansion coefficients are uniformly positive. A recurring theme in this survey will be such instances of positive basis changes between a priori unrelated bases.

A combinatorial formula manifesting the monomial-positivity of Schur polynomials was given by Littlewood. Given a partition $\lambda=\left(\lambda_{1}, \lambda_{2}, \ldots\right)$, we identify $\lambda$ with its English-orientation Young diagram, consisting of $\lambda_{1}$ left-justified boxes in the top row, $\lambda_{2}$ left-justified boxes in the second row, etc. A semistandard (Young) tableau of shape $\lambda$ is an assignment of a positive integer to each box of the Young diagram such that the labels weakly increase left to right across rows and strictly increase down columns. The weight of a tableau $T$ is the weak composition $\mathrm{wt}(T):=\left(a_{1}, a_{2}, \ldots\right)$, where $a_{i}$ records the number of boxes labeled $i$.

Theorem 2.1 (Littlewood, [47, 48]) For any partition $\lambda$, we have

$$
s_{\lambda}=\sum_{T \in \operatorname{SSYT}(\lambda)} \mathbf{x}^{\mathrm{wt}(T)}
$$

Example 2.2 We have $s_{(2,1)}\left(x_{1}, x_{2}\right)=x_{1}^{2} x_{2}+x_{1} x_{2}^{2}$, owing to the two semistandard tableaux

$$
\begin{array}{|l|l|l|l|}
\hline 1 & 1 & \begin{array}{ll|l}
\hline 1 & 2 \\
\hline 2 & & \\
\hline 2 & & \\
\hline
\end{array} . \\
\hline
\end{array}
$$

We now turn to the last two bases of $\mathrm{Sym}_{n}$ that we will consider. For a partition $\lambda=\left(\lambda_{1}, \lambda_{2}, \ldots\right)$, we define the elementary symmetric function $e_{\lambda}$ by

$$
e_{\lambda}:=\prod_{i} s_{\lambda_{i}}
$$

and the (complete) homogeneous symmetric function by

$$
h_{\lambda}:=\prod_{i} s_{1^{\lambda_{i}}} .
$$

It is not obvious that either of these families yields a basis of $\mathrm{Sym}_{n}$; nevertheless, each of them does, as was originally established by Isaac Newton. It also is not obvious that, like the Schur basis, the elementary and homogeneous bases expand positively in the $m_{\lambda}$. However, in fact, something much stronger is true: each $e_{\lambda}$ and
each $h_{\lambda}$ is a positive sum of Schur polynomials. This positivity is a consequence of an even more remarkable positivity property:

Theorem 2.3 The Schur basis of $\mathrm{Sym}_{n}$ has positive structure coefficients. In other words, for any partitions $\lambda$ and $\mu$, the product $s_{\lambda} \cdot s_{\mu}$ expands as a positive sum of Schur polynomials.

Corollary 2.4 For any $\lambda, e_{\lambda}$ and $h_{\lambda}$ are both Schur-positive, and hence monomialpositive.

Proof By Theorem 2.3, any product of Schur polynomials is Schur-positive. Since $e_{\lambda}$ and $h_{\lambda}$ are defined as products of special Schur polynomials, they are then Schurpositive. By Theorem 2.1, Schur polynomials are monomial-positive. Hence, any Schur-positive polynomial, in particular $e_{\lambda}$ or $h_{\lambda}$, is also monomial-positive.

By commutativity and the definitions, it is transparent that both the elementary and homogeneous bases of $\mathrm{Sym}_{n}$ also have positive structure coefficients.

There are a variety of distinct ways to establish Theorem 2.3. The most fundamental explanations involve interpreting the theorem representation-theoretically or geometrically.

For a representation-theoretic approach, one can establish that Sym is isomorphic to the ring of polynomial representations of the general linear group in such a way that the Schur functions are in one-to-one correspondence with the irreducible representations. Under this identification, decomposing the tensor product of two irreducible representations into irreducibles corresponds to expanding the product of two Schur functions in the Schur basis. Hence, Theorem 2.3 follows. Similarly, one can identify the Schur functions of homogeneous degree $k$ with the irreducible representations of the symmetric group $\mathcal{S}_{n}$ in such a way that multiplying Schur functions corresponds to taking an 'induction product' of the corresponding representations. Again, since the induction product representation is necessarily a direct sum of irreducible representations, we recover Theorem 2.3. For more details on these representation-theoretic proofs, see, e.g., [23, 59].

A geometric approach is to identify Sym with the Chow ring of complex Grassmannians, the classifying spaces for complex vector bundles. A Grassmannian comes with a natural cell decomposition by certain subvarieties called Schubert varieties, yielding an effective basis of the Chow ring. Under the identification with Sym, this basis corresponds to the Schur polynomials. Multiplying Schur polynomials then corresponds to intersection product on Schubert varieties, and again Theorem 2.3 follows. For more details on these geometric constructions, see, e.g., [23, 25, 59].

The above interpretations of Theorem 2.3 deepen its significance and provide relatively easy proofs. Nonetheless, Theorem 2.3 is on its face a purely combinatorial statement and so one might hope it also had a purely combinatorial proof. Indeed, such a proof exists. Even better, it gives an explicit transparently-positive formula for the positive integers appearing in the Schur expansion. This formula can then be combined with the algebraic and geometric interpretations above to compute with and to better understand aspects of representation theory and enumerative geometry.
(Indeed, the combinatorics described here can be directly interpreted in terms of representations of quantum groups via M. Kashiwara's theory of crystal bases [35-37]; see [12] for an excellent combinatorially-flavored introduction to these connections.)

We write $\lambda \subseteq \nu$ to mean that the Young diagram of the partition $\lambda$ is a subset of that for $\nu$. The set-theoretic difference is called the skew Young diagram $\nu / \lambda$. A skew semistandard tableau is a filling of a skew diagram by positive integers, such that rows weakly increase and columns strictly increase. Define the content of a skew tableau as for tableaux of partition shape. The reading word of a (skew) tableau $T$ is the word given by reading the rows of $T$ from top to bottom and from right to left (like the ordinary reading order in Arabic or Hebrew). We say that $T$ is Yamanouchi if every initial segment of its reading word contains at least as many is as $(i+1) \mathrm{s}$, for each positive integer $i$.

Theorem 2.5 (Littlewood-Richardson rule) For partitions $\lambda$ and $\mu$, we have

$$
s_{\lambda} \cdot s_{\mu}=\sum_{\nu} c_{\lambda, \mu}^{\nu} s_{\nu}
$$

where $c_{\lambda, \mu}^{\nu}$ counts the number of Yamanouchi semistandard tableaux of skew shape $\nu / \lambda$ and content $\mu$.

Example 2.6 To compute the structure $\operatorname{coefficient} c_{(2,1),(2,1)}^{(3,2,1)}$ via Theorem 2.5, we consider fillings of the skew shape $(3,2,1) /(2,1)$ :

using two 1 s and one 2 . There are three such fillings

all of which are semistandard. However, only the first two are Yamanouchi, as the reading word of the third is 211 , which has a 2 before any 1 . Hence, $c_{(2,1),(2,1)}^{(3,2,1)}=2$.

## 3 The Quasisymmetric World

In this section, we consider a third action of $\mathcal{S}_{n}$ on Poly ${ }_{n}$. Here, the simple transposition $(i, i+1)$ acts on $f \in$ Poly $_{n}$ by swapping the variables $x_{i}$ and $x_{i+1}$ in only those terms where at most one of the two variables appears. The invariants of this action are the subalgebra $\mathrm{QSym}_{n}$ of quasisymmetric polynomials. Equivalently,
quasisymmetric polynomials are those polynomials $f$ such that, for all $a_{1}, \ldots, a_{k}$, the coefficient in $f$ of $x_{i_{1}}^{a_{1}} \cdots x_{i_{k}}^{a_{k}}$ equals the coefficient in $f$ of $x_{j_{1}}^{a_{1}} \cdots x_{k}^{a_{k}}$, whenever $i_{1}<\cdots<i_{k}$ and $j_{1}<\cdots<j_{k}$. Analogously to the symmetric case, one can also define the ring QSym of quasisymmetric functions in infinitely-many variables as the inverse limit of the $\mathrm{QSym}_{n}$, but again our focus is on the essentially equivalent finite-variable case. For a much more detailed survey than we provide here of the state of the art in quasisymmetric function theory, see [61].

We will consider three important bases of $\mathrm{QSym}_{n}$ : the monomial, fundamental and quasiSchur bases. The relations among these three bases are illustrated in Fig. 2.

A strong composition $\alpha$ is a finite sequence of positive integers; we identify $\alpha$ with the weak composition obtained by appending infinitely many 0 s to the end of $\alpha$. For any weak composition $a$, its positive part is the strong composition $a^{+}$given by deleting all 0 s.

For any strong composition $\alpha$, define the monomial quasisymmetric polynomial $M_{\alpha}$ by

$$
M_{\alpha}\left(x_{1}, \ldots, x_{n}\right):=\sum_{b} \mathbf{x}^{b} \in \operatorname{QSym}_{n},
$$

where the sum is over all weak compositions $b$ with $b^{+}=\alpha$ and whose entries after position $n$ are all zero. Clearly, the monomial quasisymmetric polynomials yield a basis of QSym $_{n}$.

Example 3.1 We have

$$
M_{13}\left(x_{1}, x_{2}, x_{3}\right)=\mathbf{x}^{130}+\mathbf{x}^{103}+\mathbf{x}^{013} \in \mathrm{QSym}_{3} .
$$

Note in particular that this polynomial is not an element of $\mathrm{Sym}_{3}$.
It is clear that the monomial basis of $\mathrm{QSym}_{n}$ must have positive structure coefficients, as does the monomial basis of $\mathrm{Sym}_{n}$ discussed in Sect. 2. However, these structure coefficients are slightly more interesting than those for the monomial symmetric functions. They are given by the overlapping shuffles of M. Hazewinkel [29], which we now recall.

Let $A$ and $B$ be words in disjoint alphabets with $A$ of length $m$ and $B$ of length $n$. An overlapping shuffle of $A$ and $B$ is a surjection

$$
t:\{1,2, \ldots, m+n\} \rightarrow\{1,2, \ldots, k\}
$$

(for some $\max \{m, n\} \leq k \leq m+n$ ) such that

$$
t(i)<t(j) \text { whenever } i<j \leq m \text { or } m<i<j .
$$

We write $A \amalg_{o} B$ for the overlapping shuffle product of $A$ and $B$, the formal sum of all overlapping shuffles. The overlapping shuffle product $\alpha \amalg_{o} \beta$ of two strong
compositions $\alpha$ and $\beta$ is given by treating the strong compositions as words in disjoint alphabets. Here, we identify an overlapping shuffle $t:\{1,2, \ldots, m+n\} \rightarrow$ $\{1,2, \ldots, k\}$ of $\alpha$ and $\beta$ with the strong composition $\gamma$ defined by

$$
\gamma_{i}:=\sum_{t(j)=i}(\alpha \beta)_{j},
$$

where $\alpha \beta$ denotes the concatenation of the two strong compositions.
Example 3.2 We compute the overlapping shuffle product of $(2)$ and $(1,2)$ :

$$
(2) Ш_{o}(1,2)=(2,1,2)+2 \cdot(1,2,2)+(3,2)+(1,4) \text {. }
$$

Although the relevant combinatorial construction was somewhat involved, the following positive multiplication formula is now essentially clear.

Theorem 3.3 For strong compositions $\alpha$ and $\beta$, we have

$$
M_{\alpha} \cdot M_{\beta}=\sum_{\gamma} c_{\alpha, \beta}^{\gamma} M_{\gamma}
$$

where $c_{\alpha, \beta}^{\gamma}$ is the multiplicity of $\gamma$ in the overlapping shuffle product $\alpha Ш_{o} \beta$.
Example 3.4 To compute $M_{(2)} \cdot M_{(1,2)}$ via Theorem 3.3, we compute the overlapping shuffle product of $(2)$ and $(1,2)$ as in Example 3.2. Then, the coefficients on the various strong compositions give the coefficients on the various monomial quasisymmetric polynomials in the product:

$$
M_{(2)} \cdot M_{(1,2)}=M_{(2,1,2)}+2 M_{(1,2,2)}+M_{(3,2)}+M_{(1,4)} .
$$

Given two strong compositions $\alpha$ and $\beta$, we say $\beta$ refines $\alpha$ and write $\beta \vDash \alpha$ if $\alpha$ can be obtained by summing consecutive entries of $\beta$, e.g. $(1,2,1) \vDash(1,3)$ but $(2,1,1) \not \models(1,3)$.

Define the fundamental quasisymmetric polynomial $F_{\alpha}$ by

$$
F_{\alpha}\left(x_{1}, \ldots, x_{n}\right):=\sum_{b} \mathbf{x}^{b}
$$

where the sum is over all distinct weak compositions $b$ with $b^{+} \vDash \alpha$ and whose entries after position $n$ are all zero.

Example 3.5 We have

$$
F_{13}\left(x_{1}, x_{2}, x_{3}\right)=\mathbf{x}^{130}+\mathbf{x}^{103}+\mathbf{x}^{013}+\mathbf{x}^{112}+\mathbf{x}^{121} .
$$

Theorem 3.6 The fundamental quasisymmetric polynomials expand positively in the monomial quasisymmetric polynomials:

$$
F_{\alpha}\left(x_{1}, \ldots, x_{n}\right)=\sum_{\beta \vDash \alpha} M_{\beta}\left(x_{1}, \ldots, x_{n}\right) .
$$

A sign of the fundamental nature of the fundamental quasisymmetric polynomials is that they also have positive structure coefficients. Their multiplication is also governed by a shuffle product, that of S. Eilenberg and S. Mac Lane [19]. Let $A$ and $B$ be words in the disjoint alphabets $\mathcal{A}$ and $\mathcal{B}$, respectively. Recall a shuffle of $A$ and $B$ is a permutation of the concatenation $A B$ such that the subword on the alphabet $\mathcal{A}$ is $A$ and the subword on $\mathcal{B}$ is $B$. Alternatively, if $A$ has length $m$ and $B$ has length $n$, we can think of a shuffle of $A$ and $B$ as a bijection

$$
s:\{1,2, \ldots, m+n\} \rightarrow\{1,2, \ldots, m+n\}
$$

such that

$$
s(i)<s(j) \text { whenever } i<j \leq m \text { or } m<i<j .
$$

The shuffle product of two strong compositions $\alpha$ and $\beta$ is obtained as follows. Let $\mathcal{A}$ denote the alphabet of odd integers and let $\mathcal{B}$ denote the alphabet of even integers. Let $A$ be the word in $\mathcal{A}$ consisting of $\alpha_{1}$ copies of $2 \ell(\alpha)-1$, followed by $\alpha_{2}$ copies of $2 \ell(\alpha)-3$, all the way to $\alpha_{\ell(\alpha)}$ copies of 1 . Likewise, let $B$ denote the word in $\mathcal{B}$ consisting of $\beta_{1}$ copies of $2 \ell(\beta)$, followed by $\beta_{2}$ copies of $2 \ell(\beta)-2$, all the way to $\beta_{\ell(\beta)}$ copies of 2 . Let $\operatorname{Sh}(A, B)$ denote the set of the $\binom{|\alpha|+|\beta|}{|\beta|}$ shuffles of $A$ and $B$. For each $C \in \operatorname{Sh}(A, B)$, let $\operatorname{Des}(C)$ denote the descent composition of $C$, i.e. the strong composition obtained by decomposing $C$ into maximal runs of increasing entries and letting $\operatorname{Des}(C)_{i}$ be the number of entries in the $i$ th increasing run of $C$. Finally, define the shuffle product $\alpha \amalg \beta$ of the strong compositions $\alpha$ and $\beta$ as the formal sum of strong compositions

$$
\alpha Ш \beta:=\sum_{C \in \operatorname{Sh}(A, B)} \operatorname{Des}(C) .
$$

Example 3.7 Let $\alpha=(2)$ and $\beta=(1,2)$. Then $A=11$ and $B=422$. We compute the set of shuffles of $A$ and $B$ :

$$
\begin{aligned}
\operatorname{Sh}(A, B)= & \{4|22| 11,4|2| 12|1,4| 122|1,14| 22|1,4| 2 \mid 112, \\
& 4|12| 12,14|2| 12,4|1122,14| 122,114 \mid 22\},
\end{aligned}
$$

where we have placed bars to indicate the decomposition of each shuffle into maximally increasing runs. The corresponding descent compositions are thus, respectively,

$$
\{(1,2,2),(1,1,2,1),(1,3,1),(2,2,1),(1,1,3),(1,2,2),(2,1,2),(1,4),(2,3),(3,2)\} .
$$

Hence, we have

$$
\begin{aligned}
(2) \amalg(1,2) & =2(1,2,2)+(1,1,2,1)+(1,3,1)+(2,2,1) \\
& +(1,1,3)+(2,1,2)+(1,4)+(2,3)+(3,2) .
\end{aligned}
$$

Note that this sum is not multiplicity-free.
Theorem 3.8 For strong compositions $\alpha$ and $\beta$, we have

$$
F_{\alpha} \cdot F_{\beta}=\sum_{\gamma} c_{\alpha, \beta}^{\gamma} F_{\gamma}
$$

where $c_{\alpha, \beta}^{\gamma}$ is the multiplicity of $\gamma$ in the ordinary shuffle product $\alpha \amalg \beta$.
Proving this theorem is given as Exercise 7.93 in [78].
Example 3.9 To compute $F_{(2)} \cdot F_{(1,2)}$ via Theorem 3.8, we compute the shuffle product of (2) and $(1,2)$ as in Example 3.7. Then, the coefficients on the various strong compositions give the coefficients on the various fundamental quasisymmetric polynomials in the product:

$$
F_{(2)} \cdot F_{(1,2)}=2 F_{(1,2,2)}+F_{(1,1,2,1)}+F_{(1,3,1)}+F_{(2,2,1)}+F_{(1,1,3)}+F_{(2,1,2)}+F_{(1,4)}+F_{(2,3)}+F_{(3,2)} .
$$

The final basis for $\mathrm{QSym}_{n}$ that we will consider is the basis of quasiSchur polynomials introduced in [30]. For a detailed and readable survey of work related to this basis, see [50]. For those unfamiliar with quasiSchur polynomials, the definition may appear quite complicated; it originates as an important and tractable piece of the theory of Macdonald polynomials. It is not transparent from this definition that the quasiSchur polynomials are quasisymmetric, much less that they yield a basis of $\mathrm{QSym}_{n}$.

First, we must extend the definition of the Young diagram of a partition to a general weak composition $a=\left(a_{1}, a_{2}, \ldots\right)$ : Draw $a_{i}$ left-justified boxes in row $i$. (Here, in accordance with our English orientation on Young diagrams for partitions, row 1 is the top row.) A (composition) tableau of shape $a$ is an assignment of a positive integer to each box of the Young diagram for $a$. (Sometimes, we will augment such a tableau with an extra "zeroth" column of boxes (called the basement) immediately left of the first column, and write $b_{i}$ for the positive integer labeling the basement box in row $i$. Basement entries do not contribute to the weight of a tableau.)

A triple of boxes in a composition tableau $T$ is a set of three boxes in one of the two following configurations:


Note, in particular, that a triple has exactly two boxes sharing a row and exactly two boxes sharing a column. We say a triple is inversion if it is not the case that its labels satisfy $X \leq Y \leq Z$.

A composition tableau is semistandard if
(S.1) entries do not repeat in a column,
(S.2) rows weakly decrease from left to right,
(S.3) every triple is inversion,
(S.4) entries in the first column equal their row indices.
(Note that, in the case of partition shape, this definition unfortunately does not coincide with the definition of semistandard tableaux we have given previously.) For a weak composition $a$, let $\mathfrak{A S S T}(a)$ denote the set of semistandard tableaux of shape $a$. The quasiSchur polynomial for the strong composition $\alpha$ is then given by

$$
\begin{equation*}
S_{\alpha}\left(x_{1}, \ldots, x_{n}\right)=\sum_{a^{+}=\alpha} \sum_{T \in \mathfrak{A S S T}(a)} \mathbf{x}^{\mathrm{wt}(T)}, \tag{3.1}
\end{equation*}
$$

where the first sum is over all weak compositions $a$ of length $n$ with positive part $\alpha$.
Example 3.10 For $\alpha=(1,3)$ and $n=3$, we have
$S_{(1,3)}\left(x_{1}, x_{2}, x_{3}\right)=\mathbf{x}^{130}+\mathbf{x}^{220}+\mathbf{x}^{103}+\mathbf{x}^{202}+2 \mathbf{x}^{112}+\mathbf{x}^{121}+\mathbf{x}^{211}+\mathbf{x}^{013}+\mathbf{x}^{022}$,
where the monomials are determined by the semistandard composition tableaux shown in Fig. 3.

From the given definition of quasiSchur polynomials, it is not clear that they are natural objects that we should expect to exhibit any nice properties. Nonetheless, they participate in two beautiful positive combinatorial rules for basis expansion.

Since $\operatorname{Sym}_{n} \subset$ QSym $_{n}$, we can ask how bases of $\operatorname{Sym}_{n}$ expand in bases of QSym $n$. First, observe the following straightforward formula for the $M_{\alpha}$-expansion of the monomial symmetric polynomial $m_{\lambda}$. For a weak composition $a$, we write $\overleftarrow{a}$ for the partition formed by sorting the entries of $a$ into weakly decreasing order.

Proposition 3.11 The monomial symmetric polynomials expand positively in the monomial quasisymmetric polynomials:


Fig. 1 The four bases of $\mathrm{Sym}_{n}$ considered here. The arrows denote that the basis at the head refines the basis at the tail. All four bases have positive structure coefficients


Fig. 2 The three bases of QSym ${ }_{n}$ considered here, together with some bases of Sym $n$ from Fig. 1. The arrows denote that the basis at the head refines the basis at the tail. The star-shaped nodes have positive structure coefficients


Fig. 3 The 10 semistandard composition tableaux associated to the quasiSchur polynomial $S_{(1,3)}\left(x_{1}, x_{2}, x_{3}\right)$. The quasiYamanouchi tableaux are shaded in blue, the initial tableaux in pink, and those that are both quasiYamanouchi and initial in green

$$
m_{\lambda}\left(x_{1}, \ldots, x_{n}\right)=\sum_{\overleftarrow{\alpha}=\lambda} M_{\alpha}\left(x_{1}, \ldots, x_{n}\right)
$$

The quasiSchur expansion of a Schur polynomial is beautifully parallel to the formula of Proposition 3.11.

Theorem 3.12 ([30]) The Schur polynomials expand positively in the quasiSchur polynomials:

$$
s_{\lambda}\left(x_{1}, \ldots, x_{n}\right)=\sum_{\overleftarrow{\alpha}=\lambda} S_{\alpha}\left(x_{1}, \ldots, x_{n}\right)
$$

Remark 3.13 Considering Fig. 2 together with Proposition 3.11 and Theorem 3.12, one might be tempted to define polynomials

$$
f_{\lambda}\left(x_{1}, \ldots, x_{n}\right)=\sum_{\overleftarrow{\alpha}=\lambda} F_{\alpha}\left(x_{1}, \ldots, x_{n}\right)
$$

Extrapolating from Fig. 2, it might appear plausible that $\left\{f_{\lambda}\right\}$ should form a basis of $\mathrm{Sym}_{n}$, perhaps even with positive structure coefficients. However, the polynomials $f_{\lambda}$ are in general not even symmetric!

For example, in four or more variables, we have by Theorem 3.6 and Proposition 3.11 that

$$
\begin{aligned}
f_{31} & =F_{31}+F_{13} \\
& =M_{31}+M_{211}+2 M_{121}+M_{112}+M_{13}+2 M_{1111} \\
& =m_{31}+m_{211}+2 m_{1111}+M_{121},
\end{aligned}
$$

a symmetric polynomial plus $M_{121}$.
To describe the expansion of quasiSchur polynomials into the fundamental basis, we isolate an important subclass of semistandard composition tableaux. Fix a strong composition $\alpha$ and consider $T \in \mathfrak{A S S T}(a)$ for some $a$ with $a^{+}=\alpha$. We say that $T$ is quasiYamanouchi if for every integer $i$ appearing in $T$, either

- an $i$ appears in the first column, or
- there is an $i+1$ weakly right of an $i$.

We say $T$ is initial if the set of integers $i$ appearing in $T$ is an initial segment of $\mathbb{Z}_{>0}$.
Theorem 3.14 The quasiSchur polynomials expand positively in the fundamental quasisymmetric polynomials:

$$
S_{\alpha}\left(x_{1}, \ldots, x_{n}\right)=\sum_{T} F_{\mathrm{wt}(T)}\left(x_{1}, \ldots, x_{n}\right),
$$

where the sum is over all initial quasiYamanouchi tableaux $T$ such that $T \in$ $\mathfrak{A S S T}(a)$ for some $a$ with $a^{+}=\alpha$.

A positive formula for the expansion of quasiSchur polynomials in fundamental quasisymmetric polynomials was first given in [30] in terms of standard augmented fillings. The formula in Theorem 3.14 above follows as a consequence of a result in [77]; we state the expansion in these terms for consistency with formulas in the upcoming sections.

Remark 3.15 Unlike the other two bases of $\mathrm{QSym}_{n}$ that we have considered, the quasiSchur basis does not have positive structure coefficients. For an example, see [30, Sect. 7.1]. However, [31] proves a slightly weaker form of positivity, giving a positive combinatorial formula for the quasiSchur expansion of the product of a quasiSchur polynomial by a Schur polynomial.

Just as the combinatorics of $\mathrm{Sym}_{n}$ is related to the representation theory of symmetric groups, the combinatorics of $\mathrm{QSym}_{n}$ turns out to be related to the representation theory of 0 -Hecke algebras (in type $A$ ). First, let us recall the standard Coxeter presentation of the symmetric group $\mathcal{S}_{n}$. It is easy to see that $\mathcal{S}_{n}$ is generated by the simple transpositions $s_{i}:=(i, i+1)$ for $1 \leq i<n$. With a little more effort, one establishes that a generating set of relations is given by

- $s_{i}^{2}=\mathrm{id}$,
- $s_{i} s_{j}=s_{j} s_{i}$ for $|i-j|>1$, and
- $s_{i} s_{i+1} s_{i}=s_{i+1} s_{i} s_{i+1}$.

The 0 -Hecke algebra $\mathcal{H}_{n}$ is the unital associative algebra over $\mathbb{C}$ defined by a very similar presentation: $\mathcal{H}_{n}$ is generated by symbols $\sigma_{i}$ (for $1 \leq i<n$ ) subject to

- $\sigma_{i}^{2}=\sigma_{i}$,
- $\sigma_{i} \sigma_{j}=\sigma_{j} \sigma_{i}$ for $|i-j|>1$, and
- $\sigma_{i} \sigma_{i+1} \sigma_{i}=\sigma_{i+1} \sigma_{i} \sigma_{i+1}$.

That is, the $\left\{\sigma_{i}\right\}$ in $\mathcal{H}_{n}$ act exactly like the corresponding $\left\{s_{i}\right\}$ in $\mathcal{S}_{n}$, except that they are idempotent instead of being involutions.

The representation theory of $\mathcal{H}_{n}$ was first worked out in detail by P. Norton [65]. Despite the similarly between the descriptions of $\mathcal{S}_{n}$ and $\mathcal{H}_{n}$, their representation theory is rather different, as $\mathcal{H}_{n}$ is not semisimple. Indeed, the irreducible representations of $\mathcal{H}_{n}$ are all 1-dimensional, while $\mathcal{S}_{n}$ has irreducible representations of higher dimension. The irreducible representations of $\mathcal{H}_{n}$ are equinumerous with the set of compositions $\alpha \vDash(n)$.

There is a quasisymmetric Frobenius character map [18, 43] taking 0-Heckerepresentations to quasisymmetric functions in such a way that the irreducible representations map to the fundamental quasisymmetric functions $F_{\alpha}$. In this way, if the quasisymmetric function $f$ corresponds to the representation $M$, then decomposing $f$ as a sum of fundamental quasisymmetric functions corresponds to identifying the unique direct sum of irreducible 0 -Hecke representations that is equivalent to $M$ in the Grothendieck group of finite-dimensional representations. Certain explicit and combinatorial $\mathcal{H}_{n}$-representations are known whose quasisymmetric Frobenius characters are precisely the quasiSchur functions [79]; unfortunately, these representations are not generally indecomposable.

The geometry of QSym $n$ is much less well understood. In addition to its obvious product structure, QSym also possesses a compatible coproduct, turning it into a Hopf algebra. Although we won't describe it here, there is an important Hopf algebra NSym of noncommutative symmetric functions that is Hopf-dual to QSym. (For background on combinatorial Hopf algebras, see [27].) It is surprisingly easy to see that NSym is isomorphic to the homology of the loop space of the suspension of $\mathbb{C P}^{\infty}$, where the product structure on $H_{\star}\left(\Omega \Sigma \mathbb{C P}^{\infty}\right)$ is given by concatenation of loops [10]. Indeed, this isomorphism even holds on the level of Hopf algebras. Since $\Omega \Sigma \mathbb{C P}{ }^{\infty}$ is an H-space, its homology and cohomology are dual Hopf algebras. (See, for example, [28, 85] for background on H-spaces and their associated Hopf algebras.) From this fact and the fact that QSym is Hopf-dual to NSym, it follows that $H^{\star}\left(\Omega \Sigma \mathbb{C} \mathbb{P}^{\infty}\right)$
is isomorphic to QSym. This interpretation was used in [10] to give cohomological proofs of various properties of QSym; however, it seems that much more could be done from this perspective. A recent construction, which appears closely related, identifies QSym with the Chow ring of an algebraic stack of expanded pairs [66].

## 4 The Asymmetric World

In this section, we consider our fourth and final action of $\mathcal{S}_{n}$ on Poly ${ }_{n}$, the trivial action. Although this action is trivial the associated combinatorics is not at all trivial, but full of rich internal structure and deep connections to geometry and representation theory. The invariant ring of this trivial action is, of course, the ring Poly ${ }_{n}$ itself. However, to emphasize analogies with the previous two sections, we will refer to the invariant ring Poly ${ }_{n}$ in this context as the ring of asymmetric functions $\mathrm{ASym}_{n}$. That is, we write $\mathrm{ASym}_{n}$ for the ring $\mathbb{Z}\left[x_{1}, \ldots, x_{n}\right]$ when we think of it as an $\mathcal{S}_{n}$-module with the trivial $\mathcal{S}_{n}$ action, and we write Poly ${ }_{n}$ for the same ring considered as an object in the category of rings.

Bases of $\mathrm{ASym}_{n}$ are indexed by weak compositions $a$ of length at most $n$, with the most obvious basis of $\mathrm{ASym}_{n}$ being given by individual monomials:

$$
\mathfrak{X}_{a}:=\mathbf{x}^{a} .
$$

Just as $\left\{m_{\lambda}\right\}$ is not the most interesting basis of $\operatorname{Sym}_{n}$, the $\left\{\mathfrak{X}_{a}\right\}$ basis of $\mathrm{ASym}_{n}$ is not very interesting either! Our goal in this section is to explore seven additional bases of rather less trivial nature.

Arguably, the most interesting basis of $\mathrm{ASym}_{n}$ is given by the Schubert polynomials of A. Lascoux and M.-P. Schützenberger [52]. We will use consideration of various formulas for Schubert polynomials to organize our discussion of the bases of Fig.4. Instead of indexing Schubert polynomials by weak compositions, it is more convenient to index them by permutations. Hence, we first recall a standard way to translate between permutations and weak compositions. For a permutation $\pi \in \mathcal{S}_{n}$, let $a_{i}$ denote the number of integers $j>i$ such that $w(i)>w(j)$. (Note that

Fig. 4 The eight bases of $\mathrm{ASym}_{n}$ considered here. The arrows denote that the basis at the head refines the basis at the tail. The star-shaped nodes have positive structure coefficients

$a_{i} \leq n-i$.) The weak composition $a_{\pi}=\left(a_{1}, a_{2}, \ldots, a_{n}\right)$ is called the Lehmer code of $\pi$. Visually, one may determine the Lehmer code of a permutation $\pi$ as follows.

Consider an $n \times n$ grid of boxes and place a laser gun (or dot) in each position $(i, \pi(i))$. Each laser gun fires to the right and down, destroying all boxes directly to its right and all boxes directly below itself (including its own box). The surviving boxes are the Rothe diagram $R D(\pi)$ of the permutation $\pi$. One checks that the Lehmer code of $\pi$ records the number of boxes in each row of $R D(\pi)$.
Example 4.1 Let $\pi=2413$. The Rothe diagram $R D(\pi)$ is shown below

where the surviving boxes are shaded in grey. Hence, the Lehmer code of $\pi=2413$ is $(1,2,0,0)$.

Consider the action of $\mathcal{S}_{n}$ on Poly ${ }_{n}$ from Sect. 2, where permutations act by permuting variables. Now, for each positive integer, define an operator $\partial_{i}$ on Poly ${ }_{n}$ by

$$
\partial_{i}(f):=\frac{f-s_{i} \cdot f}{x_{i}-x_{i+1}}
$$

Note that $\partial_{i}(f)$ is symmetric in the variables $x_{i}$ and $x_{i+1}$. Now, for each permutation $w$ of the form $n(n-1) \cdots 321$ (in one-line notation), the Schubert polynomial $\mathfrak{S}_{w}$ is defined to be

$$
\begin{equation*}
\mathfrak{S}_{w}:=\prod_{i=1}^{n} x_{i}^{n-i}=\mathfrak{X}_{(n-1, n-2, \ldots, 1,0)} \tag{4.1}
\end{equation*}
$$

(These permutations are exactly those that are longest in Coxeter length in $\mathcal{S}_{n}$ for some $n$, i.e., they have the largest possible number of inversions; their Lehmer code is $(n-1, n-2, \ldots, 1,0)$.) For other permutations $w$, the corresponding Schubert polynomials are defined recursively by

$$
\mathfrak{S}_{w}:=\partial_{i} \mathfrak{S}_{w(i, i+1)}
$$

for any $i$ such that $w(i)<w(i+1)$. Amazingly, this recursive definition is selfconsistent, so there is a uniquely defined Schubert polynomial $\mathfrak{S}_{w}$ for each permutation $w$.

Our first task in this section will be to obtain a more concrete understanding of Schubert polynomials by describing how to write them non-recursively in the monomial basis $\left\{\mathfrak{X}_{a}\right\}$. We'll describe three different combinatorial formulas for this expansion, exploring some other families of polynomials along the way.

The first such formula to be proven was given by S. Billey et al. [9]. For a permutation $\pi$, a reduced factorization of $\pi$ is a way of writing $\pi$ as a product $s_{i_{1}} s_{i_{2}} \cdots s_{i_{k}}$
of simple transpositions with $k$ as small as possible. The sequence of subscripts $i_{1} i_{2} \cdots i_{k}$ is called a reduced word for $\pi$. We write $\operatorname{Red}(\pi)$ for the set of all reduced words of the permutation $\pi$. Note that every reduced word $\alpha$ is a strong composition. Given two strong compositions $\alpha$ and $\beta$, we say that $\beta$ is $\alpha$-compatible if
(R.1) $\alpha$ and $\beta$ have the same length,
(R.2) $\beta$ is weakly increasing (i.e., $\beta_{i} \leq \beta_{j}$ for $i<j$ ),
(R.3) $\beta$ is bounded above by $\alpha$ (i.e., $\beta_{i} \leq \alpha_{i}$ for all $i$ ), and
(R.4) $\beta$ strictly increases whenever $\alpha$ does (i.e., if $\alpha_{i}<\alpha_{i+1}$, then $\beta_{i}<\beta_{i+1}$ ).

In this case, we write $\beta \rightarrow \alpha$.
Example 4.2 If $\alpha$ is the strong composition 121 (a reduced word for the longest permutation in $\mathcal{S}_{3}$ ), then we claim that no strong composition is $\alpha$-compatible. Suppose $\beta$ were $\alpha$-compatible. Since $\alpha_{1}<\alpha_{2}$, we must have $\beta_{1}<\beta_{2}$ by (R.4). Hence, $\beta_{2} \geq 2$. Therefore, by (R.2), $\beta_{3} \geq 2$. But this is incompatible with (R.3), since $\alpha_{3}=1$.

On the other hand, for $\gamma=212$ (the other reduced word for this permutation), there is exactly one $\gamma$-compatible strong composition $\delta$. By (R.3), we have $\delta_{2}=1$, and hence by (R.2) we also have $\delta_{1}=1$. By (R.4), $\delta_{3}>\delta_{2}=1$, but by (R.3) $\delta_{3} \leq 2$. Hence, $\delta=112$ is the only $\gamma$-compatible strong composition. We write $112 \rightarrow 212$.

Theorem 4.3 ([9, Theorem 1.1]) The Schubert polynomials expand positively in monomials:

$$
\mathfrak{S}_{\pi}=\sum_{\alpha \in \operatorname{Red}(\pi)} \sum_{\beta \rightarrow \alpha} \prod_{i} x_{\beta_{i}} .
$$

Example 4.4 Let $\pi=321=s_{1} s_{2} s_{1}=s_{2} s_{1} s_{2}$ be the longest permutation in $\mathcal{S}_{3}$. Then, by Example 4.2, we have

$$
\begin{aligned}
\mathfrak{S}_{\pi} & =\sum_{\alpha \in \operatorname{Red}(\pi)} \sum_{\beta \rightarrow \alpha} \prod_{i} x_{\beta_{i}} \\
& =\sum_{\beta \rightarrow 121} \prod_{i} x_{\beta_{i}}+\sum_{\delta \leftrightarrow \rightarrow 212} \prod_{i} x_{\delta_{i}} \\
& =0+x_{1} x_{1} x_{2}=x_{1}^{2} x_{2}=\mathfrak{X}_{(2,1,0)} .
\end{aligned}
$$

Note that this calculation is consistent with the definition given in Eq.(4.1).
It might be reasonable to expect an important basis of $\mathrm{ASym}_{n}$ to restrict to an important basis of the subspace $\operatorname{Sym}_{n} \subset \mathrm{ASym}_{n}$. Indeed, one piece of evidence for the importance of Schubert polynomials is that those Schubert polynomials lying inside $\mathrm{Sym}_{n}$ are exactly its basis of Schur polynomials.

Definition 4.5 Suppose we have vector spaces $U \subseteq V$ with bases $B_{U}$ and $B_{V}$, respectively. If $B_{V} \cap U=B_{U}$, we say the basis $B_{V}$ lifts the basis $B_{U}$ to $V$, or equivalently that $B_{V}$ is a lift of $B_{U}$.

Since every Schur polynomial is a Schubert polynomial, in this language the Schubert basis of $\mathrm{ASym}_{n}$ is a lift of the Schur basis of $\mathrm{Sym}_{n}$. To realize the Schur polynomial $s_{\lambda} \in \operatorname{Sym}_{n}$ as a Schubert polynomial, first realize the partition $\lambda$ as a weak composition of length $n$ by padding it by an appropriate number of final 0 s. Now, reverse the letters of $\lambda$, so it becomes a weakly increasing sequence. The resulting weak composition is the Lehmer code of a unique permutation $\pi_{\lambda}$, and one has $s_{\lambda}=\mathfrak{S}_{\pi_{\lambda}}$. Equivalently, for $i \leq n$ one has $\pi_{\lambda}(i)=\lambda_{n-i+1}+i$ and for $i>n$ one has $\pi_{\lambda}(i)=\min \left(\mathbb{Z}_{>0} \backslash\left\{\pi_{\lambda}(j): j<i\right\}\right)$.

Since the Schubert polynomials lift the Schur polynomials, one might wonder whether the Littlewood-Richardson rule (Theorem 2.5) also lifts to a positive combinatorial rule for the structure coefficients of the Schubert basis. Indeed, the Schubert basis of $\mathrm{ASym}_{n}$, like the Schur basis of $\mathrm{Sym}_{n}$, has positive structure coefficients! However, no combinatorial proof of this fact is known and we lack any sort of positive combinatorial rule (even conjectural) to describe these structure coefficients (except in a few very special cases, such as when the Schubert polynomials are actually Schur polynomials). Discovering and proving such a rule is one of the most important open problems in algebraic combinatorics. Part of our motivation for studying the combinatorial theory of $\mathrm{ASym}_{n}$ is the hope that such a theory will eventually lead to a Schubert structure coefficient rule, just as the Littlewood-Richardson rule for Schur polynomial structure coefficients eventually developed from the combinatorial theory of $\mathrm{Sym}_{n}$.

Without such a combinatorial rule, how then do we know that the Schubert basis has positive structure coefficients? The answer comes, once again, from geometry and from representation theory. Geometrically, instead of looking at a complex Grassmannian, as we did for Schur polynomials, we should consider a complex flag variety Flags ${ }_{n}$, the classifying space for complete flags $V_{0} \subset V_{1} \subset \cdots \subset V_{n}$ of nested complex vector bundles with $V_{k}$ of rank $k$. This space has an analogous cell decomposition by Schubert varieties, yielding an effective basis of the Chow ring. By identifying Schubert varieties with corresponding Schubert polynomials, multiplying Schubert polynomials corresponds to the intersection product on Schubert varieties and positivity of structure coefficients follows. An alternative proof of positivity [82, 83] is given by interpreting Schubert polynomials as characters of certain Kraśkiewicz-Pragacz modules (introduced in [41, 42]) for Borel Lie algebras.

The formula of Theorem 4.3 naturally leads us to consider another family of polynomials. Suppose we fix a reduced word $\alpha \in \operatorname{Red}(\pi)$ for some $\pi \in \mathcal{S}_{n}$. Then, Theorem 4.3 suggests defining a polynomial

$$
\begin{equation*}
\mathfrak{F}(\alpha):=\sum_{\beta \leftrightarrow \alpha} \prod_{i} x_{\beta_{i}}, \tag{4.2}
\end{equation*}
$$

so that Theorem 4.3 may be rewritten as

$$
\mathfrak{S}_{\pi}=\sum_{\alpha \in \operatorname{Red}(\pi)} \mathfrak{F}(\alpha) .
$$

Indeed, the formula of Eq.(4.2) makes sense for any strong composition $\alpha$, not necessarily a reduced word of a permutation.

Labeling these polynomials by strong compositions $\alpha$ is unnatural for at least two reasons. For some $\alpha$, we have $\mathfrak{F}(\alpha)=0$; for example, we have $\mathfrak{F}(1,2,1)=0$ by Example 4.2. Those $\mathfrak{F}(\alpha)$ that are nonzero are called the fundamental slide polynomials; these were introduced in [2], although the alternate definition we give here follows [5]. Also, for $\alpha \neq \alpha^{\prime}$, we can have $\mathfrak{F}(\alpha)=\mathfrak{F}\left(\alpha^{\prime}\right) \neq 0$; for example, by Example 4.2 we have $\mathfrak{F}(2,1,2)=\mathfrak{X}_{(2,1,0)}$, but it is also equally clear that $\mathfrak{F}(3,1,2)=$ $\mathfrak{X}_{(2,1,0)}$.

For any strong composition $\alpha$, note that, if $\alpha$ has any compatible sequences, then it has a unique such compatible sequence $\beta(\alpha)$ that is termwise maximal. Define the weak composition $a(\alpha)$ by letting $a(\alpha)_{i}$ denote the multiplicity of $i$ in $\beta(\alpha)$. Then, we define

$$
\mathfrak{F}_{a(\alpha)}:=\mathfrak{F}(\alpha) .
$$

It is clear then that every fundamental slide polynomial is, in this fashion, uniquely indexed by a weak composition $a$. Moreover, every weak composition $a$ appears as an index on some $\mathfrak{F}_{a}$, and we have $\mathfrak{F}_{a} \neq \mathfrak{F}_{b}$ if $a \neq b$. It is then not hard to see by triangularity in the $\mathfrak{X}_{a}$ basis that the set of fundamental slide polynomials forms another basis of ASym $n$.

Clearly, the fundamental slide polynomials expand positively in the monomial basis $\left\{\mathfrak{X}_{a}\right\}$. It is useful to have a formula for this expansion of $\mathfrak{F}_{a}$, based only on the weak composition $a$. We first need a partial order on weak compositions: we write $a \geq b$ and say $a$ dominates $b$ if we have

$$
\sum_{i=1}^{k} a_{i} \geq \sum_{i=1}^{k} b_{i}
$$

for all $k$. (Note that the restriction of this partial order to the set of partitions recovers the usual notion of dominance order.)

Theorem 4.6 ([2, 5]) The fundamental slide polynomials expand positively in monomials:

$$
\mathfrak{F}_{a}=\sum_{\substack{b \geq a \\ b^{+} \sum_{a}^{+}}} \mathfrak{X}_{b}
$$

Essentially by definition, the fundamental slide polynomials are pieces of Schubert polynomials. Although the basis of Schubert polynomials has positive structure constants, there is no reason to expect this property to descend to this basis of pieces. Remarkably, however, the fundamental slide polynomials also have positive structure constants! The first clue that this might be the case comes from considering the intersection of the fundamental slide basis with the subring QSym $_{n} \subset \mathrm{ASym}_{n}$.

Theorem 4.7 ([2]) We have $\mathfrak{F}_{a} \in \mathrm{QSym}_{n}$ if and only if a is of the form $0^{k} \alpha$, where $\alpha$ is a strong composition of length $n-k$ and $0^{k} \alpha$ denotes the weak composition obtained from $\alpha$ by prepending $k 0 s$.

Moreover, we have $\mathfrak{F}_{0^{k} \alpha}=F_{\alpha}$. Thus, the fundamental slide basis of $\mathrm{ASym}_{n}$ is a lift of the fundamental quasisymmetric polynomial basis of $\mathrm{QSym}_{n}$.

In light of Theorem 4.7, one might hope to extend the combinatorial multiplication rule of Theorem 3.8 to fundamental slide polynomials. Indeed, this is possible. We need to extend the notion of the shuffle product of two strong compositions from Sect. 3 to the slide product or pairs of weak compositions $a, b$. Here, we borrow notation from [71]. As before, let $\mathcal{A}$ denote the alphabet of odd integers and let $\mathcal{B}$ denote the alphabet of even integers. Let $A$ be the word in $\mathcal{A}$ consisting of $a_{1}$ copies of $2 \ell(a)-1$, followed by $a_{2}$ copies of $2 \ell(a)-3$, all the way to $a_{\ell(a)}$ copies of 1 . Likewise, let $B$ denote the word in $\mathcal{B}$ consisting of $b_{1}$ copies of $2 \ell(b)$, followed by $b_{2}$ copies of $2 \ell(b)-2$, all the way to $b_{\ell(b)}$ copies of 2 .

For any word $W$ in a totally ordered alphabet $\mathcal{Z}$, let $\operatorname{Runs}(W)$ denote the sequence of successive maximally weakly increasing runs of letters of $W$ read from left to right. For a sequence $S$ of words in $\mathcal{Z}$ and any subalphabet $\mathcal{Y} \subseteq \mathcal{Z}$, write $\operatorname{Comp}_{\mathcal{Y}}(S)$ for the weak composition whose $i$ th coordinate is the number of letters of $\mathcal{Y}$ in the $i$ th word of $S$.

Let $\operatorname{Sh}(a, b)$ denote the set of those shuffles $C$ of $A$ and $B$ such that

$$
\operatorname{Comp}_{\mathcal{A}}(\operatorname{Runs}(C)) \geq a \text { and } \operatorname{Comp}_{\mathcal{B}}(\operatorname{Runs}(C)) \geq b .
$$

For $C \in \operatorname{Sh}(a, b)$, let BumpRuns( $C$ ) denote the unique dominance-minimal way to insert words of length 0 into $\operatorname{Runs}(C)$ while preserving $\operatorname{Comp}_{\mathcal{A}}(\operatorname{BumpRuns}(C)) \geq a$ and $\operatorname{Comp}_{\mathcal{B}}(\operatorname{BumpRuns}(C)) \geq b$. Finally, define the slide product $a \amalg b$ of the weak compositions $a$ and $b$ as the formal sum of weak compositions

$$
a \amalg b:=\sum_{C \in \operatorname{Sh}(a, b)} \operatorname{Comp}_{\mathbb{Z}}(\operatorname{BumpRuns}(C)) .
$$

Example 4.8 Let $a=(0,1,0,2)$ and $b=(1,0,0,1)$. Then we consider the words $A=511$ and $B=82$. The set of all shuffles of $A$ and $B$ is
$\{51182,51812,58112,85112,51821,58121,85121,58211,85211,82511\}$.

Many of these shuffles $C$ fail $\operatorname{Comp}_{\mathcal{B}}(\operatorname{Runs}(C)) \geq b$; for example, with $C=51821$, we have $\operatorname{Runs}(C)=(5,18,2,1)$ and hence $\operatorname{Comp}_{\mathcal{B}}(\operatorname{Runs}(C))=(0,1,1,0) \nsupseteq b$. Thus we have

$$
\operatorname{Sh}(a, b)=\{58112,85112,58121,85121,58211,85211,82511\} .
$$

The corresponding $\operatorname{BumpRuns}(C)$ for $C \in \operatorname{Sh}(a, b)$ are
where $\epsilon$ denotes the empty word. Thus, we have

$$
\begin{aligned}
(0,1,0,2) \amalg(1,0,0,1) & =(2,0,0,3)+(1,1,0,3)+(2,0,2,1)+(1,1,2,1)+(2,0,1,2) \\
& +(1,1,1,2)+(1,2,0,2) .
\end{aligned}
$$

Finally, we can state the multiplication rule for fundamental slide polynomials.
Theorem 4.9 ([2]) For weak compositions $a$ and $b$, we have

$$
\mathfrak{F}_{a} \cdot \mathfrak{F}_{b}=\sum_{c} C_{a, b}^{c} \mathfrak{F}_{c},
$$

where $C_{a, b}^{c}$ is the multiplicity of $c$ in the slide product $a \amalg b$.
It seems reasonable to expect that the positivity of Theorem 4.9 reflects some geometry or representation theory governed by fundamental slide polynomials. Sadly, no such interpretation of fundamental slide polynomials is known, except in the quasisymmetric case.

A second formula for Schubert polynomials uses the combinatorial model of pipe dreams. This model was successively developed in [7, 20, 21, 39]. A pipe dream $P$ is a tiling of the grid of boxes (extending infinitely to the east and south) with turning pipes $r$ and finitely many crossing pipes + . Such a tiling gives rise to a collection of lines called pipes, which one imagines traveling from the left side of the grid (the negative $y$-axis) to the top side (the positive $x$-axis). A pipe dream $P$ is called reduced if no two pipes cross each other more than once. The permutation corresponding to a reduced pipe dream is the permutation given (in one-line notation) by the columns in which the pipes end. The weight $\mathrm{wt}(P)$ of a pipe dream is the weak composition whose $i$ th entry is the number of crossing pipe tiles in row $i$ of $P$ (where row 1 is the top row).

Example 4.10 The pipe dream

corresponds to the permutation 1432. (Here, we omit the infinite collection of ${ }^{\wedge}$ rtiles extending uninterestingly to the southeast.) The pipe dream $P$ has weight ( $1,2,0$ ).

Given a permutation $\pi$, let $\operatorname{PD}(\pi)$ denote the set of reduced pipe dreams for $\pi$.
Theorem 4.11 ([7,9]) The Schubert polynomial $\mathfrak{S}_{\pi}$ is the generating function of reduced pipe dreams for $\pi$, i.e.,


Fig. 5 The 7 reduced pipe dreams associated to the permutation 15324

$$
\mathfrak{S}_{\pi}=\sum_{P \in \operatorname{PD}(\pi)} \mathbf{x}^{\mathrm{wt}(P)}
$$

Example 4.12 We have $\mathfrak{S}_{15324}=\mathbf{x}^{031}+\mathbf{x}^{121}+\mathbf{x}^{211}+\mathbf{x}^{310}+\mathbf{x}^{310}+\mathbf{x}^{130}+\mathbf{x}^{220}$, where the monomials are determined by the pipe dreams shown in Fig. 5.

A reduced pipe dream $P$ is quasiYamanouchi if the following is true for the leftmost + in every row: Either
(1) it is in the leftmost column, or
(2) it is weakly left of some + in the row below it.

For a permutation $\pi$, write $\mathrm{QPD}(\pi)$ for the set of quasiYamanouchi reduced pipe dreams for $\pi$. We obtain then the following formula for the fundamental slide polynomial expansion of a Schubert polynomial.

Theorem 4.13 ([2]) The Schubert polynomials expand positively in the fundamental slide polynomials:

$$
\mathfrak{S}_{\pi}=\sum_{P \in \mathrm{QPD}(\pi)} \mathfrak{F}_{\mathrm{wt}(P)}
$$

Before continuing to our third combinatorial formula for Schubert polynomials, let us consider another basis of $\mathrm{ASym}_{n}$, closely related to the fundamental slide polynomials. Recall from Theorem 4.7 that the fundamental slide polynomials are a lift of the fundamental quasisymmetric polynomials. One might ask for an analogous lift to $\mathrm{ASym}_{n}$ of the monomial quasisymmetric polynomials. These are provided by the monomial slide polynomials of [2], which we now discuss.

Looking back at the combinatorial formulas for fundamental and monomial quasisymmetric polynomials, observe that they are identical, except that the formula for $F_{\alpha}$ looks at weak compositions $b$ with $b^{+} \vDash \alpha$ while the formula for $M_{\alpha}$ looks at weak compositions $b$ with the more restrictive property $b^{+}=\alpha$. It is easy then to
guess the following modification of Theorem 4.6 that will yield the desired definition of monomial slide polynomials.

Definition 4.14 For any weak composition $a$, the monomial slide polynomial $\mathfrak{M}_{a}$ is defined by

$$
\mathfrak{M}_{a}=\sum_{\substack{b \geq a \\ b^{+}=a^{+}}} \mathfrak{X}_{b} .
$$

By triangularity, it is straightforward that monomial slide polynomials form a basis of $\mathrm{ASym}_{n}$. Moreover, we have the following analogue of Theorem 4.7:

Theorem 4.15 ([2]) We have $\mathfrak{M}_{a} \in \operatorname{QSym}_{n}$ if and only if a is of the form $0^{k} \alpha$, where $\alpha$ is a strong composition of length $n-k$.

Moreover, we have $\mathfrak{M}_{0^{k} \alpha}=M_{\alpha}$. Thus, the monomial slide basis of $\mathrm{ASym}_{n}$ is a lift of the monomial quasisymmetric polynomial basis of $\mathrm{QSym}_{n}$.

Just as the combinatorial multiplication rule of Theorem 3.8 for fundamental quasisymmetric polynomials lifts to that of Theorem 4.9 for fundamental slide polynomials, the combinatorial multiplication rule of Theorem 3.3 for monomial quasisymmetric polynomials lifts to a rule for monomial slide polynomials.

First, we need to extend the overlapping shuffle product of Sect. 3 from strong compositions to general weak compositions. Given weak compositions $a$ and $b$, treat them as words of some common finite length $n$ by truncating at some position past all their nonzero entries. By $a+b$ we mean the weak composition that is the coordinatewise sum of $a$ and $b$. Consider the set $S(a, b)$ of all pairs $\left(a^{\prime}, b^{\prime}\right)$ of weak compositions of equal length $k \leq n$ such that

- $\left(a^{\prime}\right)^{+}=a^{+}$and $\left(b^{\prime}\right)^{+}=b^{+}$;
- $a^{\prime} \geq a$ and $b^{\prime} \geq b$; and
- for all $1 \leq i \leq k$, we have $a_{i}^{\prime}+b_{i}^{\prime}>0$.

Fix $\left(a^{\prime}, b^{\prime}\right) \in S(a, b)$. Let $c$ be a weak composition of length $r$ with zeros in positions $z_{1}, \ldots, z_{m}$ such that $c^{+}=a^{\prime}+b^{\prime}$. Define $c_{a}$ to be the weak composition of length $r$ having zeros in the same positions $z_{1}, \ldots, z_{m}$ and the remaining positions of $c_{a}$ are the entries of $a^{\prime}$, in order from left to right. Define $c_{b}$ similarly, using the entries of $b^{\prime}$. Then we have

- $c=c_{a}+c_{b}$, and
- $\left(c_{a}\right)^{+}=\left(a^{\prime}\right)^{+}$and $\left(c_{b}\right)^{+}=\left(b^{\prime}\right)^{+}$.

For each such $\left(a^{\prime}, b^{\prime}\right) \in S(a, b)$, let $\operatorname{Bump}\left(a^{\prime}, b^{\prime}\right)$ denote the unique dominance-least weak composition satisfying

- $\operatorname{Bump}\left(a^{\prime}, b^{\prime}\right)^{+}=a^{\prime}+b^{\prime}$, and
- $\operatorname{Bump}\left(a^{\prime}, b^{\prime}\right)_{a} \geq a$ and $\operatorname{Bump}\left(a^{\prime}, b^{\prime}\right)_{b} \geq b$.

The overlapping slide product of $a$ and $b$ is then the formal sum $a \amalg_{o} b$ of the $\operatorname{Bump}\left(a^{\prime}, b^{\prime}\right)$ for all $\left(a^{\prime}, b^{\prime}\right) \in S(a, b)$.

Example 4.16 Let $a=(0,1,0,2)$ and $b=(1,0,0,1)$, as in Example 4.8. Then $S(a, b)$ consists of the seven pairs

$$
\begin{aligned}
& ((0,1,0,2),(1,0,1,0)),((0,1,2,0),(1,0,0,1)),((0,1,2),(1,0,1)),((0,1,2),(1,1,0)) \\
& ((1,0,2),(1,1,0)),((1,2,0),(1,0,1)),((1,2),(1,1))
\end{aligned}
$$

The seven corresponding weak compositions $\operatorname{Bump}\left(a^{\prime}, b^{\prime}\right)$ are

$$
(1,1,1,2),(1,1,2,1),(1,1,0,3),(1,2,0,2),(2,0,1,2),(2,0,2,1),(2,0,0,3)
$$

Theorem 4.17 ([2]) For weak compositions $a$ and $b$, we have

$$
\mathfrak{M}_{a} \cdot \mathfrak{M}_{b}=\sum_{c} C_{a, b}^{c} \mathfrak{M}_{c},
$$

where $C_{a, b}^{c}$ is the multiplicity of $c$ in the overlapping slide product $a Ш_{o} b$.
The fundamental slide polynomials expand positively in the monomial slide basis. Say that $b \unrhd a$ if $b \geq a$, and $c \geq b$ whenever $c \geq a$ and $c^{+}=b^{+}$.

Theorem 4.18 [2] The fundamental slide polynomials expand positively in the monomial slide polynomials:

$$
\mathfrak{F}_{a}=\sum_{\substack{b \unrhd a \\ b^{+}=a^{+}}} \mathfrak{M}_{b} .
$$

A third combinatorial formula for Schubert polynomials comes from a model introduced (conjecturally) by Axel Kohnert [40]. Let $D$ be a box diagram, i.e., any subset of the boxes in an $n \times n$ grid. A Kohnert move on $D$ selects the rightmost box in some row and moves it to the first available empty space above it in the same column (if such an empty space exists). Let $\mathrm{KD}(D)$ denote the set of all box diagrams that can be obtained from $D$ by some sequence (possibly empty) of Kohnert moves. Define the weight $\mathrm{wt}(D)$ of a box diagram to be the weak composition where $\mathrm{wt}(D)_{i}$ records the number of boxes in the $i$ th row of $D$ from the top.

Example 4.19 Let $D$ be the leftmost diagram in the top row of Fig. 6 (which happens to be the Rothe diagram $R D(15324)$ ). The set of diagrams in Fig. 6 is exactly $\operatorname{KD}(D)$.

Theorem 4.20 ([4, 86, 87]) The Schubert polynomial $\mathfrak{S}_{\pi}$ is the generating function of the Kohnert diagrams for the Rothe diagram of $\pi$, i.e.,

$$
\mathfrak{S}_{\pi}=\sum_{D \in \operatorname{KD}(\operatorname{RD}(\pi))} \mathbf{x}^{\mathrm{wt}(D)}
$$



Fig. 6 The 7 Kohnert diagrams associated to $R D$ (15324)

Example 4.21 We have $\mathfrak{S}_{15324}=\mathbf{x}^{031}+\mathbf{x}^{121}+\mathbf{x}^{211}+\mathbf{x}^{310}+\mathbf{x}^{310}+\mathbf{x}^{130}+\mathbf{x}^{220}$, where the monomials are determined by the Kohnert diagrams shown in Fig. 6. Note this is consistent with the computation in Example 4.12.

Kohnert diagrams also yield another natural basis of $\mathrm{ASym}_{n}$, which we now consider. Given a weak composition $a$, we can associate a box diagram $D$ by first obtaining the permutation $w$ corresponding to $a$ and then taking $D=\mathrm{RD}(w)$. This construction, combined with the Kohnert moves, leads us to the characterization of Schubert polynomials from Theorem 4.20. However, there is also a much easier way to associate a box diagram $D(a)$ to a weak composition $a$-namely, just take $D(a)$ to be the Young diagram of $a$, as described in Sect.3. This leads us to the following definition (really a theorem of Kohnert [40]): For a weak composition $a$, the key polynomial $\mathfrak{D}_{a}$ is the generating function

$$
\mathfrak{D}_{a}:=\sum_{D \in K D(D(a))} \mathbf{x}^{\mathrm{wt}(D)} .
$$

Example 4.22 Let $a=(0,2,1)$. Then $\mathfrak{D}_{a}=\mathbf{x}^{021}+\mathbf{x}^{111}+\mathbf{x}^{201}+\mathbf{x}^{210}+\mathbf{x}^{120}$, as computed by the Kohnert diagrams in Fig. 7.

The definition of key polynomials that we have given here is not the original one. These polynomials were first introduced in [17] where they were realized as characters of (type A) Demazure modules; for this reason, they are often referred to as


Fig. 7 The 5 Kohnert diagrams associated to the weak composition ( $0,2,1$ )

Demazure characters. Later they were studied from a more combinatorial perspective by Lascoux and Schützenberger [54], who coined the term 'key polynomial'. Key polynomials also arise as a specialization [34, 76] of the nonsymmetric Macdonald polynomials introduced in [15, 57, 67].

An alternative description of key polynomials is via a modification of the $\partial_{i}$ operators that define Schubert polynomials. For each positive integer, define an operator $\pi_{i}$ on Poly ${ }_{n}$ by

$$
\pi_{i}(f):=\partial_{i}\left(x_{i} f\right)
$$

Let $w$ be a permutation and let $s_{i_{1}} \cdots s_{i_{r}}$ be any reduced word for $w$. Then define $\pi_{w}=$ $\pi_{i_{1}} \cdots \pi_{i_{r}}$. This is independent of the choice of reduced word since these operators satisfy $\pi_{i}^{2}=\pi_{i}$ and the usual commutation and braid relations for the symmetric group. (That is to say, the action of the $\pi_{i}{\text { operators on } \text { Poly }_{n} \text { is a representation }}^{\text {a }}$ of the type A 0 -Hecke algebra.) Let $s_{i}$ act on a weak composition $a$ by exchanging the $i$ th and $(i+1)$ st entries of $a$. Given a weak composition $a$, let $w(a)$ denote the permutation of minimal Coxeter length such that $w(a) \cdot a=\overleftarrow{a}$. Finally, the key polynomial $\mathfrak{D}_{a}$ is given by

$$
\mathfrak{D}_{a}=\pi_{w(a)} \mathbf{x}^{\overleftarrow{a}}
$$

## Example 4.23

$$
\begin{aligned}
\mathfrak{D}_{(0,2,1)} & =\pi_{1} \pi_{2}\left(x_{1}^{2} x_{2}\right) \\
& =\pi_{1}\left(x_{1}^{2} x_{2}+x_{1}^{2} x_{3}\right) \\
& =x_{1}^{2} x_{2}+x_{1} x_{2}^{2}+x_{1}^{2} x_{3}+x_{1} x_{2} x_{3}+x_{2}^{2} x_{3} .
\end{aligned}
$$

Compare this calculation with that of Example 4.22.
As is the case for Schubert polynomials, there are several additional combinatorial formulas for key polynomials! An excellent overview can be found in [74]. Another such combinatorial formula for the key polynomials is given in [31] in terms of semiskyline fillings. Let $a$ be a weak composition of length $n$. We recall the definition of a triple of entries from the previous section; this extends verbatim from diagrams of compositions to diagrams of weak compositions.

A semi-skyline filling of $D(a)$ is a filling of the boxes of $D(a)$ with positive integers, one per box, such that
(S.1) entries do not repeat in a column,
(S.2) entries weakly decrease from left to right along rows, and
(S.3) every triple of entries is inversion.

The weight $\mathrm{wt}(T)$ of a semi-skyline filling is weak composition whose $i$ th entry records the number of occurrences of the entry $i$ in $T$.

Let $\bar{D}(a)$ denote the diagram of $D(a)$ augmented with an additional 0th column called a basement, and let $\operatorname{rev}(a)$ denote the weak composition obtained by reading the entries of $a$ in reverse. Define a key semi-skyline filling for $a$ to be a filling of


Fig. 8 The five key semi-skyline fillings associated to 021
$\bar{D}(\operatorname{rev}(a))$, where the $i$ th basement entry is $n+1-i$, satisfying (S.1), (S.2) and (S.3) above (including on basement entries). Let $\mathfrak{D S S T}(a)$ denote the set of key semi-skyline fillings for $a$.

Example 4.24 Let $a=(0,2,1)$. The key semi-skyline fillings associated to $a$ are shown in Fig. 8. The basement boxes are shaded in grey.

Theorem 4.25 ([31]) The key polynomial $\mathfrak{D}_{a}$ is given by

$$
\mathfrak{D}_{a}=\sum_{T \in \mathfrak{D S S T}(a)} \mathbf{x}^{\mathrm{wt}(T)}
$$

Example 4.26 From Example 4.24, we again compute $\mathfrak{D}_{(0,2,1)}=\mathbf{x}^{021}+\mathbf{x}^{111}+$ $\mathbf{x}^{201}+\mathbf{x}^{210}+\mathbf{x}^{120}$.

Yet another formula for key polynomials is given in terms of Kohnert tableaux [3]. Kohnert tableaux associate a canonical path in Kohnert's algorithm from the diagram of a weak composition to a given Kohnert diagram. In fact, the Kohnert tableaux are equivalent to the key semi-skyline fillings turned upside down, but are described by quite different local rules, which arise from Kohnert's algorithm as opposed to considerations in Macdonald polynomial theory. We omit the details of this construction, for which see [3].

One might naturally ask, analogously for Schubert polynomials, whether key polynomials restrict to an important basis of $\operatorname{Sym}_{n} \subset \mathrm{ASym}_{n}$. Exactly the same as for the Schubert polynomials, the key polynomials in $\mathrm{Sym}_{n}$ are exactly the basis of Schur polynomials. Thus the key basis of $\mathrm{ASym}_{n}$ is also a lift of the Schur basis of $\operatorname{Sym}_{n}$. The key polynomial $\mathfrak{D}_{a}$ is a Schur polynomial if and only if $a$ is weakly increasing; in this case, we have $\mathfrak{D}_{a}=s_{\overleftarrow{a}}$.

Moreover, the stable limits of the key polynomials are exactly the Schur functions. Given $m \in \mathbb{Z}_{\geq 0}$ and a weak composition $a$, recall that $0^{m} a$ denotes the weak composition obtained by prepending $m$ zeros to $a$, e.g., $0^{2}(1,0,3)=(0,0,1,0,3)$. The stable limit of $\mathfrak{D}_{a}$ is the formal power series $\lim _{m \rightarrow \infty} \mathfrak{D}_{0^{m} a}$. The following result is implicit in work of Lascoux and Schützenberger [54]; an explicit proof is given in [3] with further details.

Theorem 4.27 ([3,54]) Let a be a weak composition. Then the stable limit of the key polynomial $\mathfrak{D}_{a}$ is the Schur function associated to the partition obtained by rearranging the entries of a into decreasing order. That is,

$$
\lim _{m \rightarrow \infty} \mathfrak{D}_{0^{m} a}=s_{\overleftarrow{a}}(X)
$$

Unlike the Schubert basis, however, the key basis does not have positive structure constants.

Schubert polynomials expand in the key basis with positive coefficients, although we omit the details of this decomposition. Let $T$ be a semistandard Young tableau. Define the column reading word colword $(T)$ of $T$ to be the word obtained by writing the entries of each column of $T$ from bottom to top, starting with the leftmost column and proceeding rightwards. Then, as given in [74, Theorem 4], we have

$$
\mathfrak{S}_{\pi}=\sum_{\operatorname{colword}(T) \in \operatorname{Red}\left(\pi^{-1}\right)} \mathfrak{D}_{\mathrm{wt}\left(K_{-}^{0}(T)\right)}
$$

where the sum is over all semistandard Young tableaux $T$ whose column reading word is a reduced word for $\pi^{-1}$, and $K_{-}^{0}(T)$ is a semistandard Young tableau of the same shape as $T$ called the left nil key of $T$, as defined in [54, 74]. For another approach to this decomposition, see [6].

The Demazure atoms $\mathfrak{A}_{a}$ form another basis of $\mathrm{ASym}_{n}$, introduced and studied in [54], where they are referred to as standard bases. Demazure atoms are characters of quotients of Demazure modules, and, like key polynomials, also arise as specializations of nonsymmetric Macdonald polynomials [60]. Just as we defined fundamental slide polynomials as the pieces of Schubert polynomials given by the summands of Theorem 4.3, we can define the Demazure atoms as the pieces of quasiSchur polynomials given by the summands of the definition in Eq. (3.1): Given a weak composition $a$, the Demazure atom $\mathfrak{A}_{a}$ is given by

$$
\mathfrak{A}_{a}:=\sum_{T \in \mathfrak{A} S S T(a)} \mathbf{x}^{\mathrm{wt}(T)} .
$$

Example 4.28 We have

$$
\mathfrak{A}_{(1,0,3)}=\mathbf{x}^{103}+\mathbf{x}^{112}+\mathbf{x}^{202}+\mathbf{x}^{121}+\mathbf{x}^{211}
$$

where the monomials are determined by the five semistandard composition tableaux of shape (1, 0, 3) from Fig. 3.

The Demazure atom basis does not have positive structure coefficients. However, it does exhibit a variety of surprising positivity properties. First, notice that the definition that we have given immediately implies that the quasiSchur polynomial $S_{\alpha}\left(x_{1}, \ldots, x_{n}\right)$ expands positively in Demazure atoms.

Proposition 4.29 ([30]) The quasiSchur polynomials expand positively in the Demazure atoms:

$$
S_{\alpha}\left(x_{1}, \ldots, x_{n}\right)=\sum_{a^{+}=\alpha} \mathfrak{A}_{a}
$$

where the sum is over weak compositions a of length $n$.
It is also the case that key polynomials expand positively in Demazure atoms. Let $S_{n}$ act on weak compositions of length $n$ via $v \cdot\left(a_{1}, \ldots, a_{n}\right)=\left(a_{v^{-1}(1)}, \ldots, a_{v^{-1}(n)}\right)$. Given a weak composition $a$, let $v(a)$ denote the permutation of minimal Coxeter length such that $v(a) \cdot a=\overleftarrow{a}$.

Theorem 4.30 ([54]) The key polynomials expand positively in the Demazure atoms:

$$
\mathfrak{D}_{a}=\sum_{\substack{v(b) \leq v(a) \\ \overleftarrow{b}=\overleftarrow{a}}} \mathfrak{A}_{b},
$$

where $\leq$ denotes the (strong) Bruhat order on permutations.
Since every Schur polynomial is also a key polynomial, Theorem 4.30 implies that Schur polynomials expand positively in the Demazure atoms.

Example 4.31 Let $a=(1,0,3)$. Then $v(a)=231$ and

$$
\mathfrak{D}_{(1,0,3)}=\mathfrak{A}_{(1,0,3)}+\mathfrak{A}_{(1,3,0)}+\mathfrak{A}_{(3,0,1)}+\mathfrak{A}_{(3,1,0)} .
$$

Finally, we mention the following remarkable conjecture of V. Reiner and M. Shimozono; for more details on this conjecture, see the work of A. Pun [72]. For a generalization, see [64]. Observe that the conjecture would follow trivially from Theorem 4.30 if either the key polynomial or the Demazure atom basis had positive structure coefficients; however, neither does, so the conjecture is quite mysterious.

Conjecture 4.32 (Reiner-Shimozono) The product $\mathfrak{D}_{a} \cdot \mathfrak{D}_{b}$ expands positively in Demazure atoms.

At this point, we have considered lifts to $\mathrm{ASym}_{n}$ of the Schur polynomials (two distinct lifts even), the fundamental quasisymmetric polynomials, and the monomial quasisymmetric polynomials. It is natural then to hope for an appropriate lift of the remaining basis of $\mathrm{QSym}_{n}$ that we considered in Sect. 3, namely the quasiSchur polynomials. The next basis we consider is exactly this desired lift, the quasikey polynomials of [3]. The quasikey polynomials are a lifting of the quasiSchur basis of $\mathrm{QSym}_{n}$ to $\mathrm{ASym}_{n}$, and simultaneously a common coarsening of the fundamental slide polynomial and Demazure atom bases.

Let $a$ be a weak composition of length $n$. The quasikey polynomial associated to $a$ is given by

$$
\mathfrak{Q}_{a}:=\sum_{\substack{b^{+}=a^{+} \\ b \geq a}} \sum_{T \in \mathfrak{A S S T}(a)} \mathbf{x}^{\operatorname{wt}(T)}
$$

where the first sum is over all weak compositions $b$ of length $n$ satisfying $b \geq a$ in dominance order and whose positive part is $a^{+}$. The form of this definition is due
to [77]. The quasikey polynomials were originally defined in [3] as a weighted sum of quasi-Kohnert tableaux; we omit this alternate formulation.

Example 4.33 Let $a=(1,0,3)$. Then

$$
\mathfrak{Q}_{(1,0,3)}=\mathbf{x}^{130}+\mathbf{x}^{220}+\mathbf{x}^{103}+\mathbf{x}^{112}+\mathbf{x}^{202}+\mathbf{x}^{121}+\mathbf{x}^{211}
$$

where the monomials are determined by the first seven semistandard composition tableaux shown in Fig. 3.

The definition we have given here immediately implies that each quasikey polynomial $\mathfrak{Q}_{a}$ expands positively in Demazure atoms. Specifically, we have the following.

Theorem 4.34 ([77, Theorem 3.4]) The quasikey polynomials expand positively in the Demazure atoms: For a weak composition a of length n, we have

$$
\mathfrak{Q}_{a}=\sum_{\substack{b^{+}=a^{+} \\ b \geq a}} \mathfrak{A}_{b},
$$

where the sum is over all weak compositions $b$ of length $n$ satisfying $b \geq a$ in dominance order and whose positive part is $a^{+}$.

Example 4.35 One easily calculates from Theorem 4.34 that

$$
\mathfrak{Q}_{(1,0,3)}=\mathfrak{A}_{(1,0,3)}+\mathfrak{A}_{(1,3,0)}
$$

Compare this calculation to the tableaux of Fig. 3.
Less clear from our definition is the following additional positivity property of quasikey polynomials.

Theorem 4.36 ([2]) The quasikey polynomials expand positively in the fundamental slide polynomials: For a weak composition a of length $n$, we have

$$
\mathfrak{Q}_{a}=\sum_{T} \mathfrak{F}_{\mathrm{wt}(T)}
$$

where the sum is over all quasiYamanouchi tableaux $T$ whose support contains the support of $a$, and such that $T \in \mathfrak{A S S T}(b)$ for some weak composition $b$ of length $n$ with $b^{+}=a^{+}$and $b \geq a$.

## Example 4.37

$$
\mathfrak{Q}_{(1,0,3)}=\mathfrak{F}_{(1,0,3)}+\mathfrak{F}_{(2,0,2)}
$$

where the two fundamental slides correspond to the 3rd and 5th composition tableaux in Fig. 3.

Although we claimed that the quasikey polynomials were to be a lift from $\mathrm{QSym}_{n}$ to $\mathrm{ASym}_{n}$ of the quasiSchur polynomials, we have not yet explained this fact. The sense of this lift is given in the following proposition.

Proposition 4.38 ([3, Theorem 4.16]) We have $\mathfrak{Q}_{a} \in \mathrm{QSym}_{n}$ if and only if $a$ is of the form $0^{k} \alpha$, where $\alpha$ is a strong composition of length $n-k$. Moreover, we have

$$
\mathfrak{Q}_{0^{k} \alpha}=S_{\alpha}\left(x_{1}, \ldots, x_{n}\right) .
$$

Thus, the quasikey polynomial basis of $\mathrm{ASym}_{n}$ is a lift of the quasiSchur polynomial basis of $\mathrm{QSym}_{n}$.

Moreover, the quasikey polynomials stabilize to the quasiSchur functions:
Theorem 4.39 ([3, Theorem 4.17]) Let a be a weak composition. Then

$$
\lim _{m \rightarrow \infty} \mathfrak{Q}_{0^{m} a}=S_{a^{+}} .
$$

As Schur polynomials expand positively in the quasiSchur polynomial basis, one might hope for the same positivity to hold for their respective lifts, the key polynomials (or Schubert polynomials) and quasikey polynomials. Indeed, these expansions are positive, with positive combinatorial formulas mirroring the expansion of Theorem 3.12. To provide a formula for this expansion, we need the concept of a left swap on weak compositions. A left swap on a weak composition $a$ exchanges two entries $a_{i}$ and $a_{j}$ such that $a_{i}<a_{j}$ and $i<j$. In essence, left swaps move larger entries leftwards. Given a weak composition $a$ of length $n$, define the set lswap $(a)$ to be all weak compositions $b$ of length $n$ that can be obtained from $a$ by a sequence of left swaps.

Example 4.40 We have

$$
\operatorname{lswap}(1,0,3)=\{(1,0,3),(1,3,0),(3,0,1),(3,1,0)\} .
$$

In fact, the elements of $\operatorname{lswap}(a)$ are exactly those weak compositions $b$ such that $\overleftarrow{b}=\overleftarrow{a}$ and $w(b) \leq w(a)$ in Bruhat order. Hence, in this new language, the formula in Theorem 4.30 for expanding key polynomials in Demazure atoms may be re-expressed as follows.

Proposition 4.41 ([77, Lemma 3.1])

$$
\mathfrak{D}_{a}=\sum_{b \in \operatorname{lswap}(a)} \mathfrak{A}_{b} .
$$

Define $\operatorname{Qlswap}(a)$ to be those $b \in \operatorname{lswap}(a)$ such that for all $c \in \operatorname{lswap}(a)$ with $c^{+}=b^{+}$, one has $c \geq b$ in dominance order.

Example 4.42 The set $\operatorname{lswap}(1,0,3)$ is computed in Example 4.40. Partitioning $\operatorname{lswap}(1,0,3)$ into equivalence classes under the relation $c \sim b$ when $c^{+}=b^{+}$yields classes $\{(1,0,3),(1,3,0)\}$ and $\{(3,0,1),(3,1,0)\}$. Taking the weak composition smallest in dominance order from each class then yields

$$
\operatorname{Qlswap}(1,0,3)=\{(1,0,3),(3,0,1)\}
$$

Theorem 4.43 ([3, Theorem 3.7]) Let a be a weak composition of length $n$. Then

$$
\mathfrak{D}_{a}=\sum_{b \in \operatorname{Qlswap}(a)} \mathfrak{Q}_{b} .
$$

A key polynomial $\mathfrak{D}_{a}$ is a Schur polynomial if and only if the entries of $a$ are weakly increasing. In this case, by Proposition 4.38 and the definition of Qlswap, the formula of Theorem 4.43 reduces to the formula of Theorem 3.12 for the quasiSchur polynomial expansion of a Schur polynomial.

Quasikey polynomials do not have positive structure constants. This is immediate from the fact that the quasikey polynomial basis contains all the quasiSchur polynomials, which themselves do not have positive structure constants. However, [77] gives a positive combinatorial formula for the quasikey expansion of the product of a quasikey polynomial and a Schur polynomial, extending the analogous formula of [31] for quasiSchur polynomials.

The final basis of $\mathrm{ASym}_{n}$ that we consider here is the basis of fundamental particles introduced in [77]. Note that the formula for the Demazure atom expansion of a quasikey polynomial given in Theorem 4.34 is identical to the formula for the monomial expansion of a monomial slide polynomial given in Definition 4.14. The main motivating property of fundamental particles is that this same formula will give the fundamental particle expansion of a fundamental slide polynomial.

Given a weak composition $a$, let $\mathfrak{P S S T}(a)$ denote the set of those semistandard composition tableaux of shape $a$ satisfying the property that whenever $i<j$, every label in row $i$ is smaller than every label in row $j$. For example, $\mathfrak{P S S T}(1,0,3)$ consists of the 3rd, 4th and 6th tableaux in Fig. 3.

Let $a$ be a weak composition of length $n$. The fundamental particle associated to $a$ is given by

$$
\mathfrak{P}_{a}:=\sum_{T \in \mathfrak{P S S T}(a)} \mathbf{x}^{\mathrm{wt}(T)}
$$

Example 4.44 Let $a=(1,0,3)$. Then

$$
\mathfrak{P}_{(1,0,3)}=\mathbf{x}^{103}+\mathbf{x}^{112}+\mathbf{x}^{121}
$$

The following is straightforward from the definitions.

Theorem 4.45 ([77]) The fundamental slide polynomials expand positively in the fundamental particles: For a weak composition a of length n, we have

$$
\mathfrak{F}_{a}=\sum_{\substack{b^{+}=a^{+} \\ b \geq a}} \mathfrak{P}_{b},
$$

where the sum is over all weak compositions $b$ of length $n$ satisfying $b \geq a$ in dominance order and whose positive part is $a^{+}$.

The fundamental particles were constructed to be a refinement of the fundamental slide polynomials (with a particular positive expansion formula), as in Theorem 4.45; remarkably, the fundamental particles are also a refinement of the Demazure atoms. Given a weak composition $a$, define the set $\operatorname{HSST}(a)$ of particle-highest semistandard composition tableaux of shape $a$ to be the set of those $T \in \mathfrak{A S S T}(a)$ such that for each integer $i$ appearing in $T$, either

- an $i$ appears in the first column, or
- there is an $i^{+}$weakly right of an $i$, where $i^{+}$is the smallest integer greater than $i$ appearing in $T$.

Notice the particle-highest condition is a weakening of the quasiYamanouchi condition: every quasiYamanouchi tableau is necessarily particle-highest.

Establishing the last of the arrows shown in Fig. 4, we have the following additional positivity.

Theorem 4.46 ([77]) The Demazure atoms expand positively in the fundamental particles:

$$
\mathfrak{A}_{a}=\sum_{T \in \operatorname{HSST}(a)} \mathfrak{P}_{\mathrm{wt}(T)} .
$$

Fundamental particles do not have positive structure constants, however, in analogy with quasikey polynomials and Demazure atoms, there is a positive combinatorial formula for the fundamental particle expansion of the product of a fundamental particle and a Schur polynomial given in [77].

In this section, we have given three formulas for Schubert polynomials. Each formula lead us to consider certain other families of polynomials described by similar combinatorics. There is additionally a fourth, fundamentally different, formula for Schubert polynomials given by A. Lascoux [45] in terms of the square-ice/6-vertex model of statistical physics. This formula was recently rediscovered in [49] in the guise of bumpless pipe dreams. (The connection between these two combinatorial descriptions is detailed in [84].) We won't describe this rule here, as its context is currently unclear, although we think it is likely to be important. Connections to Gröbner geometry appear in [32].

## 5 The Mirror Worlds: $\boldsymbol{K}$-Theoretic Polynomials

A trend in modern Schubert calculus is to look at Flags $_{n}$, Grassmannians, and other generalized flag varieties, not through the lens of ordinary cohomology as in Sects. 2-4, but through the sharper yet more mysterious lenses of other complex oriented cohomology theories [16, 26, 55, 56].

Particularly well studied over the past 20 years are combinatorial aspects of the $K$-theory rings of these spaces. Early work here includes [20, 22, 52, 53]; however, the area only became very active after the influential work of [14, 46]. Following [8, 20, 33], it turns out that one can slightly generalize this setting to connective $K$-theory with almost no extra combinatorial complexity (indeed, in some ways the more general combinatorics seems easier). In general, each complex oriented cohomology theory is determined by its formal group law, which describes how to write the Chern class of a tensor product of two line bundles in terms of the two original Chern classes. In the case of connective $K$-theory, the formal group law is

$$
\begin{equation*}
c_{1}(L \otimes M)=c_{1}(L)+c_{1}(M)+\beta c_{1}(L) c_{1}(M), \tag{5.1}
\end{equation*}
$$

where $\beta$ is a formal parameter and $L, M$ are complex line bundles on the space in question. In this notation, the ordinary cohomology ring is recovered by setting $\beta=0$ and the ordinary $K$-theory ring is recovered (up to convention choices) by setting $\beta=-1$. For more background on connective $K$-theory, see the appendix to [1].

Just as the Schubert classes in the ordinary cohomology of Flags ${ }_{n}$ are represented by the Schubert polynomials (as described in Sect.4), we would like to have such polynomial representatives for the corresponding connective $K$-theory classes. These are provided by the $\beta$-Grothendieck polynomials $\left\{\overline{\mathfrak{S}}_{a}\right\}$ of S. Fomin and A. Kirillov [20], as identified in [33]. These polynomials form a basis of $\operatorname{ASym}_{n}[\beta]$, where $\beta$ is the formal parameter from Eq. (5.1). This basis is homogeneous if the parameter $\beta$ is understood to live in degree -1 . Specializing at $\beta=0$, one recovers the Schubert basis $\left\{\mathfrak{S}_{a}\right\}$ of $\operatorname{ASym}_{n}$. The usual Grothendieck polynomials of A. Lascoux and M.-P. Schützenberger [52] are realized at $\beta=-1$. (To help the reader track relations among bases, we deviate from established practice by denoting connective $K$-analogues by applying an 'overbar' to their cohomological specializations.) Like the Schubert basis of $\mathrm{ASym}_{n}$, the $\beta$-Grothendieck polynomial basis has positive structure coefficients; this is, of course, currently only known by geometric arguments [11], as there is no combinatorial proof of this fact even in the case $\beta=0$.

Intersecting $\left\{\bar{S}_{a}\right\}$ with $\operatorname{Sym}_{n}[\beta]$ yields the basis $\left\{\bar{s}_{\lambda}\right\}$ of symmetric Grothendieck polynomials. These represent connective $K$-theory Schubert classes on Grassmannians. In this setting, like the Schur polynomial setting, a number of LittlewoodRichardson rules for $\left\{\bar{s}_{\lambda}\right\}$ are now known (e.g., [73, 80, 81]), following the first found by A. Buch [14].

There is a beautiful combinatorial formula for $\bar{s}_{\lambda}$ given by Buch [14], as a direct extension of the Littlewood formula of Theorem 2.1 for Schur polynomials. A set-
valued tableau $T$ of shape $\lambda$ is a filling of each cell of the Young diagram $\lambda$ by a nonempty set of positive integers. Deleting all but one number from each set resolves $T$ to an ordinary tableau, and $T$ is called semistandard if all all such resolutions are semistandard in the sense of Sect. 2. In other words, $T$ is semistandard if the greatest entry in each box is not larger than the least entry in the box to its right and is strictly smaller than the least entry in the box directly below it. Let $\operatorname{SV}(\lambda)$ denote the set of semistandard set-valued tableaux of shape $\lambda$. The weight of a set-valued tableau $T$ is the weak composition $\operatorname{wt}(T):=\left(a_{1}, a_{2}, \ldots\right)$, where $a_{i}$ records the number of instances of the number $i$ among all boxes of $T$.

Theorem 5.1 ([14, Theorem 3.1]) For any partition $\lambda$, we have

$$
\bar{s}_{\lambda}=\sum_{T \in \operatorname{SV}(\lambda)} \mathbf{x}^{\mathrm{wt}(T)}
$$

The remaining families of polynomials discussed in Sects. 2-4 are not currently understood well in term of cohomology. Remarkably, however, from a combinatorial perspective they all appear to have natural 'connective $K$-analogues'. That is, for each basis, there is a combinatorially-natural $\beta$-deformation that is homogeneous (with the understanding that $\beta$ has degree -1 ), forms a basis of $\operatorname{ASym}_{n}[\beta]$, and (at least conjecturally) shares the positivity properties of the original basis. These deformed bases are presented in Table 1.

It is a mystery as to whether these various apparently $K$-theoretic families of polynomials in fact have a geometric interpretation in terms of $K$-theory. If they did, then presumably the $\beta=0$ specialization considered in the previous sections would similarly have a cohomological interpretation. Such an interpretation would be rather surprising, as currently these specializations are only understood through combinatorics and (in some cases) representation theory. Alternatively, perhaps there is a general representation-theoretic construction that yields all of these various $\beta$ deformations. No such construction is currently known, but for some ideas along these lines see [24, 63, 69].

For details of the definitions of the bases from Table 1 and their relations, see the references given there, especially [64] which contains a partial survey. Here, we only briefly sketch hints of this theory. (However, the theory is in many ways exactly parallel to that given in Sects.2-4, so the astute reader can likely guess approximations to many of the structure theorems.)

The notion of semistandard set-valued skyline fillings was introduced in [62], and employed there to provide an explicit combinatorial definition of the Lascoux polynomials $\overline{\mathfrak{D}}_{a}$ and the Lascoux atoms $\overline{\mathfrak{A}}_{a}$ : K-theoretic analogues of key polynomials and Demazure atoms, respectively. Such $K$-theoretic analogues have also been studied in [13, 38, 44, 63, 64, 70, 75].

The basis $\overline{\mathfrak{F}}_{a}$ of glide polynomials was introduced in [71]. Glide polynomials are simultaneously a $K$-theoretic analogue of the fundamental slide basis and a polynomial lift of the multi-fundamental quasisymmetric basis [51] of quasisymmetric polynomials. Remarkably, the glide basis also has positive structure constants, which

Table 1 Bases from Sects. 2-4, together with their corresponding connective $K$-analogues. For each $K$-theoretic family of polynomials, we have a given a few major references; however, these references are generally not exhaustive

| Sym $_{n}$ | Schur polynomial $s_{\lambda}$ | Symmetric Grothendieck polynomial $\bar{s}_{\lambda}$ <br> $[14,63]$ |
| :--- | :--- | :--- |
| QSym $_{n}$ | Monomial quasisymmetric polynomial <br> $M_{\alpha}$ | Multimonomial polynomial $\bar{M}_{\alpha}[51]$ |
|  | Fundamental quasisymmetric polynomial |  |
| $F_{\alpha}$ | QuasiSchur polynomial $S_{\alpha}$ | Multifundamental polynomial $\bar{F}_{\alpha}[51$, <br> $68,71]$ <br> QuasiGrothendieck polynomial $\bar{S}_{\alpha}[62$, <br> $64]$ |
| $\operatorname{ASym}_{n}$ | Schubert polynomial $\mathfrak{S}_{a}$ <br> Demazure character/key polynomial $\mathfrak{D}_{a}$ <br> Quasikey polynomial $\mathfrak{Q}_{a}$ | Grothendieck polynomial $\overline{\mathfrak{S}}_{a}[20,39,52]$ <br> Lascoux polynomial $\overline{\mathfrak{D}}_{a}[38,62,64,75]$ <br> QuasiLascoux polynomial $\overline{\mathfrak{Q}}_{a}[64]$ <br> Demazure atom/standard basis $\mathfrak{A}_{a}$ <br> Fundamental particle/pion $\mathfrak{P}_{a}$ <br> Fundamental slide polynomial $\mathfrak{F}_{a}$ |
| Lascoux atom $\overline{\mathfrak{A}}_{a}[13,62,64,70]$ <br> kaon $\overline{\mathfrak{P}}_{a}[64]$ <br> glide polynomial $\overline{\mathfrak{F}}_{a}[64,71]$ |  |  |

can be described in terms the glide product [71] of weak compositions. The glide product is a simultaneous generalization of the slide product of [2] and the multishuffle product of [51] on strong compositions.

The quasiLascoux basis $\overline{\mathfrak{Q}}_{a}$ and kaon basis $\overline{\mathfrak{P}}_{a}$ were introduced in [64]. These are $K$-analogues of the quasi-key and fundamental particle bases. Remarkably, the positivity relations between the bases in Fig. 4 have been proven to hold for their $K$-analogues mentioned above [64, 71], with the exception of the expansion of Grothendieck polynomials in Lascoux polynomials, whose positivity remains conjectural.

As outlined in the previous section, the product of an element of any basis in Fig. 4 with a Schur polynomial expands positively in that basis. It would be interesting to know if the analogous result is true in the $K$-theory world: that product of an element of a $K$-theoretic analogue and a symmetric Grothendieck polynomial expands positively in that $K$-theoretic basis. This is obviously true for the Grothendieck basis by geometry. It is also true for the glide basis, since the glide basis has positive structure constants and refines Grothendieck polynomials. We believe this question remains, however, open for the Lascoux, quasiLascoux, Lascoux atom and kaon bases.

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# Minuscule Schubert Calculus and the Geometric Satake Correspondence 

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#### Abstract

We describe a relationship between work of Gatto, Laksov, and their collaborators on realizations of (generalized) Schubert calculus of Grassmannians, and the geometric Satake correspondence of Lusztig, Ginzburg, and Mirković and Vilonen. Along the way we obtain new proofs of equivariant Giambelli formulas for the ordinary and orthogonal Grassmannians, as well as a simple derivation of the "rim-hook" rule for computing in the equivariant quantum cohomology of the Grassmannian.


Keywords Geometric Satake correspondence - Schubert calculus • Affine
Grassmannian • Pfaffian • Equivariant quantum cohomology

## 1 Introduction

The goal of this article is to illustrate connections between several circles of ideas in Schubert calculus, representation theory, and symmetric functions. The driving force behind the connections we explore is the geometric Satake correspondence-which, in the special cases we examine, matches Schubert classes in the Grassmannian (and related spaces) with weight vectors in exterior products (and related representations). We will focus especially on an equivariant version of this correspondence, where additional bases appear on both sides. An analysis of the transition matrices for these bases leads us directly to well-known symmetric functions: the Schur $S$ - and $P$-functions, along with their factorial generalizations.

[^10]The approach we take is primarily expository, at least in the sense that all the main theorems have appeared previously. However, in describing a relationship among notions which do not often appear side by side, we obtain some novel consequences: new and simple proofs of the equivariant Giambelli formulas for ordinary and orthogonal Grassmannians, as well as of a rule for computing in the equivariant quantum cohomology of the Grassmannian. This perspective should find further applications within Schubert calculus. One theme we wish to emphasize is this: when some aspect of $H_{T}^{*} G r(k, n)$ appears related to exterior algebra, one can expect generalizations via the Satake correspondence, either to other minuscule spaces, or to other subvarieties of the affine Grassmannian.

Our starting point is a very simple observation, which must be quite old. The Grassmannian $\operatorname{Gr}(k, n)$ has a decomposition into Schubert cells, indexed by $k$-element subsets of $[n]:=\{1, \ldots, n\}$. As $\mathbb{C}$-vector spaces, therefore, one has

$$
\begin{equation*}
H_{*}(G r(k, n), \mathbb{C})=H^{*}(G r(k, n), \mathbb{C})=\bigwedge^{k} \mathbb{C}^{n} \tag{1}
\end{equation*}
$$

and the cohomology ring $H^{*}(\operatorname{Gr}(k, n), \mathbb{C})$ acts as a certain ring of operators on the exterior power. The rest of this story is an attempt to add as much structure as possible to this identification.

As a first step, for $k=1$, let us identify the basis of linear subspaces of $\mathbb{P}^{n-1}$ with the standard basis of $\mathbb{C}^{n}$ :

$$
\left[\mathbb{P}^{n-i}\right]=\varepsilon_{i} \quad \text { in } \quad H_{n-i} \mathbb{P}^{n-1}=H^{i-1} \mathbb{P}^{n-1}
$$

(Since the spaces we consider have no odd-degree singular (co)homology, we will economize notation by always writing $H_{i}$ and $H^{i}$ for singular homology and cohomology in degree $2 i$.) This induces a grading on $\mathbb{C}^{n}$, and makes (1) an isomorphism of graded vector spaces. More generally, writing $\Omega_{I}$ for the Schubert variety corresponding to $I=\left\{i_{1}<\cdots<i_{k}\right\} \subseteq[n]$, we identify

$$
\left[\Omega_{I}\right]=\varepsilon_{I}:=\varepsilon_{i_{1}} \wedge \cdots \wedge \varepsilon_{i_{k}}
$$

in $H_{*} \operatorname{Gr}(k, n)=H^{*} G r(k, n)$.
To be a little more specific, since it will matter later, these are the opposite Schubert varieties: $\Omega_{I}$ is the closure of the cell $\Omega_{I}^{\circ}$ whose representing $k \times n$ matrices have pivots in columns $I$, with zeroes to the left of the pivots. For example, in $\operatorname{Gr}(3,8)$ we have

$$
\Omega_{\{2,4,7\}}^{\circ}=\left[\begin{array}{ccccccc}
0 & 1 & * & 0 & * & * & 0
\end{array}\right]+\left[\begin{array}{cccccc}
0 & 0 & 0 & 1 & * & * \\
0 & 0 & 0 & 0 & 0 & 0
\end{array}\right]
$$

There is a standard bijection between $I \subseteq[n]$ and partitions $\lambda$ fitting inside the $k \times(n-k)$ rectangle; one puts

$$
\lambda_{k+1-j}=i_{j}-j
$$

Then, writing $\Omega_{\lambda}=\Omega_{I}$, the grading is realized by $\operatorname{codim} \Omega_{\lambda}=|\lambda|=\lambda_{1}+\cdots+\lambda_{k}$. We will write $\sigma_{\lambda}$ for the (co)homology class $\left[\Omega_{\lambda}\right]$.

The essential idea is to exploit the isomorphism (1) and use linear algebra to develop the basic ingredients of Schubert calculus-including Giambelli and Pieri formulas. This can be done using elementary (though still nontrivial) methods.

Again, the first step is to examine the easiest case, where $k=1$. The action of the divisor class $\sigma_{\square}$ translates into an operator $\xi$ on $\mathbb{C}^{n}$, given by

$$
\xi \cdot \varepsilon_{i}= \begin{cases}\varepsilon_{i+1} & \text { if } i<n \\ 0 & \text { if } i=n\end{cases}
$$

So this has matrix

$$
\xi=\left[\begin{array}{lllll}
0 & & & &  \tag{2}\\
1 & \ddots & & \\
0 & \ddots & \ddots & \\
0 & 0 & 1 & 0
\end{array}\right]
$$

and one can view it as an element of the Lie algebra $\mathfrak{g l}_{n}\left(\right.$ or $\left.\mathfrak{s l}_{n}\right)$. This Lie algebra acts naturally on exterior powers $\bigwedge^{k} \mathbb{C}^{n}$, and the basic observation is that for any $k$, $\sigma_{\square}$ acts on $H^{*} G r(k, n)$ (via cup product) just as the matrix $\xi$ acts on $\bigwedge^{k} \mathbb{C}^{n}$ (via the Lie algebra action).

For example, consider $\sigma_{1}$ acting on $H^{*} G r(2, n)$. We have

$$
\begin{aligned}
\xi \cdot \varepsilon_{1} \wedge \varepsilon_{2} & =\varepsilon_{2} \wedge \varepsilon_{2}+\varepsilon_{1} \wedge \varepsilon_{3} \\
& =\varepsilon_{1} \wedge \varepsilon_{3}
\end{aligned}
$$

corresponding to $\sigma_{\square} \cdot \sigma_{\emptyset}=\sigma_{(1)}$. Similarly,

$$
\xi \cdot \varepsilon_{1} \wedge \varepsilon_{3}=\varepsilon_{2} \wedge \varepsilon_{3}+\varepsilon_{1} \wedge \varepsilon_{4}
$$

corresponding to $\sigma_{\square} \cdot \sigma_{(1)}=\sigma_{(1,1)}+\sigma_{(2)}$. Note, however, that the Lie algebra action of the matrix $\xi^{2}$ does not correspond to multiplication by $\sigma_{\square}^{2}$. For instance, the Lie algebra action gives

$$
\begin{aligned}
\xi^{2} \cdot \varepsilon_{2} \wedge \varepsilon_{3} & =\varepsilon_{4} \wedge \varepsilon_{3}+\varepsilon_{2} \wedge \varepsilon_{5} \\
& =\varepsilon_{2} \wedge \varepsilon_{5}-\varepsilon_{3} \wedge \varepsilon_{4}
\end{aligned}
$$

which corresponds to the computation $p_{2} \cdot \sigma_{(1,1)}=\sigma_{(3,1)}-\sigma_{(2,2)}$, where $p_{2}$ is the power sum symmetric polynomial.

The question thus arises: are the above calculations merely coincidental? In the next two sections, we will see that they are not, by viewing them as a shadow of the geometric Satake correspondence, a major construction in modern representation theory.

In Sect.4, we show how the equivariant Giambelli formula-which computes a Schubert class in $H_{T}^{*} G r(k, n)$ as a factorial Schur polynomial $s_{\lambda}(x \mid t)$-follows directly from the defining properties of the exterior power, putting computations of Gatto, Laksov, Thorup, and others into the general context of the Satake correspondence. Having done this, it is natural to proceed to the minuscule spaces of type D: in Sect. 6, we apply similar methods to compute equivariant Schubert classes via the factorial Schur $P$-functions $P_{\lambda}(x \mid t)$, using computations on even-dimensional quadrics (carried out in Sect. 5) in place of projective spaces. The functions $s_{\lambda}(x \mid t)$ are frequently defined as a certain "Jacobi-Trudi" determinant, but they may also be written as a ratio of two determinants (as was originally done by Cauchy); similarly, $P_{\lambda}(x \mid t)$ may be written either as a Pfaffian or as a ratio of two Pfaffians. A curious aspect of our arguments is that the ratio description of these functions appears naturally, in contrast to most geometric arguments (dating to Giambelli), where the Jacobi-Trudi formulation is used.

We turn to quantum cohomology in Sect. 7, where we give a short proof of the equivariant rim-hook rule for computing in $Q H_{T}^{*} G r(k, n)$. Here the Satake isomorphism serves only a psychological function, and is not logically necessary: the main point is that the combinatorial operation of removing a rim hook from a partition (and picking up a corresponding sign) is precisely that of reducing the indices of a pure wedge modulo $n$.

Most of these ideas have appeared in the work of other authors, at least in some form; as mentioned above, our primary aim is to indicate connections and extract a few new consequences. We first learned of the possibility of a "formal" Schubert calculus on exterior powers from a series of papers by Gatto and Laksov-Thorup in the 2000s, and this point of view has been developed further by these authors and their collaborators [14-17, 30-33]. More detailed references are given throughout the article, and we point to further connections in a closing section (Sect. 8).

## 2 The Geometric Satake Correspondence

A second simple observation is the following: On one hand, the vector space $\bigwedge^{k} \mathbb{C}^{n}$ is a fundamental (and in fact, minuscule) representation $V_{\varpi_{k}}$ of $\mathfrak{s l}_{n}$. On the other hand, $\operatorname{Gr}(k, n)=P G L_{n} / P_{\varpi_{k}}$, where $P_{\varpi_{k}}$ is the parabolic corresponding ${ }^{1}$ to the cocharacter $\varpi_{k}$ of $P G L_{n}$. In fact, if $V=\mathbb{C}^{n}$ is the standard representation of $\mathfrak{s l}_{n}$, it is most natural to regard $\operatorname{Gr}(k, n)=G r\left(k, V^{*}\right)$ as parametrizing $k$-planes in the dual vector space.

[^11]Work from the 1990's by Ginzburg [18] and Mirković and Vilonen [41]—which in turn builds on work of Lusztig [35] from the early 1980's—puts this into a more general context. To describe it we need some terminology and basic facts.

Any reductive group $G$ with maximal torus $T$ comes with a root datum $R$. Root data have a built-in duality, and exchanging $R$ with $R^{\vee}$ yields a Langlands dual group $G^{\vee}$. The details will not be too important for now, beyond this: the roots of $G^{\vee}$ are the coroots of $G$, and the characters of the maximal torus $T^{\vee} \subseteq G^{\vee}$ are the co-characters of $T \subseteq G$. For example, $\left(P G L_{n}\right)^{\vee} \cong S L_{n},\left(S p_{2 n}\right)^{\vee} \cong S O_{2 n+1}$, $\left(P S O_{2 n}\right)^{\vee}=\operatorname{Spin}_{2 n}$, and $\left(G L_{n}\right)^{\vee} \cong G L_{n}$. (A significant part of Ginzburg and Mirković-Vilonen's program was to give a more intrinsic construction of $G^{\vee}$.)

Let $\mathcal{K}=\mathbb{C}((z))$ and $\mathcal{O}=\mathbb{C} \llbracket z \rrbracket$. The affine Grassmannian of a complex reductive group $G$ is the infinite-dimensional orbit space

$$
\mathcal{G} r_{G}=G(\mathcal{K}) / G(\mathcal{O}),
$$

topologized as an ind-variety. The essence of geometric Satake is to relate the geometry of $\mathcal{G} r_{G}$ with the representation theory of $G^{\vee}$. We will describe a very small part of this correspondence which suffices for our purposes.

The group $G(\mathcal{O})$ acts on $\mathcal{G} r_{G}$ via left multiplication, and its orbits are naturally parametrized by dominant co-characters $\varpi: \mathbb{C}^{*} \rightarrow T$. These, by duality, are the same as dominant characters of $T^{\vee}$. Since there is a well-known indexing of irreducible representations of $G^{\vee}$ by dominant characters, we have a bijection of sets

$$
\begin{aligned}
\left\{G(\mathcal{O}) \text {-orbit closures in } \mathcal{G} r_{G}\right\} & \leftrightarrow \\
\overline{\mathcal{G} r^{\varpi}} & \leftrightarrow
\end{aligned}
$$

and as before the goal is to endow this with more structure.
Let $\mathfrak{g}^{\vee}=\operatorname{Lie}\left(G^{\vee}\right)$, and take a regular nilpotent element $\xi \in \mathfrak{g}^{\vee}$. (Up to conjugation, in $\mathfrak{s l}_{n}$ such an element is the matrix from (2). More generally, one can write $\xi=\sum a_{i} E_{\alpha_{i}}$ as a sum of simple root vectors.) Let $\mathfrak{a} \subseteq \mathfrak{g}^{\vee}$ be the centralizer of $\xi$, an abelian Lie subalgebra of dimension equal to the rank of $\mathfrak{g}$. (In the case of $\xi \in \mathfrak{s l}_{n}$, this subalgebra is spanned by the matrix powers $\xi, \xi^{2}, \ldots, \xi^{n-1}$.) Its universal enveloping algebra, denoted $\mathfrak{U}(\mathfrak{a})$, acts naturally on any representation of $\mathfrak{a}$. Since $\mathfrak{a}$ is abelian, $\mathfrak{U}(\mathfrak{a})$ is just the polynomial algebra $\operatorname{Sym}_{\mathbb{C}}^{*} \mathfrak{a}$.

Finally, $I H_{*} X$ denotes the (middle-perversity) intersection homology of a space $X$, with coefficients in $\mathbb{C}$. This is a graded vector space which exhibits Poincaré duality, and which comes with an action of $H^{*} Y$, for any $X \rightarrow Y$, via cap product.

Theorem 1 (Geometric Satake [18, 41]) There are graded isomorphisms of algebras

$$
H^{*}\left(\mathcal{G} r_{G}^{\circ}\right) \cong \mathfrak{U}(\mathfrak{a})=\operatorname{Sym}_{\mathbb{C}}^{*} \mathfrak{a}
$$

and vector spaces

$$
I H_{*}\left(\overline{\mathcal{G} r^{\varpi}}\right) \cong V_{\varpi}
$$

for all dominant $\varpi$, and these isomorphisms are compatible with the natural actions of $H^{*}\left(\mathcal{G} r_{G}\right)$ on $I H_{*}\left(\overline{\mathcal{G} r^{\varpi}}\right)$ (via cap product) and of $\mathfrak{U}(\mathfrak{a})$ on $V_{\varpi}$ (via the representation of $\left.\mathfrak{g}^{\vee}\right)$. Furthermore, there is a natural basis of $M V$-cycles in $I H_{*}\left(\overline{\mathcal{G r} r^{\varpi}}\right)$ which corresponds to a weight basis of $V_{\varpi}$.

The statements proved by Ginzburg and Mirković-Vilonen are vastly stronger: they establish an equivalence of tensor categories between the category of $G(\mathcal{O})$ equivariant perverse sheaves and the representation category of $G^{\vee}$. We will not need this level of generality, however.

The connected components of $\mathcal{G} r_{G}$ are indexed by elements of $\pi_{1}(G)$, and in fact there is a natural group structure on the set of components. In the statement of the theorem, $\mathcal{G} r_{G}^{\circ}$ means the identity component, although in fact all components are isomorphic-as spaces, but not compatibly with the left $G(\mathcal{O})$-action.

In general, the orbit closure $\overline{\mathcal{G} r^{\varpi}}$ is singular, hence the appearance of intersection homology. However, the minimal orbit in each connected component of $\mathcal{G r} r_{G}$ is closed, so such $\overline{\mathcal{G} r^{\varpi}}=\mathcal{G} r^{\varpi}$ are smooth, and one has $I H_{*}=H_{*}=H^{*}$. When $G$ is adjoint (so $G^{\vee}$ is simply connected), these minimal orbits correspond to the minuscule weights of $G^{\vee}$. For a minuscule weight $\varpi$, one has $\mathcal{G} r^{\varpi}=G / P_{\varpi}$. Furthermore, in this case the MV-cycles are precisely the Schubert varieties in $G / P_{\varpi}$ (as noted in [26, Sect. 1]).

Example 2 The minuscule weights of $\mathfrak{s l}_{n}$ are $0, \varpi_{1}, \ldots, \varpi_{n-1}$, corresponding to the $n$ elements of $\pi_{1}\left(P G L_{n}\right)=\mathbb{Z} / n \mathbb{Z}$. The representations are the exterior powers $\bigwedge^{k} \mathbb{C}^{n}$, and the orbits are the Grassmannians $P G L_{n} / P_{\varpi_{k}}=\operatorname{Gr}(k, n)$, for $0 \leq k \leq$ $n-1$.

The minuscule weights of $\mathfrak{s o}_{2 n}$ are $0, \varpi_{1}, \varpi_{n-1}, \varpi_{n}$, where the nonzero ones are the fundamental weights corresponding to the three end-nodes of the $D_{n}$ Dynkin diagram. The representations are the standard one, $V_{\varpi_{1}}=\mathbb{C}^{2 n}$, and the half-spin representations, $V_{\varpi_{n-1}}=\mathbb{S}_{n}^{+}$and $V_{\varpi_{n}}=\mathbb{S}_{n}^{-}$. The orbits are the quadric $\mathcal{Q}^{2 n-2}=$ $P S O_{2 n} / P_{\varpi_{1}}$, and the two maximal orthogonal Grassmannians $O G^{+}(n, 2 n)=$ $P S O_{2 n} / P_{\varpi_{n-1}}$ and $O G^{-}(n, 2 n)=P S O_{2 n} / P_{\varpi_{n}}$.

In type $C_{n}$, the weight $\varpi_{1}$ corresponding to the standard representation of $\mathfrak{s p}_{2 n}$ is minuscule, and in type $B_{n}$, the weight $\varpi_{n}$ corresponding to the spin representation of $\mathfrak{s o}_{2 n+1}$ is minuscule. The minimal orbits are isomorphic to $\mathbb{P}^{2 n-1}$ and $O G^{+}(n+$ $1,2 n+2$ ), respectively, so they already occur in types $A$ and $D$.

Among simple groups $G$, there are only a few other instances of nonzero minuscule weights. In type $E_{6}$, the weights $\varpi_{1}$ and $\varpi_{6}$ are minuscule, corresponding to the 27dimensional Jordan algebra representation and its dual; the corresponding varieties are the octonionic projective plane and its dual. In type $E_{7}$, there is one nonzero minuscule weight, whose corresponding representation is the 56-dimensional Brown algebra, and whose corresponding homogeneous space is known as the Freudenthal variety.

As a special case of Theorem 1, we have the isomorphism

$$
H_{*} G r(k, n)=\bigwedge^{k} \mathbb{C}^{n}
$$

together with the compatible actions by divisor class and regular nilpotent, described in the introduction-in particular, it is no coincidence that one can do this. We will push this further to obtain a new perspective on Laksov's computation of the equivariant cohomology of $\operatorname{Gr}(k, n)$ in Sect. 4 .

Remark 3 Let $\Lambda_{\mathbb{C}}$ be the ring of symmetric functions with coefficients in $\mathbb{C}$. It can be identified with the infinite polynomial ring $\mathbb{C}\left[p_{1}, p_{2}, \ldots\right]$, where $p_{r}=x_{1}^{r}+x_{2}^{r}+\cdots$ is the power sum symmetric function. For $G=P G L_{n}$, Bott [5] showed that there is a natural map

$$
\Lambda_{\mathbb{C}} \rightarrow H^{*} \mathcal{G} r_{G}^{\circ},
$$

identifying the RHS as $\Lambda_{\mathbb{C}} /\left(p_{n}, p_{n+1}, \ldots\right) \cong \mathbb{C}\left[p_{1}, \ldots, p_{n-1}\right]$. Furthermore, this identifies the subspace $\mathfrak{P} \subseteq H^{*} \mathcal{G} r_{G}^{\circ}$ of primitive classes with the space spanned by $\left\{p_{1}, \ldots, p_{n-1}\right\}$. Ginzburg's proof of the first part of Theorem 1 establishes an isomorphism $\mathfrak{a} \cong \mathfrak{P}$. So the power sum symmetric functions play a central role in this story; we will see an echo of this in the rim-hook rule for quantum cohomology (Sect. 7). (A word of caution: this isomorphism does not hold when one takes cohomology with coefficients in $\mathbb{Z}$. See [5, Proposition 8.1] for a more precise statement.)

## 3 The Equivariant Correspondence

There is an equivariant version of Theorem 1, whose proof is sketched in [18]. We will write $t$ for the generic element of the Cartan subalgebra $\mathfrak{t}^{\vee} \subseteq \mathfrak{g}^{\vee}$ and use the notation $\mathfrak{g}^{\vee}[t]=\mathfrak{g}^{\vee} \otimes \mathbb{C}\left[\mathfrak{t}^{\vee}\right]$ for the Lie algebra over the polynomial ring. Given a $\mathfrak{g}^{\vee}$-module $V$, there is an induced $\mathfrak{g}^{\vee}[t]$-module $V[t]:=V \otimes \mathbb{C}\left[\mathfrak{t}^{\vee}\right]$, where the action is given by

$$
(x \otimes f) \cdot(v \otimes g)=(x \cdot v) \otimes(f g),
$$

for $x \in \mathfrak{g}^{\vee}, v \in V$, and $f, g \in \mathbb{C}\left[\mathfrak{t}^{\vee}\right]$.
Next suppose $\mathfrak{b}^{\vee} \subseteq \mathfrak{g}^{\vee}$ is a Borel subalgebra containing $\mathfrak{t}^{\vee}$. Any character $\chi$ of $\mathfrak{t}^{\vee}$ extends to one of $\mathfrak{b}^{\vee}$, and also to $\mathfrak{b}^{\vee}[t]$. If $V$ is a $\mathfrak{b}^{\vee}$-module, we can twist it by the character $\chi$ to obtain modules $V(\chi)$ and $V(\chi)[t]$ for $\mathfrak{b}^{\vee}$ and $\mathfrak{b}^{\vee}[t]$, respectively. Concretely, if one writes an element of $\mathfrak{b}^{\vee}=\mathfrak{n}^{\vee} \oplus \mathfrak{t}^{\vee}$ as $x=n+t$, then for $f, g \in$ $\mathbb{C}\left[\mathfrak{t}^{\vee}\right]$ and $v \in V(\chi)$ a weight vector for $\mathfrak{t}^{\vee}$, we have

$$
\begin{aligned}
(x \otimes f) \cdot(v \otimes g) & =(n \otimes f+t \otimes f) \cdot(v \otimes g) \\
& =(n \cdot v) \otimes(f g)+(t \cdot v) \otimes(f g)+\chi(t \otimes f)(v \otimes g) .
\end{aligned}
$$

Now let $\xi_{t}=\xi-t$ in $\mathfrak{g}^{\vee}[t]=\mathfrak{g}^{\vee} \otimes \mathbb{C}\left[\mathfrak{t}^{\vee}\right]$, where $\xi$ is a principal nilpotent as before, and $t$ is the generic element of the Cartan subalgebra $\mathfrak{t}^{\vee} \subseteq \mathfrak{g}^{\vee}$. Concretely, for $\mathfrak{g l}_{n}$ this is

$$
\xi_{t}=\left[\begin{array}{cccc}
-t_{1} & & &  \tag{3}\\
1 & \ddots & & \\
0 & \ddots & \ddots & \\
0 & 0 & 1 & -t_{n}
\end{array}\right]
$$

Let $\mathfrak{a}_{t} \subseteq \mathfrak{g}^{\vee}[t]$ be the centralizer of $\xi_{t}$. For $\mathfrak{g l}_{n}$, this subalgebra is spanned over $\mathbb{C}\left[\mathfrak{t}^{\vee}\right]=$ Sym $^{*} \mathfrak{t}^{\vee}$ by the matrix powers $1, \xi_{t}, \xi_{t}^{2}, \ldots, \xi_{t}^{n-1}$.

Theorem 4 (Equivariant Geometric Satake) There are isomorphisms

$$
H_{T}^{*} \mathcal{G} r_{G}^{\circ} \cong \mathfrak{U}_{\mathbb{C}\left[t^{\vee}\right]}\left(\mathfrak{a}_{t}\right)=\operatorname{Sym}_{\mathbb{C}\left[\mathrm{t}^{\vee}\right]}^{*} \mathfrak{a}_{t}
$$

inducing compatible actions on

$$
I H_{*}^{T}\left(\overline{\mathcal{G} r^{\varpi}}\right) \cong V_{\varpi}(-\varpi)[t] .
$$

The effect of twisting by the character $-\varpi$ is to move the highest weight vector of $V_{\varpi}$ to weight zero. This corresponds to endowing $I H_{*}^{T}\left(\mathcal{G} r^{\varpi}\right)$ with a $\mathfrak{t}^{\vee}$-module structure so that the fundamental class $\left[\overline{\mathcal{G} r^{\varpi}}\right]$ has weight 0 . This choice has the advantage of identifying the action of the element $\xi_{t}$ with equivariant multiplication by the divisor class $\sigma_{\square}$. (Of course, a similar isomorphism holds without the twist; noting that $\sigma_{\square}=c_{1}^{T}(\mathcal{O}(1) \otimes(-\varpi))$, where $\mathcal{O}(1)$ is the ample line bundle corresponding to the weight $\varpi$, the untwisted version identifies the action of $\xi_{t}$ with multiplication by $c_{1}^{T}(\mathcal{O}(1))$.)

Example 5 For $\mathfrak{g}=\mathfrak{g l}_{n}$, the action of $\xi_{t}$ on $V_{\varpi_{2}}\left(-\varpi_{2}\right)[t]=\left(\bigwedge^{2} \mathbb{C}^{n} \otimes\left(-t_{1}-\right.\right.$ $\left.\left.t_{2}\right)\right) \otimes \mathbb{C}\left[\mathfrak{t}^{\vee}\right]$ is as follows.

$$
\begin{aligned}
\xi_{t} \cdot\left(\varepsilon_{1} \wedge \varepsilon_{2}\right) & =\varepsilon_{2} \wedge \varepsilon_{2}+\varepsilon_{1} \wedge \varepsilon_{3}-\left(t_{1}+t_{2}-t_{1}-t_{2}\right) \varepsilon_{1} \wedge \varepsilon_{2} \\
& =\varepsilon_{1} \wedge \varepsilon_{3}
\end{aligned}
$$

corresponding to $\sigma_{\square} \cdot \sigma_{\emptyset}=\sigma_{(1)}$ in $H_{T}^{*} G r(2, n)$. (The cancellation of the last term shows why the twist by $-\varpi$ is necessary.) Similarly,

$$
\begin{aligned}
\xi_{t} \cdot\left(\varepsilon_{1} \wedge \varepsilon_{3}\right) & =\varepsilon_{2} \wedge \varepsilon_{3}+\varepsilon_{1} \wedge \varepsilon_{4}-\left(t_{1}+t_{3}-t_{1}-t_{2}\right) \varepsilon_{1} \wedge \varepsilon_{3} \\
& =\varepsilon_{2} \wedge \varepsilon_{3}+\varepsilon_{1} \wedge \varepsilon_{4}+\left(t_{2}-t_{3}\right) \varepsilon_{1} \wedge \varepsilon_{3},
\end{aligned}
$$

corresponding to $\sigma_{\square} \cdot \sigma_{(1)}=\sigma_{(1,1)}+\sigma_{(2)}+\left(t_{2}-t_{3}\right) \sigma_{(1)}$.
As in the non-equivariant case, one needs to beware of the notation: matrix powers $\xi_{t}^{j}$ do not correspond to iterates of the Lie algebra action, e.g., $\xi_{t}^{2} \cdot \varepsilon_{I}$ is generally not equal to $\xi_{t} \cdot\left(\xi_{t} \cdot \varepsilon_{I}\right)$.

Example 6 Still in the case $\mathfrak{g}=\mathfrak{g l}_{n}$, let us consider higher powers of $\xi_{t}$. One computes the entries of the matrix powers as

$$
\begin{aligned}
\xi_{t}^{j} \varepsilon_{i} & =\varepsilon_{i+j}-h_{1}\left(t_{i}, \ldots, t_{i+j-1}\right) \varepsilon_{i+j-1}+\cdots+(-1)^{j} h_{j}\left(t_{i}\right) \varepsilon_{i} \\
& =\sum_{a=0}^{j}(-1)^{a} h_{a}\left(t_{i}, \ldots, t_{i+j-a}\right) \varepsilon_{i+j-a},
\end{aligned}
$$

where the $h_{a}$ are complete homogeneous symmetric polynomials in the indicated variables. (That is, the $(i+j-a, i)$ matrix entry of $\xi_{t}^{j}$ is $(-1)^{a} h_{a}\left(t_{i}, \ldots, t_{i+j-a}\right)$.) Incorporating the twist by $-\varpi_{k}$, the Lie algebra action on $V_{\varpi_{k}}\left(-\varpi_{k}\right)[t]$ is

$$
\begin{aligned}
\xi_{t}^{j} \cdot \varepsilon_{I}= & \left(\sum_{a=0}^{j}(-1)^{a} h_{a}\left(t_{i_{1}}, \ldots, t_{i_{1}+j-a}\right) \varepsilon_{i_{1}+j-a}\right) \wedge \varepsilon_{i_{2}} \wedge \cdots \wedge \varepsilon_{i_{k}} \\
& +\varepsilon_{i_{1}} \wedge\left(\sum_{a=0}^{j}(-1)^{a} h_{a}\left(t_{i_{2}}, \ldots, t_{i_{2}+j-a}\right) \varepsilon_{i_{2}+j-a}\right) \wedge \cdots \wedge \varepsilon_{i_{k}} \\
& +\cdots+\varepsilon_{i_{1}} \wedge \varepsilon_{i_{2}} \wedge \cdots \wedge\left(\sum_{a=0}^{j}(-1)^{a} h_{a}\left(t_{i_{k}}, \ldots, t_{i_{k}+j-a}\right) \varepsilon_{i_{k}+j-a}\right) \\
& -\left(\left(-t_{1}\right)^{j}+\cdots+\left(-t_{k}\right)^{j}\right) \varepsilon_{I} .
\end{aligned}
$$

For instance,

$$
\begin{aligned}
\xi_{t}^{2} \cdot\left(\varepsilon_{2} \wedge \varepsilon_{3}\right)= & \left(\varepsilon_{4}-\left(t_{2}+t_{3}\right) \varepsilon_{3}+t_{2}^{2} \varepsilon_{2}\right) \wedge \varepsilon_{3}+\varepsilon_{2} \wedge\left(\varepsilon_{5}-\left(t_{3}+t_{4}\right) \varepsilon_{4}+t_{3}^{2} \varepsilon_{3}\right) \\
& -\left(t_{1}^{2}+t_{2}^{2}\right) \varepsilon_{2} \wedge \varepsilon_{3} \\
= & \varepsilon_{2} \wedge \varepsilon_{5}-\varepsilon_{3} \wedge \varepsilon_{4}-\left(t_{3}+t_{4}\right) \varepsilon_{2} \wedge \varepsilon_{4}+\left(t_{3}^{2}-t_{2}^{2}\right) \varepsilon_{2} \wedge \varepsilon_{3}
\end{aligned}
$$

The leading term agrees with the computation $p_{2} \cdot \sigma_{(1,1)}=\sigma_{(3,1)}-\sigma_{(2,2)}$ in $H^{*} \operatorname{Gr}(2, n)$ done in the introduction.

A new feature appears in the equivariant correspondence. Let us pass to the fraction field $\mathbb{C}\left(\mathfrak{t}^{\vee}\right)$, and consider $\mathfrak{g}^{\vee}(t)$, etc., as Lie algebras over this field. Since the element $\xi_{t}$ is regular semisimple, its centralizer $\mathfrak{h} \subseteq \mathfrak{g}^{\vee}(t)$ is a Cartan subalgebra. (In fact, $\mathfrak{h}$ is just the extension of $\mathfrak{a}_{t}$ to $\mathbb{C}\left(\mathfrak{t}^{\vee}\right)$.) So our setup leads naturally to another basis for $V_{\varpi}(-\varpi)(t)$, a basis of weight vectors for $\mathfrak{h}(t)$, diagonalizing $\xi_{t}$.

What is this basis on the geometric side of the correspondence? By the localization theorem (see $[22,(6.3)]$ ), there is a fixed-point basis for $I H_{*}^{T}\left(\mathcal{G} r^{\varpi}\right) \otimes \mathbb{C}\left(\mathfrak{t}^{\vee}\right)$, and in fact this basis corresponds to a (suitably chosen) weight basis for $\mathfrak{h}$.

Theorem 7 (Equivariant Satake, continued) Under the Satake isomorphism $I H_{*}^{T}\left(\overline{\mathcal{G} r^{\varpi}}\right) \cong V_{\varpi}(-\varpi)[t]$, equivariant $M V$-cycles correspond to a weight basis of $V_{\varpi}(-\varpi)[t]$ with respect to $\mathfrak{t}^{\vee}$, and the fixed point basis corresponds to a weight basis with respect to $\mathfrak{h}$.

In general, there is ambiguity in choosing a weight basis. However, for minuscule $\varpi$, all weight spaces of $V_{\varpi}$ are one-dimensional, so a weight basis is determined (up
to scaling) by the Cartan. As noted before, in this case $\mathcal{G} r^{\varpi}=G / P_{\varpi}$ is homogeneous, and the MV basis consists of (opposite) Schubert classes. This is the situation we will consider for the remainder of the paper. Let us write $X=G / P_{\varpi}$.

Let us write $\left\{\sigma_{\lambda}\right\}$ for the basis of Schubert classes in $H_{T}^{*} X$. The fixed point set is $X^{T}=\left\{p_{\lambda}\right\}$ (with $\lambda$ running over the same set indexing Schubert classes), and we will write $\left\{\mathbf{1}_{\lambda}\right\}$ for the corresponding idempotent basis of $H_{T}^{*} X^{T}=\bigoplus H_{T}^{*}\left(p_{\lambda}\right)$. The localization theorem says that the restriction homomorphism (of $\mathbb{C}\left[\mathfrak{t}^{\vee}\right]$-algebras)

$$
\iota^{*}: H_{T}^{*} X \rightarrow H_{T}^{*} X^{T}
$$

becomes an isomorphism after tensoring with $\mathbb{C}\left(\mathfrak{t}^{\vee}\right)$. An important part of equivariant Schubert calculus is to compute the restriction of a Schubert class $\sigma_{\lambda}$ to a fixed point $p_{\mu}$. Formulas for these restrictions have been given by Billey for complete flag varieties [4], and by Ikeda-Naruse, who consider special cases that are related to the focus of this article [24].

The fact that the fixed-point classes form a basis of eigenvectors for $\xi_{t}$ is part of a general phenomenon, with a simple proof. Consider any nonsingular variety $X$ with finite fixed locus $X^{T}$, and any class $\alpha \in\left(H_{T}^{*} X\right) \otimes \mathbb{C}\left(\mathfrak{t}^{\vee}\right)$.

Lemma 8 The idempotent classes $\mathbf{1}_{p} \in\left(H_{T}^{*} X\right) \otimes \mathbb{C}\left(\mathfrak{t}^{\vee}\right)=\left(H_{T}^{*} X^{T}\right) \otimes \mathbb{C}\left(\mathfrak{t}^{\vee}\right)$ form a basis of eigenvectors for the endomorphism $x \mapsto \alpha \cdot x$.

This is almost a tautology. Simply observe that for distinct fixed points $p \neq q$, we have $\mathbf{1}_{p} \cdot \mathbf{1}_{q}=0$. Writing $\alpha=\sum \alpha_{p} \cdot \mathbf{1}_{p}$, the statement follows.

Returning to minuscule Schubert calculus, the restrictions $\left.\sigma_{\lambda}\right|_{\mu}$ may be regarded as matrix entries for the homomorphism $\iota^{*}$, with respect to the Schubert and fixedpoint bases. The Satake correspondence translates the problem of computing this matrix into the following:

Find the change-of-basis matrix relating weight bases of the minuscule representation $V_{\varpi}(-\varpi)(t)$, with respect to two (specific) Cartan subalgebras, $\mathfrak{t}^{\vee}$ and $\mathfrak{h}$, of $\mathfrak{g}^{\vee}(t)$.

This perspective also suggests a framework for setting up and solving the problem of computing restrictions of Schubert classes $\left.\sigma_{\lambda}\right|_{\mu}$. We will work this out in types A and D below.

## 4 A Giambelli Formula for Grassmannians

We will describe a proof of the "equivariant Giambelli formula"

$$
\begin{equation*}
\sigma_{\lambda}=s_{\lambda}(x \mid t) \tag{4}
\end{equation*}
$$

identifying the Schubert class $\sigma_{\lambda} \in H_{T}^{*} G r(k, n)$ with a factorial Schur polynomial, in the spirit of Laksov's approach to equivariant Schubert calculus [30]. The following
definition of the factorial Schur polynomial can be found in Macdonald's book [36, Sect.I.3, Ex. 20]. The (generalized) factorial power is defined as

$$
(x \mid t)^{a}=\left(x+t_{1}\right)\left(x+t_{2}\right) \cdots\left(x+t_{a}\right) .
$$

Let $I=\left\{i_{1}, \ldots, i_{k}\right\} \subseteq[n]$ be the subset corresponding to the partition $\lambda$; recall that this means $\lambda_{k+1-a}=i_{a}-a$. One defines

$$
s_{\lambda}(x \mid t)=\frac{\operatorname{det}\left(\left(x_{j} \mid t\right)^{i-1}\right)_{i \in I, 1 \leq j \leq k}}{\operatorname{det}\left(\left(x_{j} \mid t\right)^{i-1}\right)_{1 \leq i, j \leq k}} .
$$

An easy computation shows the denominator is

$$
\operatorname{det}\left(\left(x_{j} \mid t\right)^{i-1}\right)_{1 \leq i, j \leq k}=\operatorname{det}\left(x_{j}^{i-1}\right)=\prod_{1 \leq a<b \leq k}\left(x_{a}-x_{b}\right)=: \Delta,
$$

the Vandermonde determinant, so the factorial Schur polynomial can also be written as

$$
s_{\lambda}(x \mid t)=\frac{\operatorname{det}\left(\left(x_{j} \mid t\right)^{i-1}\right)_{i \in I, 1 \leq j \leq k}}{\Delta}
$$

The meaning of the Giambelli formula is this. By the localization theorem, the equivariant cohomology of $\operatorname{Gr}(k, n)$ embeds in that of its fixed locus:

$$
H_{T}^{*} G r(k, n) \hookrightarrow H_{T}^{*}\left(G r(k, n)^{T}\right)=\bigoplus_{J} \mathbb{C}[t],
$$

the sum being over all $k$-element subsets $J \subset[n]$. On the other hand, there is a presentation of $H_{T}^{*} \operatorname{Gr}(k, n)$ as a quotient of $\mathbb{C}[t]\left[x_{1}, \ldots, x_{k}\right]^{S_{k}}$ (symmetric polynomials in $x$, with coefficients in $\mathbb{C}[t])$. Composing with the localization homomorphism gives

$$
\mathbb{C}[t]\left[x_{1}, \ldots, x_{k}\right]^{S_{k}} \rightarrow \bigoplus_{J} \mathbb{C}[t]
$$

defined on the $J$ th summand by sending $x_{a} \mapsto-t_{j_{a}}$. The precise statement is this:
Theorem 9 Under the homomorphism $\mathbb{C}[t]\left[x_{1}, \ldots, x_{k}\right]^{S_{k}} \rightarrow H_{T}^{*} \operatorname{Gr}(k, n)$, we have $s_{\lambda}(x \mid t) \mapsto \sigma_{\lambda}$. Equivalently, for each $J=\left\{j_{1}<\cdots<j_{k}\right\}$, we have

$$
\left.\sigma_{\lambda}\right|_{p_{J}}=s_{\lambda}\left(-t_{j_{1}}, \ldots,-t_{j_{k}} \mid t\right) .
$$

Proof There are three simple steps. We describe them informally first, since we will follow the same pattern in proving a type D formula later.
(1) Work out the case $k=1$, corresponding to projective space. Here, by construction, the element $\xi_{t}$ corresponds to multiplication by the hyperplane class $\sigma_{\square}$
on $H_{T}^{*} \mathbb{P}^{n-1}$, written in the Schubert basis $\varepsilon_{i}=\left[\mathbb{P}^{n-i}\right]$. We choose a basis $\bar{f}_{i}$ diagonalizing the semisimple element $\xi_{t}$; by Lemma 8, this basis coincides with the basis of idempotents $\mathbf{1}_{i}$, up to scalar. We normalize the $\bar{f}_{i}$ so that $\bar{f}_{i}=\mathbf{1}_{i}$, by requiring $\varepsilon_{1}=\bar{f}_{1} \pm \cdots+\bar{f}_{n}$ (since $\varepsilon_{1}$ corresponds to $\mathbf{1} \in H_{T}^{*} \mathbb{P}^{n-1}$ ). The expansion of $\varepsilon_{i}$ in the $\bar{f}_{j}$ basis is then a localization calculation, which is easy for projective space.
(2) For each $k>1$, take the weight basis $\left\{\varepsilon_{I}\right\}$ of $V_{\varpi_{k}}\left(-\varpi_{k}\right)[t]=\bigwedge_{\mathbb{C}[t]}^{k} \mathbb{C}[t]^{n}$ to be $\varepsilon_{I}=\varepsilon_{i_{1}} \wedge \cdots \wedge \varepsilon_{i_{k}}$. Verify that the action of $\xi_{t}$ agrees with the known formula for multiplication by $\sigma_{\square}$ on the Schubert basis of $H_{T}^{*} \operatorname{Gr}(k, n)$, so that we can identify $\sigma_{I}=\varepsilon_{I}$. (The latter formula is often called the equivariant Chevalley formula.)
(3) By Lemma 8 again, the vectors $\bar{f}_{j_{1}} \wedge \cdots \wedge \bar{f}_{i_{k}}$ agree with the basis of idempotents $\mathbf{1}_{J}$, up to scalar; normalize it so that $\bar{f}_{J}=\mathbf{1}_{J}$ by requiring $\varepsilon_{\{1, \ldots, k\}}=\sum \bar{f}_{J}$ (since $\varepsilon_{\{1, \ldots, k\}}$ corresponds to $\mathbf{1} \in H_{T}^{*} \operatorname{Gr}(k, n)$ ). On the other hand, formulas from Step (1) expressing $\varepsilon_{i}$ in terms of $\bar{f}_{j}$ yield (determinantal) formulas for $\varepsilon_{I}$ in terms of $\bar{f}_{j_{1}} \wedge \cdots \wedge \bar{f}_{j_{k}}$; comparing with the normalized vectors $\bar{f}_{J}$ proves the theorem.

Now we proceed to work this out in detail. It is not hard to see that

$$
f_{i}=\varepsilon_{i}+\frac{1}{\left(t_{i+1}-t_{i}\right)} \varepsilon_{i+1}+\cdots+\frac{1}{\left(t_{n}-t_{i}\right) \cdots\left(t_{i+1}-t_{i}\right)} \varepsilon_{n}
$$

is a basis of eigenvectors for $\xi_{t}$ acting on $\mathbb{C}^{n} \otimes \mathbb{C}(t)$ (inverting nonzero characters). This is related to $\varepsilon_{i}$ by a unitriangular change of basis. However, note that $\left.\sigma_{i}\right|_{p_{i}}=$ $\left(t_{1}-t_{i}\right) \cdots\left(t_{i-1}-t_{i}\right)$ (since $\Omega_{i}=\mathbb{P}^{n-i}$ is defined by the vanishing of the first $i-1$ coordinates); this means that we must rescale to obtain the idempotent basis. In fact,

$$
\bar{f}_{i}=\frac{1}{\left(t_{1}-t_{i}\right) \cdots\left(t_{i-1}-t_{i}\right)} f_{i}
$$

identifies with the idempotent basis $\mathbf{1}_{i} \in H_{T}^{*}\left(\mathbb{P}^{n-1}\right)^{T}$. Since we know the restrictions of $\varepsilon_{i}=\left[\mathbb{P}^{n-i}\right]$ to fixed points, we see

$$
\begin{aligned}
\varepsilon_{i} & =\sum_{j}\left(t_{1}-t_{j}\right) \cdots\left(t_{i-1}-t_{j}\right) \bar{f}_{j} \\
& =\left.\sum_{j}\left(x_{j} \mid t\right)^{i-1}\right|_{x_{j}=-t_{j}} \bar{f}_{j}
\end{aligned}
$$

using the generalized factorial power notation. This completes the first step.
For the second step, we take $\varepsilon_{I}=\varepsilon_{i_{1}} \wedge \cdots \wedge \varepsilon_{i_{k}}$ as our basis for $V_{\varpi_{k}}\left(-\varpi_{k}\right)[t]=$ $H_{T}^{*} G r(k, n)$. The verification that $\xi_{t}$ acts on this basis as $\sigma_{\square}$ does on the Schubert basis is left to the reader. (Illustrative examples were done above.) We note that the Chevalley formula says

$$
\sigma_{\square} \cdot \sigma_{I}=\sum_{I^{+}} \sigma_{I^{+}}+\left(t_{1}+\cdots+t_{k}-t_{i_{1}}-\cdots-t_{i_{k}}\right) \sigma_{I},
$$

where the sum is over $I^{+}$obtained from $I$ by replacing some $i_{a} \in I$ such that $i_{a}+1 \notin$ $I$ by $i_{a}+1$. (In terms of the corresponding partitions, $\lambda\left(I^{+}\right)$is obtained from $\lambda(I)$ by adding a single box.)

Finally, from the definition of exterior product, we get

$$
\varepsilon_{I}:=\varepsilon_{i_{1}} \wedge \cdots \wedge \varepsilon_{i_{k}}=\sum_{J} \operatorname{det}\left(\left.\left(x_{j} \mid t\right)^{i-1}\right|_{x_{j}=-t_{j}}\right)_{i \in I, j \in J} \bar{f}_{j_{1}} \wedge \cdots \wedge \bar{f}_{j_{k}}
$$

In particular,

$$
\begin{aligned}
\varepsilon_{\{1, \ldots, k\}} & =\sum_{J} \operatorname{det}\left(\left.\left(x_{j} \mid t\right)^{i-1}\right|_{x_{j}=-t_{j}}\right)_{1 \leq i \leq k, j \in J} \bar{f}_{j_{1}} \wedge \cdots \wedge \bar{f}_{j_{k}} \\
& =\sum_{J} \Delta_{J} \bar{f}_{j_{1}} \wedge \cdots \wedge \bar{f}_{j_{k}},
\end{aligned}
$$

where

$$
\Delta_{J}=\prod_{1 \leq a<b \leq k}\left(t_{j_{a}}-t_{j_{b}}\right)
$$

is the specialization of the Vandermonde determinant $\Delta$. Since $\varepsilon_{\{1, \ldots, k\}}$ should be identified with $\mathbf{1}=[G r(k, n)]$ in $H_{T}^{*} G r(k, n)$, this tells us that the idempotent classes are

$$
\mathbf{1}_{J}=\bar{f}_{J}:=\Delta_{J} \bar{f}_{j_{1}} \wedge \cdots \wedge \bar{f}_{j_{k}}
$$

and we can rewrite the above formula as

$$
\begin{aligned}
\varepsilon_{I}:=\varepsilon_{i_{1}} \wedge \cdots \wedge \varepsilon_{i_{k}} & =\sum_{J}\left(\frac{\operatorname{det}\left(\left.\left(x_{j} \mid t\right)^{i-1}\right|_{x_{j}=-t_{j}}\right)_{i \in I, j \in J}}{\Delta_{J}}\right) \bar{f}_{J} \\
& =\left.\sum_{J} s_{\lambda}\left(x_{1}, \ldots, x_{k} \mid t\right)\right|_{x_{a}=-t_{j a}} \bar{f}_{J},
\end{aligned}
$$

as required.
Remark 10 It is hard to give clear attribution to the equivariant Giambelli formula; certainly it was known by around 2000. Reference to it appears in [28], and a proof is in [40]. In retrospect, the Kempf-Laksov formula [27] is equivalent to (4). See also [1] for more discussion and an alternative proof.

Likewise, it is difficult to identify the earliest appearance of the connection between the $\mathfrak{g l}_{n}$-module $\bigwedge^{k} \mathbb{C}^{n}$ and the cohomology of $H^{*} \operatorname{Gr}(k, n)$. While surely known long before, it appears in several sources by the 2000s [14, 26, 30, 48].

Remark 11 Recall that the Grassmannian $\operatorname{Gr}(k, n)=\operatorname{Gr}\left(k, V^{*}\right)$ embeds naturally in $\mathbb{P}\left(\bigwedge^{k} V^{*}\right)$ as the locus of "pure wedges" $v_{1} \wedge \cdots \wedge v_{k}$. Since our basis $\varepsilon_{1}, \ldots, \varepsilon_{n}$ for $V$ is a weight basis for $\mathfrak{t}^{\vee}$, the dual basis $\varepsilon_{1}^{*}, \ldots, \varepsilon_{n}^{*}$ is a weight basis for the action of $T$ on $V^{*}$; the points of $\operatorname{Gr}\left(k, V^{*}\right)$ corresponding to $\varepsilon_{I}^{*}=\varepsilon_{i_{1}}^{*} \wedge \cdots \wedge e_{i_{k}}^{*}$ are therefore precisely the $T$-fixed points. Thus the Satake correspondence exchanges $T$ fixed points in $\operatorname{Gr}\left(k, V^{*}\right) \subseteq \mathbb{P}\left(\bigwedge^{k} V^{*}\right)$ with Schubert classes in $\mathbb{P}\left(H^{*} \operatorname{Gr}(k, V)\right)=$ $\mathbb{P}\left(\bigwedge^{k} V\right)$. When considered as cohomology classes on $\operatorname{Gr}\left(k, V^{*}\right)$, do pure wedges in $\mathbb{P}\left(\bigwedge^{k} V\right)$ have a natural geometric meaning? What does this correspondence look like when upgraded to the equivariant setting?

## 5 Quadrics

As another example, we describe the correspondence for minuscule varieties of type $D$. To set things up, fix a basis

$$
\varepsilon_{\overline{n-1}}, \ldots, \varepsilon_{\overline{1}}, \varepsilon_{\overline{0}}, \varepsilon_{0}, \varepsilon_{1}, \ldots, \varepsilon_{n-1}
$$

for $V \cong \mathbb{C}^{2 n}$, and equip this vector space with the symmetric bilinear form defined by $\left\langle\varepsilon_{\bar{\imath}}, \varepsilon_{j}\right\rangle=\delta_{i, j}$. (The barred indices should be regarded as notation for negative integers.) The form identifies the dual basis for $V^{*}=V$ as $\varepsilon_{i}^{*}=\varepsilon_{\bar{l}}$. The subspaces $E=\operatorname{span}\left\{\varepsilon_{0}, \ldots, \varepsilon_{n-1}\right\}$ and $E=\operatorname{span}\left\{\varepsilon_{\overline{n-1}}, \ldots, \varepsilon_{\overline{0}}\right\}$ are maximal isotropic subspaces of $V$, and we have $V=\bar{E} \oplus E$.

We will take $\mathfrak{s o}_{2 n} \subseteq \mathfrak{s l}_{2 n}$ to be the algebra preserving the given bilinear form; there is also a canonical identification $\mathfrak{s o}_{2 n}=\bigwedge^{2} V$ (see, e.g., [12, Sect. 20]). We take our principal nilpotent element $\xi$ and generic $t$ so that
$\operatorname{Using} \varepsilon_{\bar{l}}^{*}=\varepsilon_{i}, \xi$ can also be written as

$$
\begin{aligned}
\xi & =-\sum_{i=1}^{n-1} \varepsilon_{\bar{l}}^{*} \otimes \varepsilon_{i-1}-\varepsilon_{\overline{1}}^{*} \otimes \varepsilon_{0}+\varepsilon_{\overline{0}}^{*} \otimes \varepsilon_{1}+\sum_{i=1}^{n-1} \varepsilon_{i-1}^{*} \otimes \varepsilon_{i} \\
& =\varepsilon_{0} \wedge \varepsilon_{1}+\sum_{i=1}^{n-1} \varepsilon_{\overline{i-1}} \wedge \varepsilon_{i},
\end{aligned}
$$

which exhibits it as an element of $\bigwedge^{2} V$. Similarly, we have

$$
\begin{equation*}
\xi_{t}=\varepsilon_{0} \wedge \varepsilon_{1}+\sum_{i=1}^{n-1} \varepsilon_{\overline{i-1}} \wedge \varepsilon_{i}-\sum_{i=0}^{n-1} t_{i} \varepsilon_{\bar{l}} \wedge \varepsilon_{i} \tag{6}
\end{equation*}
$$

Note that our indexing conventions and choice of form $\langle$,$\rangle make it natural to identify$ elements of $\bigwedge^{2} V$ with matrices which are skew-symmetric about the anti-diagonal.

The odd matrix powers $\xi_{t}, \xi_{t}^{3}, \ldots, \xi_{t}^{2 n-3}$ all lie in $\mathfrak{s o}_{2 n}$ as well, and they are easily seen to be linearly independent elements of the centralizer $\mathfrak{a}_{t}$ of $\xi_{t}$. Since $\xi_{t}$ is regular, one knows $\operatorname{dim} \mathfrak{a}_{t}=n$; the missing element is

$$
\begin{equation*}
\eta_{t}=-\sum_{j=1}^{n-1}\left(t_{j+1} \cdots t_{n-1}\right) \varepsilon_{0} \wedge \varepsilon_{j}+\sum_{0 \leq i \leq j \leq n-1}\left(t_{0} \cdots t_{i-1} t_{j+1} \cdots t_{n-1}\right) \varepsilon_{\bar{\imath}} \wedge \varepsilon_{j} . \tag{7}
\end{equation*}
$$

For example, when $n=4$, this is

$$
\eta_{t}=\left[\begin{array}{cccccccc}
-t_{0} t_{1} t_{2} & & & & & & & \\
-t_{0} t_{1} & -t_{0} t_{1} t_{3} & & & & & & \\
-t_{0} & -t_{0} t_{3} & -t_{0} t_{2} t_{3} & & & & & \\
-1 & -t_{3} & -t_{2} t_{3} & -t_{1} t_{2} t_{3} & & & & \\
1 & t_{3} & t_{2} t_{3} & 0 & t_{1} t_{2} t_{3} & & & \\
0 & 0 & 0 & -t_{2} t_{3} & t_{2} t_{3} & t_{0} t_{2} t_{3} & & \\
0 & 0 & 0 & -t_{3} & t_{3} & t_{0} t_{3} & t_{0} t_{1} t_{3} & \\
0 & 0 & 0 & -1 & 1 & t_{0} & t_{0} t_{1} & t_{0} t_{1} t_{2}
\end{array}\right] .
$$

When discussing homogeneous spaces, we will assume $n \geq 3$ to avoid setting conventions for special cases. ${ }^{2}$ Consider the ( $2 n-2$ )-dimensional quadric $\mathcal{Q}=\mathcal{Q}^{2 n-2} \subseteq \mathbb{P}\left(V^{*}\right) \cong \mathbb{P}^{2 n-1}$ of isotropic vectors for the given bilinear form. (In coordinates, $\mathcal{Q}$ is defined by the vanishing of the quadratic form $\sum_{i=0}^{n-1} X_{\bar{\imath}} X_{i}$, where $X_{i}=\varepsilon_{i}^{*}$.) The torus $T \cong\left(\mathbb{C}^{*}\right)^{n}$ acts on $V$ with weights $-t_{n-1}, \ldots,-t_{0}, t_{0}, \ldots, t_{n-1}$, inducing an action on $\mathcal{Q}$.

The quadric $\mathcal{Q}$ is homogeneous for $\mathrm{PSO}_{2 n}$, and the Satake correspondence identifies $H_{T}^{*} \mathcal{Q}$ with $V_{\varpi_{1}}\left(-\varpi_{1}\right)[t]$, where $\varpi_{1}=-t_{n-1}$ and $V_{\varpi_{1}}=V$ is the standard

[^12]representation of $\mathfrak{S o}_{2 n}$. To see this explicitly, define Schubert varieties in $\mathcal{Q}$ by
\[

$$
\begin{aligned}
& \Omega_{\bar{l}}=\left\{\varepsilon_{\overline{n-1}}^{*}=\cdots=\varepsilon_{\overline{i+1}}^{*}=0\right\} \\
& \Omega_{i}=\left\{\varepsilon_{\frac{*}{n-1}}^{*}=\cdots=\varepsilon_{\overline{0}}^{*}=\varepsilon_{0}^{*}=\cdots=\varepsilon_{i-1}^{*}=0\right\}
\end{aligned}
$$
\]

for $i>0$; and

$$
\begin{aligned}
& \Omega_{\overline{0}}=\left\{\varepsilon_{\frac{*}{n-1}}^{*}=\cdots=\varepsilon_{\overline{1}}^{*}=\varepsilon_{\overline{0}}^{*}=0\right\} \\
& \Omega_{0}=\left\{\varepsilon_{\frac{*}{n-1}}^{*}=\cdots=\varepsilon_{\overline{1}}^{*}=\varepsilon_{0}^{*}=0\right\}
\end{aligned}
$$

Identify the Schubert classes $\sigma_{i}=\left[\Omega_{i}\right]$ in $H_{n-1-i}^{T} \mathcal{Q}=H_{T}^{n-1+i} \mathcal{Q}$ with basis elements $\varepsilon_{i}$ by

$$
\begin{equation*}
\sigma_{\bar{\iota}}=(-1)^{i} \varepsilon_{\bar{\iota}} \quad \text { and } \quad \sigma_{i}=\varepsilon_{i} \tag{8}
\end{equation*}
$$

for $i \geq 0$. Using the twist by $t_{n-1}$, we have

$$
\begin{aligned}
& \xi_{t} \cdot \varepsilon_{\bar{l}}=-\varepsilon_{\bar{\imath}+1}-\left(-t_{i}+t_{n-1}\right) \varepsilon_{\bar{l}} \quad \text { for } i>1, \\
& \xi_{t} \cdot \varepsilon_{\overline{1}}=-\varepsilon_{\overline{0}}-\varepsilon_{0}-\left(-t_{1}+t_{n-1}\right) \varepsilon_{\overline{1}}, \\
& \xi_{t} \cdot \varepsilon_{\overline{0}}=\varepsilon_{1}-\left(-t_{0}+t_{n-1}\right) \varepsilon_{\overline{0}}, \quad \text { and } \\
& \xi_{t} \cdot \varepsilon_{i}=\varepsilon_{i+1}-\left(t_{i}+t_{n-1}\right) \varepsilon_{i} \quad \text { for } i \geq 0 .
\end{aligned}
$$

On the other hand, taking $\sigma_{\square}=\sigma_{\overline{n-2}} \in H_{T}^{1} \mathcal{Q}$ to be the hyperplane class,

$$
\begin{aligned}
\sigma_{\square} \cdot \sigma_{\bar{\imath}} & =\sigma_{\bar{\imath}+1}+\left(t_{i}-t_{n-1}\right) \sigma_{\bar{\imath}} \quad \text { for } i>1, \\
\sigma_{\square} \cdot \sigma_{\overline{1}} & =\sigma_{\overline{0}}+\sigma_{0}+\left(t_{1}-t_{n-1}\right) \sigma_{\overline{1}}, \\
\sigma_{\square} \cdot \sigma_{\overline{0}} & =\sigma_{1}+\left(t_{0}-t_{n-1}\right) \sigma_{\overline{0}}, \quad \text { and } \\
\sigma_{\square} \cdot \sigma_{i} & =\sigma_{i+1}+\left(-t_{i}-t_{n-1}\right) \sigma_{i} \quad \text { for } i \geq 0,
\end{aligned}
$$

so (8) compatibly identifies the action of $\xi_{t}$ with the product by $\sigma_{\square}$.
Iterating the above computation of $\xi_{t} \cdot \varepsilon_{i}$ leads to a formula for the matrix entries of $\xi_{t}^{2 j-1}$, for $j=1, \ldots, n-1$ :

$$
\begin{align*}
\xi_{t}^{2 j-1}= & \sum_{k=0}^{n-1} \sum_{i=0}^{n-1-k}(-1)^{2 j-1-k} h_{2 j-1-k}\left(t_{i}, \ldots, t_{i+k}\right) \varepsilon_{\bar{\imath}} \wedge \varepsilon_{i+k} \\
& +\sum_{k=1}^{n-1}(-1)^{2 j-1-k} h_{2 j-1-k}\left(-t_{0}, t_{1}, \ldots, t_{k}\right) \varepsilon_{0} \wedge \varepsilon_{k} \\
& +2 \sum_{i=1}^{n-2} \sum_{k=i+1}^{n-1}(-1)^{2 j-1-k} h_{2 j-1-i-k}\left(-t_{i},-t_{i-1}, \ldots,-t_{0}, t_{0}, \ldots, t_{k}\right) \varepsilon_{i} \wedge \varepsilon_{k}, \tag{9}
\end{align*}
$$

where in the last sum, the complete homogeneous symmetric functions are in $i+$ $k+2$ variables, specialized as indicated to consecutive $t$ 's.

The fixed points in the quadric $\mathcal{Q}$ are the $2 n$ coordinate points:

$$
\mathcal{Q}^{T}=\left\{p_{\overline{n-1}}, \ldots, p_{\overline{0}}, p_{0}, \ldots, p_{n-1} .\right\}
$$

That is, $p_{j} \in \mathcal{Q} \subseteq \mathbb{P}^{2 n-1}$ is the point with 1 in the $j$ th coordinate and 0 elsewhere. For each $i$, one computes the weights of $T$ acting on the tangent space $T_{p_{i}} \mathcal{Q}$ to be $\left\{t_{j}-t_{i} \mid j \neq i, \bar{\imath}\right\}$, using the notation $t_{\bar{J}}=-t_{j}$.

Using the defining equations, it is easy to write down the restrictions of Schubert classes to fixed points. For $i>0$, and any $j$, we have

$$
\begin{aligned}
\left.\sigma_{\bar{\imath}}\right|_{p_{j}} & =\left(-t_{n-1}-t_{j}\right) \cdots\left(-t_{i+1}-t_{j}\right) \\
& =\prod_{k=i+1}^{n-1}\left(-t_{k}-t_{j}\right),
\end{aligned}
$$

Thus $\left.\sigma_{\bar{\imath}}\right|_{p_{j}}=0$ for $j<\bar{\imath}$, since the factor $\left(-t_{\bar{J}}-t_{j}\right)$ is zero-and indeed, in this case $p_{j} \notin \Omega_{\bar{i}}$.

For $i=\overline{0}$, we have

$$
\begin{aligned}
& \left.\sigma_{\overline{0}}\right|_{p_{\overline{0}}}=\prod_{k=1}^{n-1}\left(-t_{k}+t_{0}\right) ; \\
& \left.\sigma_{\overline{0}}\right|_{p_{j}}=\left(t_{0}-t_{j}\right) \prod_{\substack{k=0 \\
k \neq j}}^{n-1}\left(-t_{k}-t_{j}\right) \quad \text { for } j>0 ;
\end{aligned}
$$

and all other restrictions are zero. Finally, for $i \geq 0$,

$$
\left.\sigma_{i}\right|_{p_{j}}=\left(\prod_{\substack{k=0 \\ k \neq j}}^{n-1}\left(-t_{k}-t_{j}\right)\right) \prod_{k=0}^{i-1}\left(t_{k}-t_{j}\right) \quad \text { for } j \geq i
$$

and all other restrictions are zero.
Finally, we translate these calculations into a change of bases for the representation $V$. Using $\varepsilon_{\bar{l}}=(-1)^{i} \sigma_{\bar{l}}$ and $\varepsilon_{i}=\sigma_{i}$ for $i \geq 0$, we have $[\mathcal{Q}]=\sigma_{\overline{n-1}}=(-1)^{n-1} \varepsilon_{\overline{n-1}}$, so by setting $\bar{f}_{i}=(-1)^{n-1} \mathbf{1}_{i}$, where $\mathbf{1}_{i} \in H_{T}^{*} \mathcal{Q}^{T}$ is the idempotent class at $p_{i}$, we have

$$
\varepsilon_{\overline{n-1}}=\bar{f}_{\overline{n-1}}+\cdots+\bar{f}_{\overline{0}}+\bar{f}_{0}+\cdots+\bar{f}_{n-1} .
$$

More generally, for $i>0$ we have

$$
\varepsilon_{\bar{l}}=\sum_{j=\bar{l}}^{n-1}\left(\prod_{k=i+1}^{n-1}\left(t_{k}+t_{j}\right)\right) \bar{f}_{j}
$$

for $i=\overline{0}$ we have

$$
\varepsilon_{\overline{0}}=\left(\prod_{k=1}^{n-1}\left(t_{k}-t_{0}\right)\right) \bar{f}_{\overline{0}}+\sum_{j=1}^{n-1}\left(\left(t_{j}-t_{0}\right) \prod_{k>0, k \neq j}\left(t_{k}+t_{j}\right)\right) \bar{f}_{j}
$$

and for $i \geq 0$ we have

$$
\varepsilon_{i}=\sum_{j=i}^{n-1}\left(\prod_{k=0}^{i-1}\left(t_{k}^{2}-t_{j}^{2}\right) \prod_{\substack{k=i \\ k \neq j}}^{n-1}\left(t_{k}+t_{j}\right)\right) \bar{f}_{j}
$$

It will be convenient to rescale the basis $\left\{\bar{f}_{i}\right\}$ so that $\left\{\varepsilon_{i}\right\}$ is related by a unitriangular change of basis. To this end, for each $i$ let

$$
f_{i}=\alpha_{i} \bar{f}_{i}
$$

where, for $i \geq 0$, the scaling coefficients are $\alpha_{\bar{l}}=\left.(-1)^{n-1-i} \sigma_{\bar{l}}\right|_{p_{\bar{\tau}}}$ and $\alpha_{i}=(-1)^{n-1}$ $\left.\sigma_{i}\right|_{p_{i}}$. Now we may write, for $i \geq 0$,

$$
\varepsilon_{\bar{\imath}}=\sum_{j=0}^{i} \bar{c}_{j i} f_{\bar{J}}+\sum_{j=0}^{n-1} c_{j i} f_{j},
$$

and

$$
\varepsilon_{i}=\sum_{j=0}^{i} \bar{b}_{j i} f_{j}
$$

Explicitly, the matrices $\bar{C}=\left(\bar{c}_{j i}\right), C=\left(c_{j i}\right)$, and $\bar{B}=\left(\bar{b}_{j i}\right)$ are computed as follows. For $i>0$,

$$
\begin{aligned}
& \bar{c}_{j i}=\frac{1}{\prod_{k=j+1}^{i}\left(t_{k}-t_{j}\right)} \quad \text { for } 0 \leq j \leq i ; \\
& c_{j i}=\frac{1}{\prod_{k=0}^{j-1}\left(t_{k}^{2}-t_{j}^{2}\right) \prod_{k=j+1}^{i-1}\left(t_{k}+t_{j}\right)} \quad \text { for } 0 \leq j \leq i ;
\end{aligned}
$$

and

$$
c_{j i}=\frac{2 t_{j}}{\prod_{k=0}^{i}\left(t_{k}^{2}-t_{j}^{2}\right) \prod_{k=i+1}^{j-1}\left(t_{k}-t_{j}\right)} \quad \text { for } j>i .
$$

For $i=0$,

$$
\begin{aligned}
& \bar{c}_{00}=1 ; \\
& c_{00}=0 ;
\end{aligned}
$$

and

$$
c_{j 0}=\frac{1}{\left(-t_{0}-t_{j}\right) \prod_{k=1}^{j-1}\left(t_{k}-t_{j}\right)} .
$$

Finally, for $i \geq 0$,

$$
\bar{b}_{j i}=\frac{1}{\prod_{k=i}^{j-1}\left(t_{k}-t_{j}\right)} .
$$

Note the matrices $\bar{C}$ and $\bar{B}$ are indeed unitriangular.

## 6 Orthogonal Grassmannians

Now we turn to the maximal orthogonal Grassmannians $O G^{ \pm}(n, 2 n)$, also known as spinor varieties. We will maintain the notation from the previous section, so $V=\mathbb{C}^{2 n}$ has a bilinear form and basis $\varepsilon_{i}$ so that $\left\langle\varepsilon_{\bar{i}}, \varepsilon_{j}\right\rangle=\delta_{i j}$. As noted above, the subspace $E \subseteq V$ spanned by $\varepsilon_{0}, \varepsilon_{1}, \ldots, \varepsilon_{n-1}$ is isotropic with respect to the bilinear form, as is the complementary subspace $\bar{E}=\operatorname{span}\left\{\varepsilon_{\overline{0}}, \ldots, \varepsilon_{\overline{n-1}}\right\}$. The orthogonal Grassmannian $O G^{+}(n, 2 n)$ (respectively, $O G^{-}(n, 2 n)$ ) parametrizes all $n$-dimensional isotropic subspaces $L \subseteq V$ such that $\operatorname{dim}(E \cap L)$ is even (resp., odd). We will focus on the " + " case, and write $O G(n)=O G^{+}(n, 2 n)$ from now on.

The torus $T=\left(\mathbb{C}^{*}\right)^{n}$ acts on $O G(n)$ via its action on $V=\mathbb{C}^{2 n}$; recall that this is given by weights $-t_{n-1}, \ldots,-t_{0}, t_{0}, \ldots, t_{n-1}$. The $T$-fixed points in $O G(n)$ are indexed by subsets $I \subseteq\{0, \ldots, n-1\}$ such that the cardinality of $I$ is even. For such a subset, the fixed point $p_{I}$ corresponds to the subspace

$$
E_{I}=\operatorname{span}\left(\left\{\varepsilon_{i} \mid i \in I\right\} \cup\left\{\varepsilon_{\bar{J}} \mid j \notin I\right\}\right)
$$

For example, $p_{\emptyset}=\bar{E}$.
Schubert varieties in $O G(n)$ are also indexed by subsets $I \subseteq\{0, \ldots, n-1\}$ of even cardinality. As for the ordinary Grassmannian, the elements of $I$ index pivots for Schubert cells $\Omega_{I}^{\circ}$, and $\Omega_{I}$ is the closure. The isotropicity conditions mean that exactly one of $i$ or $\bar{l}$ occurs as a pivot, for $0 \leq i \leq n-1$, and we record the positive ones. For example, in $O G(4)$ we have

$$
\Omega_{\{1,3\}}^{\circ}=\left[\begin{array}{cccccccc}
0 & 1 & * & 0 & * & 0 & \bullet & 0 \\
0 & 0 & 0 & 1 & \bullet & 0 & \bullet & 0 \\
0 & 0 & 0 & 0 & 0 & 1 & \bullet & 0 \\
0 & 0 & 0 & 0 & 0 & 0 & 0 & 1
\end{array}\right]
$$

(From left to right, the columns are numbered $\overline{3}, \overline{2}, \overline{1}, \overline{0}, 0,1,2,3$. Stars are free entries, and bullets indicate entries that are dependent on the others, by the isotropicity condition.) Similarly, $\Omega_{\emptyset}=O G(n)$, and $\Omega_{\{(0), 1, \ldots, n-1\}}=\left\{p_{\{(0), 1, \ldots, n-1\}}\right\}$, where 0 is included or not, depending on the parity of $n$. Schubert varieties are $T$-invariant.

Frequently one interprets the subsets $I$ as strict partitions $\lambda$, simply by reversing order from increasing to decreasing. In this context, we will usually prefer the subset notation to partition notation, although the latter is useful for indicating containment relations: if $I$ and $J$ are subsets corresponding to partitions $\lambda$ and $\mu$, respectively, then $\Omega_{I} \subseteq \Omega_{J}$ if and only if $\lambda \supseteq \mu$ as Young diagrams. We will write $I \geq J$ in this case. (If $I=\left\{i_{1}<\cdots<i_{r}\right\}$ and $J=\left\{j_{1}<\cdots<j_{s}\right\}$, then $I \geq J$ is equivalent to $r \geq s$ and $i_{a} \geq j_{a}$ for $1 \leq a \leq s$.) We write $\sigma_{I}=\left[\Omega_{I}\right]$ for the equivariant class of a Schubert variety; it has degree $|I|:=\sum i_{a}$.

Our main goal in this section is to compute formulas for the restrictions $\left.\sigma_{I}\right|_{p_{J}}$. Since $p_{I}$ is the unique fixed point in the Schubert cell $\Omega_{I}^{\circ}$, we have $p_{I} \in \Omega_{J}$ if and only if $I \geq J$. From matrix representatives, it is easy to see that the normal space to $\Omega_{I}^{\circ} \subset O G(n)$ at the point $p_{I}$ has weights $\left\{-t_{i}+t_{j} \mid i>j ; i \in I, j \notin I\right\} \cup\left\{-t_{i}-\right.$ $\left.t_{j} \mid i>j ; i, j \in I\right\}$. It follows that

$$
\begin{equation*}
\left.\sigma_{I}\right|_{p_{I}}=\prod_{i \in I}\left(\prod_{\substack{j \notin I \\ j<i}}\left(-t_{i}+t_{j}\right) \prod_{\substack{j \in I \\ j<i}}\left(-t_{i}-t_{j}\right)\right) \tag{10}
\end{equation*}
$$

The corresponding minuscule representation is the half-spin representation $\mathbb{S}^{+}$of $\mathfrak{s o}_{2 n}$. A brief description, suitable for our purposes, is in the appendix; to see this worked out in detail, we recommend [12, Sect. 20], [10], or [37].

Recall our standard representation $V$ of $\mathfrak{s o}_{2 n}$ splits into maximal isotropic subspaces $V=\bar{E} \oplus E$, and we have fixed a basis $\varepsilon_{i}$ so that $\varepsilon_{\overline{n-1}}, \ldots, \varepsilon_{\overline{0}}$ span $\bar{E}$, and $e_{0}, \ldots, e_{n-1}$ span $E$. As noted in the appendix, $\mathbb{S}=\mathbb{S}^{+} \oplus \mathbb{S}^{-}$is an ideal of the Clifford algebra $C l(V)$, and $\mathbb{S}^{+}$has a basis of elements

$$
\begin{equation*}
\varepsilon_{I}:=\varepsilon_{\bar{l}_{1}^{\prime}} \cdots \varepsilon_{\bar{l}_{n-r}^{\prime}} \cdot \varepsilon, \tag{11}
\end{equation*}
$$

for $I=\left\{i_{1}<\cdots<i_{r}\right\} \subseteq\{0, \ldots, n-1\}$ of even cardinality, with complement $I^{\prime}=$ $\left\{i_{1}^{\prime}<\cdots<i_{n-r}^{\prime}\right\}$. Here $\varepsilon=\varepsilon_{0} \cdots \varepsilon_{n-1}$.

The Cartan subalgebra is spanned by $\varepsilon_{\bar{\imath}} \wedge \varepsilon_{i}$, and its eigenvalue on the weight vector $\varepsilon_{I}$ is computed to be

$$
\frac{1}{2}\left(\sum_{j \notin I} t_{j}-\sum_{i \in I} t_{i}\right)
$$

In particular, the highest weight vector $\varepsilon_{\emptyset}$ has weight $\varpi_{n}=\frac{1}{2} \sum_{j} t_{j}$. Twisting by $-\varpi_{n}$, the action of $\mathfrak{t}^{\vee}$ on $V_{\varpi_{n}}\left(-\varpi_{n}\right)[t]$ has $t \cdot \varepsilon_{I}=\left(-\sum_{i \in I} t_{i}\right) \varepsilon_{I}$. Straightforward computations also show

$$
\begin{aligned}
\varepsilon_{\bar{l}} \wedge \varepsilon_{i+1} \cdot \varepsilon_{I} & =\varepsilon_{\bar{l}} \cdot \varepsilon_{i+1} \cdot \varepsilon_{I} \\
& = \begin{cases}\varepsilon_{I^{+}} & \text {if } i \in I \text { and } i+1 \notin I \\
0 & \text { otherwise }\end{cases}
\end{aligned}
$$

here $I^{+}=(I \backslash\{i\}) \cup\{i+1\}$. Similarly, $\varepsilon_{0} \wedge \varepsilon_{1} \cdot \varepsilon_{I}=\varepsilon_{I \cup\{0,1\}}$ if $0,1 \notin I$, and is zero otherwise.

This is enough to compute the action of $\xi_{t}$ on $\mathbb{S}^{+}$. For example,

$$
\begin{aligned}
\xi_{t} \cdot \varepsilon_{\emptyset} & =\varepsilon_{\{0,1\}} ; \\
\xi_{t} \cdot \varepsilon_{\{0,1\}} & =\left(\varepsilon_{\overline{1}} \wedge \varepsilon_{2}\right) \cdot \varepsilon_{\{0,1\}}-t \cdot \varepsilon_{\{0,1\}} \\
& =\varepsilon_{\{0,2\}}+\left(t_{0}+t_{1}\right) \varepsilon_{\{0,1\}} ; \\
\xi_{t} \cdot \varepsilon_{\{0,2\}} & =\left(\varepsilon_{\overline{0}} \wedge \varepsilon_{1}+\varepsilon_{\overline{2}} \wedge \varepsilon_{3}\right) \cdot \varepsilon_{\{0,2\}}-t \cdot \varepsilon_{\{0,2\}} \\
& =\varepsilon_{\{1,2\}}+\varepsilon_{\{0,3\}}+\left(t_{0}+t_{2}\right) \varepsilon_{\{0,2\}} .
\end{aligned}
$$

The Schubert classes are identified by

$$
\sigma_{I}=\varepsilon_{I},
$$

so writing $\sigma_{\square}=\sigma_{\{0,1\}}$ for the divisor class, the above calculation agrees with

$$
\begin{aligned}
\sigma_{\square} \cdot 1 & =\sigma_{\{0,1\}} ; \\
\sigma_{\square} \cdot \sigma_{\{0,1\}} & =\sigma_{\{0,2\}}+\left(t_{0}+t_{1}\right) \sigma_{\{0,1\}} ; \\
\sigma_{\square} \cdot \sigma_{\{0,2\}} & =\sigma_{\{1,2\}}+\sigma_{\{0,3\}}+\left(t_{0}+t_{2}\right) \sigma_{\{0,2\}},
\end{aligned}
$$

and once again the action of $\xi_{t}$ corresponds to multiplication by $\sigma_{\square}$.

Now we are ready for the equivariant Giambelli formula for orthogonal Grassmannians. This says

$$
\begin{equation*}
\sigma_{I}=P_{\lambda}(x \mid t) \tag{12}
\end{equation*}
$$

where $\lambda$ is the strict partition corresponding to $I$ (i.e., write $I$ in decreasing order), and $P_{\lambda}(x \mid t)$ is the factorial Schur $P$-function. These polynomials were studied by Ivanov [25] and were shown to represent Schubert classes in $H_{T}^{*} O G(n)$ by Ikeda [23]; setting $t=0$, they specialize to Schur's $P$-functions (see [36, III.8]), which were shown to represent Schubert classes by Pragacz [45].

Ivanov gives a formula for $P_{\lambda}(x \mid t)$ as a ratio of Pfaffians, similar to the one defining $s_{\lambda}(x \mid t)$, and inspired by Nimmo's formula for $P_{\lambda}(x)=P_{\lambda}(x \mid 0)$ [42]. Let $A(x)=\left(a_{i j}(x)\right)$ be the $n \times n$ skew-symmetric matrix with

$$
a_{i j}(x)=\frac{x_{i}-x_{j}}{x_{i}+x_{j}},
$$

for $0 \leq i, j \leq n-1$. Given $I=\left\{i_{1}<\cdots<i_{r}\right\}$, let $B_{I}(x \mid t)=\left(b_{k l}(x \mid t)\right)$ be the $n \times$ $r$ matrix with

$$
b_{k l}(x \mid t)=\left(x_{k} \mid t\right)^{i_{l}}
$$

Form the skew-symmetric matrix ${ }^{3}$

$$
A_{I}(x \mid t)=\left[\begin{array}{c|c}
A(x) & B_{I}(x \mid t) \\
\hline-B_{I}(x \mid t)^{t} & 0
\end{array}\right] .
$$

Then the Ivanov-Nimmo formula is

$$
\begin{equation*}
P_{\lambda}(x \mid t)=\frac{\operatorname{Pf}\left(A_{I}(x \mid t)\right)}{\operatorname{Pf}(A(x))} \tag{13}
\end{equation*}
$$

By Schur's identity, the denominator is

$$
\operatorname{Pf}(A(x))=\prod_{i<j} \frac{x_{i}-x_{j}}{x_{i}+x_{j}} .
$$

Writing $\operatorname{Pf}_{K}(A)$ for the Pfaffian of the submatrix on any subset of rows and columns $K \subseteq[n]$, note that

$$
\operatorname{Pf}_{K}(A(x))=\prod_{k<k^{\prime} \text { in } K} \frac{x_{k}-x_{k^{\prime}}}{x_{k}+x_{k^{\prime}}}
$$

[^13]Now we can state the theorem. The setting is analogous to that for ordinary Grassmannians. For each even subset $K \subseteq[n]$, there is a fixed point $p_{K} \in O G(n)$ and corresponding restriction homomorphism $H_{T}^{*} O G(n) \rightarrow \mathbb{C}[t]$. On the other hand, there is a presentation of the cohomology in terms of symmetric functions, i.e., a surjective homomorphism $\mathbb{C}[t]\left[x_{0}, \ldots, x_{n-1}\right]^{S_{n}} \rightarrow H_{T}^{*} O G(n)$, and composition with the fixed-point restriction to $p_{K}$ is the evaluation

$$
\mathbb{C}[t]\left[x_{0}, \ldots, x_{n-1}\right]^{S_{n}} \rightarrow \mathbb{C}[t]
$$

given by $x_{i} \mapsto-t_{i}$ if $i \in K$, and $x_{i} \mapsto 0$ if $i \notin K$.
Theorem 12 Under the homomorphism $\mathbb{C}[t]\left[x_{0}, \ldots, x_{n-1}\right]^{S_{n}} \rightarrow H_{T}^{*} O G(n)$, we have $P_{\lambda}(x \mid t) \mapsto \sigma_{I}$, where I is the subset corresponding to the strict partition $\lambda$. Equivalently, for each $K=\left\{k_{1}<\cdots<k_{r}\right\}$, we have

$$
\left.\sigma_{I}\right|_{p_{K}}=\left.P_{\lambda}(x \mid t)\right|_{x=-t_{K}},
$$

where the specialization $x=-t_{K}$ means $x_{i} \mapsto-t_{i}$ if $i \in K$, and $x_{i} \mapsto 0$ if $i \notin K$.
Proof We follow the same outline as in type A. The first step has been done in the previous section, where we worked out the cohomology of quadrics and the corresponding change of basis between $\varepsilon_{i}$ and $f_{i}$. The second step has been done above: we identify the spinor $\varepsilon_{I} \in \mathbb{S}^{+}[t]=V_{\varpi_{n}}\left(-\varpi_{n}\right)[t]$ with the Schubert class $\sigma_{I} \in H_{T}^{*} O G(n)$.

It remains to compute the idempotent basis, and the expansion of $\varepsilon_{I}$ in this basis. With $f_{i}$ as before, we define spinors $f_{I}$ by the same formula (11) defining $\varepsilon_{I}$ :

$$
\begin{equation*}
f_{I}=f_{\bar{t}_{1}^{\prime}} \cdots f_{t_{n-r}^{\prime}} \cdot f, \tag{14}
\end{equation*}
$$

where $f=f_{0} \cdots f_{n-1}$. In fact, $f=\varepsilon$, since the change of basis is unitriangular.
Since $\varepsilon_{\emptyset}=\sigma_{\emptyset}=\mathbf{1} \in H_{T}^{*} O G(n)$, we compute this case first (to normalize the $\bar{f}_{I}$ basis). Using notation of the previous section,

$$
\varepsilon_{\bar{\imath}}=\bar{C} \cdot f_{\bar{l}}+C \cdot f_{i} .
$$

Since $\bar{C}$ is unitriangular, if we multiply by its inverse and instead consider

$$
\tilde{\varepsilon}_{\bar{\imath}}=f_{\bar{\imath}}+\bar{C}^{-1} C \cdot f_{i},
$$

we have

$$
\varepsilon_{\emptyset}=\varepsilon_{\overline{0}} \cdots \varepsilon_{\overline{n-1}} \cdot \varepsilon=\widetilde{\varepsilon}_{\overline{0}} \cdots \widetilde{\varepsilon}_{\overline{n-1}} \cdot \varepsilon
$$

in $\mathbb{S}^{+}$. On the other hand, now we are in the situation of the Theorem of the Appendix, which says

$$
\varepsilon_{\emptyset}=\sum_{K} \operatorname{Pf}_{K}\left(\bar{C}^{-1} C\right) f_{K} .
$$

To compute these Pfaffians, it helps to introduce the $n \times n$ diagonal matrix $S$, with diagonal entries

$$
\left(t_{0}-t_{n-1}\right) \cdots\left(t_{n-2}-t_{n-1}\right),\left(t_{0}-t_{n-2}\right) \cdots\left(t_{n-3}-t_{n-2}\right), \ldots, t_{0}-t_{1}, 1
$$

Then $S \bar{C}^{-1} C S=A=\left(a_{j i}\right)$, where $a_{j i}=\frac{t_{j}-t_{i}}{t_{j}+t_{i}}$. So for any $K \subseteq[n]$,

$$
\operatorname{Pf}_{K}(A)=\operatorname{Pf}_{K}\left(S \bar{C}^{-1} C S\right)=\operatorname{det}_{K}(S) \cdot \operatorname{Pf}_{K}\left(\bar{C}^{-1} C\right)
$$

where $\operatorname{det}_{K}(S)$ means the determinant of the submatrix on rows and columns $K$. Combining the formulas

$$
\operatorname{det}_{K}(S)=\prod_{\substack{i<k \\ k \in K}}\left(t_{i}-t_{k}\right) \quad \text { and } \operatorname{Pf}_{K}(A)=\prod_{\substack{i<j \\ i, j \in K}} \frac{t_{j}-t_{i}}{t_{j}+t_{i}}
$$

we obtain

$$
\operatorname{Pf}_{K}\left(\bar{C}^{-1} C\right)=\frac{1}{\prod_{\substack{i<k \\ i, k \in K}}\left(-t_{i}-t_{k}\right) \prod_{\substack{i<k \\ i \notin K, k \in K}}\left(t_{i}-t_{k}\right)}
$$

Therefore, the idempotent classes $\bar{f}_{K}$ are determined by writing

$$
\varepsilon_{\emptyset}=\sum_{K} \bar{f}_{K},
$$

that is, by setting

$$
\bar{f}_{K}=\operatorname{Pf}_{K}\left(\bar{C}^{-1} C\right) f_{K}
$$

Equivalently,

$$
f_{K}=\left(\prod_{\substack{i<k \\ i, k \in K}}\left(-t_{i}-t_{k}\right) \prod_{\substack{i<k \\ i \notin K, k \in K}}\left(t_{i}-t_{k}\right)\right) \bar{f}_{K} .
$$

Now we compute $\varepsilon_{I}$. Using notation from our computations at the end of Sect. 5, the transition between $\varepsilon_{i}$ and $f_{i}$ has the form

$$
\left[\begin{array}{c|c}
C & \bar{B} \\
\hline \bar{C} & 0
\end{array}\right] .
$$

Since the Theorem of the Appendix computes $\varepsilon_{I}$ as $(-1)^{|I|} \varepsilon_{i_{1}} \cdots \varepsilon_{i_{r}} \cdot \varepsilon_{\emptyset}$, we may replace $\varepsilon_{\bar{l}}$ by $\widetilde{\varepsilon}_{\bar{l}}$, which amounts to using the matrix

$$
\left[\begin{array}{c|c}
\bar{C}^{-1} C & \bar{B} \\
\hline w_{\circ} & 0
\end{array}\right],
$$

where $w_{\circ}$ is the $n \times n$ matrix with 1 's on the antidiagonal and 0 's elsewhere. Now formula (A.4) says

$$
\begin{equation*}
\varepsilon_{I}=\sum_{K} \operatorname{Pf}_{K}\left(\bar{A}_{I}\right) f_{K}, \tag{15}
\end{equation*}
$$

where

$$
\bar{A}_{I}=\left[\begin{array}{c|c}
\bar{C}^{-1} C & \bar{B}_{I} \\
\hline-\bar{B}_{I}^{t} & 0
\end{array}\right]
$$

and $\bar{B}_{I}$ is the submatrix of $\bar{B}$ on columns $I$.
(Note that unitriangularity of the matrix relating $\varepsilon_{i}$ and $f_{i}$ implies that

$$
\begin{aligned}
\varepsilon_{I} & =f_{I}+\sum_{K>I} \operatorname{Pf}_{K}\left(\bar{A}_{I}\right) f_{K} \\
& =\left(\prod_{\substack{j<i \\
i, j \in I}}\left(-t_{j}-t_{i}\right) \prod_{\substack{j<i \\
j \notin I, i \in I}}\left(t_{j}-t_{i}\right)\right) \bar{f}_{I}+\cdots,
\end{aligned}
$$

which agrees with the formula (10) for $\left.\sigma_{I}\right|_{p_{I}}$.)
To conclude the proof, we must relate the coefficients $\operatorname{Pf}_{K}\left(\bar{A}_{I}\right)$ to the evaluations $\left.P_{\lambda}(x \mid t)\right|_{x=-t_{K}}$. Observe first that

$$
\begin{aligned}
\left.\operatorname{Pf}(A(x))\right|_{x=-t_{K}} & =\prod_{k^{\prime}<k \text { in } K} \frac{-t_{k^{\prime}}+t_{k}}{-t_{k^{\prime}}-t_{k}} \\
& =\operatorname{Pf}_{K}(A) \\
& =\operatorname{det}_{K}(S) \operatorname{Pf}_{K}\left(\bar{C}^{-1} C\right) .
\end{aligned}
$$

Furthermore, scaling the first $n$ rows and columns of $\bar{A}_{I}$ by $S$, one computes

$$
\left[\begin{array}{c|c}
S \bar{C}^{-1} C S & S \bar{B}_{I} \\
\hline-\bar{B}_{I}^{t} S & 0
\end{array}\right]=\left.\left[\begin{array}{c|c}
A(x) & B_{I}(x \mid t) \\
\hline-B_{I}(x \mid t)^{t} & 0
\end{array}\right]\right|_{x_{0}=-t_{0}, \ldots, x_{n-1}=-t_{n-1}},
$$

and also that

$$
\left.\left.\begin{array}{rl}
\operatorname{Pf}_{K}\left(\left.\left[\begin{array}{c|c}
A(x) & B_{I}(x \mid t) \\
\hline-B_{I}(x \mid t)^{t} & 0
\end{array}\right]\right|_{x_{j}=-t_{j} \text {, all } j}\right) \\
\quad=\left(\operatorname { P f } \left[\left.\frac{A(x)}{} \right\rvert\, B_{I}(x \mid t)\right.\right. \\
\hline-B_{I}(x \mid t)^{t} & 0
\end{array}\right]\right)\left.\right|_{x=-t_{K}} .
$$

Taking Pfaffians on rows and columns $K$, it follows that

$$
\operatorname{Pf}_{K}\left(\bar{A}_{I}\right) \cdot \operatorname{det}_{K}(S)=\left.\left(\operatorname{Pf}\left[\begin{array}{c|c}
A(x) & B_{I}(x \mid t) \\
\hline-B_{I}(x \mid t)^{t} & 0
\end{array}\right]\right)\right|_{x=-t_{K}}
$$

So the Ivanov-Nimmo formula gives

$$
\left.P_{\lambda}(x \mid t)\right|_{x=-t_{K}}=\frac{\operatorname{Pf}_{K}\left(\bar{A}_{I}\right)}{\operatorname{Pf}_{K}\left(\bar{C}^{-1} C\right)}
$$

which is precisely the coefficient of $\bar{f}_{K}=\operatorname{Pf}_{K}\left(\bar{C}^{-1} C\right) f_{K}$ in the expansion of $\varepsilon_{I}$ in (15).

## 7 Rim Hook Rules for Quantum Cohomology

Consider a homogeneous space $X=G / P$ such that $H^{2}(X) \cong \mathbb{C}$. (This includes all minuscule $G / P$.) The (small) quantum cohomology ring of $X$ is a $\mathbb{C}[q]$-algebra $Q H^{*}(X)=\mathbb{C}[q] \otimes H^{*}(X)$, where $q$ is a formal parameter whose degree depends on $X$, equipped with a product which deforms the usual cup product on $H^{*}(X)$; the structure constants in $Q H^{*}(X)$ are 3-point Gromov-Witten invariants, counting certain rational curves in $X$. Using the action of a maximal torus $T \subseteq G$ on $X$, there are also equivariant versions $Q H_{T}^{*}(X)$. We will write $T=\left(\mathbb{C}^{*}\right)^{n}$, and

$$
S=S_{n}=H_{T}^{*}(\mathrm{pt}) \cong \mathbb{C}\left[t_{1}, \ldots, t_{n}\right]
$$

so $Q H_{T}^{*}(X)$ is an algebra over $S[q]$. (We continue to use $\mathbb{C}$ coefficients for cohomology, but all results of this section are equally valid with $\mathbb{Z}$ coefficients.)

For any $N$, consider the embedding of Grassmannians

$$
\iota: G r(k, N) \hookrightarrow G r(k, N+1)
$$

corresponding to the standard embedding of $\mathbb{C}^{N}$ in $\mathbb{C}^{N+1}$ (as the span of the first $N$ standard basis vectors). Let us write $T_{N}=\left(\mathbb{C}^{*}\right)^{N}$, and $S_{N}=H_{T_{N}}^{*}(\mathrm{pt})$. There is a similar embedding of $T_{N}$ in $T_{N+1}$ (corresponding to the map on character groups $\mathbb{Z}^{N+1} \rightarrow \mathbb{Z}^{N}$ sending $t_{N+1}$ to 0 and all other $t_{i}$ to $t_{i}$ ). The embedding $\iota$ is equivariant with respect to the inclusion of tori and their natural actions


Fig. 1 A 7-rim hook of height 2
on the Grassmannians, so we have an inverse system of graded homomorphisms $\iota^{*}: H_{T_{N+1}}^{*} \operatorname{Gr}(k, N+1) \rightarrow H_{T_{N}}^{*} \operatorname{Gr}(k, N)$. Let $H_{T_{\infty}}^{*} G r(k, \infty)$ be the graded inverse limit. This ring can be identified with the ring of symmetric polynomials in variables $x_{1}, \ldots, x_{k}$, with coefficients in $S_{\infty}$.

Writing $\Omega_{\lambda}^{(N)}$ for the Schubert variety in $\operatorname{Gr}(k, N)$ and $\sigma_{\lambda}^{(N)}$ for its equivariant cohomology class, one checks that $\iota^{-1} \Omega_{\lambda}^{(N+1)}=\Omega_{\lambda}^{(N)}$, and it follows that $\iota^{*} \sigma_{\lambda}^{(N+1)}=$ $\sigma_{\lambda}^{(N)}$. Therefore we have well-defined classes $\sigma_{\lambda} \in H_{T_{\infty}}^{*} \operatorname{Gr}(k, \infty)$.

Fixing $n$ and $T=T_{n}$, for $N>n$, consider a (different) inclusion of tori $T \hookrightarrow T_{N}$ given by

$$
\left(z_{1}, z_{2}, \ldots, z_{n}\right) \mapsto\left(z_{1}, z_{2}, \ldots, z_{n}, z_{1}, z_{2}, \ldots\right)
$$

equivalently, take the map on character groups $\mathbb{Z}^{N} \rightarrow \mathbb{Z}^{n}$ given by $t_{i} \mapsto t_{i}(\bmod n)$, where the representatives $\bmod n$ are taken to be $1, \ldots, n$. This also defines a ring homomorphism $S_{\infty} \rightarrow S_{n}$, inducing an algebra homomorphism $H_{T_{\infty}}^{*} G r(k, \infty) \rightarrow$ $H_{T}^{*} \operatorname{Gr}(k, \infty)$. We will see how to extend this to a homomorphism $H_{T}^{*} G r(k, \infty) \rightarrow$ $Q H_{T}^{*} G r(k, n)$.

The rim hook rule ${ }^{4}$ for the equivariant quantum cohomology of $\operatorname{Gr}(k, n)$ is a recipe for computing quantum products in terms of ordinary products on $\operatorname{Gr}(k, N)$, for $N>n$. Given any partition $\lambda$, an $n$-rim hook (also called an $n$-border strip) is a connected collection of $n$ boxes on the southeast border of its diagram, with exactly one box in each diagonal, such that the complement $\mu$ is also the diagram of a partition. If the boxes appear in rows $i$ through $j$ of the diagram, the height of the rim hook is $j-i$. See Fig. 1 for an example, and [36, Sect. I.1] for more details.

Any partition $\lambda$ has a well-defined $n$-core $\mu$, which is the partition obtained by removing all possible $n$-rim hooks from $\lambda$, in any order. The parity of the sum of the heights of these rim hooks is also independent of choices, so we can define $\epsilon(\lambda / \mu)=\sum_{\delta}$ height $(\delta)(\bmod 2)$, taking the sum over a pieces of a decomposition of the skew shape $\lambda / \mu$ into $n$-rim hooks $\delta$.

Under the bijection between partitions $\lambda$ with $k$ parts and $k$-element subsets $I$ of $\{1,2, \ldots\}$, the operation of removing an $n$-rim hook corresponds to replacing some element $i \in I$ with $i-n$ (and an $n$-rim hook exists only if there is an $i>n$ such that $i-n$ is not in $I$ ); the height of such a rim hook is the length of the permutation needed to sort the resulting set into increasing order. Taking the $n$-core of $\lambda$ can be described as follows. Write the elements of $I$ as $s_{i} n+r_{i}$, for $1 \leq i \leq k$ and $1 \leq r_{i} \leq n$. Next consider the multiset of residues $\left\{r_{i}\right\}$. For each $r$ that appears in

[^14]this multiset, replace the $j$ th occurrence of $r$ by $r+(j-1) n$, obtaining a set $\bar{I}$ of $k$ distinct integers with the same multiset of residues modulo $n$. The $n$-core of $\lambda$ is the partition $\mu$ corresponding to $\bar{I}$, and the sign $(-1)^{\epsilon(\lambda / \mu)}$ is the sign of the permutation needed to sort $\bar{I}$ (see [36, Sect. I.1, Ex. 8]).

For example, consider $\lambda=(7,6,3)$. For $k=3$, this corresponds to the set $I=$ $\{4,8,10\}$. There are two ways to remove a 7 -rim hook: replace 10 with 3 (obtaining $I^{\prime}=\{4,8,3\}$ ) or replace 8 with 1 (obtaining $I^{\prime}=\{4,1,10\}$ ). In the first case, the permutation which sorts $I^{\prime}$ has length 2 , and in the second it has length 1 . To find the 7 -core, write residues modulo 7 to obtain $\{4,1,3\}$; this is sorted to $\bar{I}=\{1,3,4\}$ by a permutation of length 2 , so $\mu=(1,1)$ and $\epsilon(\lambda / \mu)$ is even.

Given a class $\sigma_{\lambda} \in H_{T_{\infty}}^{*} \operatorname{Gr}(k, \infty)$, for $\lambda=\left(\lambda_{1} \geq \cdots \geq \lambda_{k} \geq 0\right)$, let $\mu$ be the $n$ core of $\lambda$ and define a linear map

$$
\begin{aligned}
\varphi: H_{T_{\infty}}^{*} G r(k, \infty) & \rightarrow Q H_{T}^{*} G r(k, n) \\
\sigma_{\lambda} & \mapsto \begin{cases}(-1)^{(k-1) s+\epsilon(\lambda / \mu)} q^{s} \sigma_{\mu} & \text { if } \mu \subseteq \rho_{k, n-k} \\
0 & \text { otherwise }\end{cases}
\end{aligned}
$$

where $s=(|\lambda|-|\mu|) / n$ is the number of rim hooks removed from $\lambda$ to obtain $\mu$. This map factors into the specialization homomorphism $H_{T_{\infty}}^{*} \operatorname{Gr}(k, \infty) \rightarrow H_{T}^{*} G r(k, \infty)$ described above, followed by a linear map $\psi: H_{T}^{*} \operatorname{Gr}(k, \infty) \rightarrow Q H_{T}^{*} G r(k, n)$ which is given by the same formula as $\varphi$. The first is clearly an algebra homomorphism (with respect to the ring map $S_{\infty} \rightarrow S$ which cyclically specializes the $t$ variables), while the second is a priori a homomorphism of $S$-modules.

Theorem 13 (Equivariant rim hook rule $[2,3]$ ) The map $\varphi$ respects multiplication: it is a surjective homomorphism of algebras, compatible with the cyclic specialization $S_{\infty} \rightarrow S$.

It suffices to show that the second factor $\psi$ is a surjective ring homomorphism. We will give a simple proof, inspired by the Satake correspondence. (A similar construction was described by Gatto [14]. The phrasing in terms of reduction modulo $n$ also appears in work by Buch [7, Corollary 1] and Sottile [47].) The basic idea is to show that the kernel of $\psi$ is an ideal, so that $\psi$ induces a ring structure on its image; then apply Mihalcea's characterization of the quantum product [39] to conclude that the product induced by $\psi$ is the quantum product.

Before proving the theorem, we first consider the easier case of projective space, so $k=1$. Writing a class in $H_{T}^{*} \mathbb{P}^{\infty}$ as $\sigma_{s n+i}$, for $0 \leq i \leq n-1$, the rim hook rule says $\psi\left(\sigma_{s n+i}\right)=q^{s} \sigma_{i}$. (In this case, $\psi$ is an isomorphism of $S$-modules.) Both $H_{T_{n}}^{*} \mathbb{P}^{\infty}$ and $Q H_{T}^{*} \mathbb{P}^{n-1}$ are free $S$-algebras, generated by the divisor class $\sigma_{\square}=\sigma_{1}$, so to see that $\psi$ is an isomorphism of algebras it suffices to check that it respects multiplication by $\sigma_{\square}$. This is a simple application of a special case of the equivariant (quantum) Pieri rule [39]: in $H_{T}^{*} \mathbb{P}^{\infty}$ one has

$$
\sigma_{\square} \cdot \sigma_{i}=\sigma_{i+1}+\left(t_{1}-t_{(i+1)}(\bmod n)\right) \sigma_{i}
$$

and in $Q H_{T_{n}}^{*} \mathbb{P}^{n-1}$ one has

$$
\sigma_{\square} \cdot \sigma_{i}= \begin{cases}\sigma_{i+1}+\left(t_{1}-t_{i+1}\right) \sigma_{i} & \text { for } i<n-1 \\ q \sigma_{0}+\left(t_{1}-t_{n}\right) \sigma_{n-1} & \text { for } i=n-1\end{cases}
$$

This computation can be rephrased to make it analogous to the one we did in the introduction. There are homomorphisms with compatible identifications

where the first horizontal map is given by the cyclic specialization $t_{i} \mapsto t_{i(\bmod n)}$ as above. Identifying the standard basis $\varepsilon_{i}$ with $\sigma_{i-1}$ as before, the rim hook rule says $\psi$ is given by $\varepsilon_{s n+i} \mapsto q^{s} \varepsilon_{i}$, for $1 \leq i \leq n$. (Recalling that $\varepsilon_{i}$ is a weight vector with weight $t_{i}$, one sees that the cyclic specialization is necessary to make the second map compatible with the $S_{n}$-module structure.)

Multiplication by $\sigma_{\square}$ on $H_{T_{\infty}}^{*} \mathbb{P}^{\infty}$ is given by the infinite matrix $\xi_{t}$ with $-t_{1},-t_{2}, \ldots$ on the diagonal, 1's on the subdiagonal, and 0's elsewhere. (One should twist ( $\left.S_{\infty}\right)^{\infty}$ by e ${ }^{-t_{1}}$ as in Sect. 3 to get the correct action.) Specializing the diagonal variables modulo $n$ and applying the homomorphism $\psi$, the action of $\xi_{t}$ on $\left(S_{\infty}\right)^{\infty}$ is transformed into the action of the $n \times n$ matrix

$$
\xi_{q, t}=\left[\begin{array}{cccc}
-t_{1} & & & q  \tag{17}\\
1 & \ddots & & \\
0 & \ddots & \ddots & \\
0 & 0 & 1 & -t_{n}
\end{array}\right]
$$

on $(S[q])^{n}$, which agrees with multiplication by $\sigma_{\square}$ in $Q H_{T}^{*} \mathbb{P}^{n-1}$.
Lemma 14 The $S$-module homomorphism $\psi: H_{T}^{*} \operatorname{Gr}(k, \infty) \rightarrow Q H_{T}^{*} \operatorname{Gr}(k, n)$ is surjective, and its kernel is an ideal.

Proof It is straightforward to see that $\psi$ is surjective, and we leave this to the reader. It is also easy to see that the kernel is generated (as an $S$-module) by two types of elements:
(1) classes $\sigma_{\lambda}$ such that the $n$-core of $\lambda$ does not fit in $\rho_{k, n-k}$; and
(2) differences $(-1)^{\epsilon(\lambda / \mu)} \sigma_{\lambda}-(-1)^{\epsilon\left(\lambda^{\prime} / \mu\right)} \sigma_{\lambda^{\prime}}$, for two partitions $\lambda, \lambda^{\prime}$ of the same size and with the same $n$-core $\mu$.

We will show that the products of such elements with $S$-algebra generators for $H_{T}^{*} \operatorname{Gr}\left(k, \mathbb{C}^{\infty}\right)$ are also in the kernel of $\psi$.

Making identifications

up to sign the homomorphism $\psi$ is induced by the corresponding map described above for projective space: if one sends $\varepsilon_{s n+i}$ to $\left((-1)^{k-1} q\right)^{s} \varepsilon_{i}$, then $\psi$ is the induced map on the $k$ th exterior power. Furthermore, from this point of view, the operators $\xi_{t}, \xi_{t}^{2}, \ldots, \xi_{t}^{k}$-or rather, the cohomology classes they correspond to-form a set of $S$-algebra generators for $H_{T}^{*} G r(k, \infty)$. (To see this, recall that $H_{T}^{*} G r(k, \infty)$ is isomorphic to the ring of symmetric polynomials in $k$ variables with coefficients in $S$, and that the leading term of $\xi_{t}^{j}$ corresponds to multiplication by the power sum function $p_{j}$; these generate the ring of symmetric polynomials, for $1 \leq j \leq k$.)

Now let us consider the generators of the kernel of $\psi$, using the correspondence between partitions and $k$-element sets to write classes in $H_{T}^{*} G r(k, \infty)$ as $\sigma_{\lambda}=\varepsilon_{I}$. For $I=\left\{i_{1}<\cdots<i_{k}\right\}$, extract residues by writing $i_{a}=s_{a} n+r_{a}$. The $n$-core of $\lambda$ fits inside $\rho_{k, n-k}$ if and only if the residues $r_{a}$ are all distinct, and two partitions $\lambda, \lambda^{\prime}$ share the same $n$-core if and only if the corresponding sets $I, I^{\prime}$ share the same multiset of residues.

Suppose $\varepsilon_{I}$ is the first type of generator for the kernel, so $I$ has at least two elements $i, i^{\prime}$ with the same residue modulo $n$. Reordering as necessary, we can write

$$
\varepsilon_{I}= \pm\left(\varepsilon_{i} \wedge \varepsilon_{i^{\prime}} \wedge \cdots\right)
$$

Using the formula from Example 6,

$$
\begin{aligned}
\xi_{t}^{j} \cdot\left(\varepsilon_{i} \wedge \varepsilon_{i^{\prime}} \wedge \cdots\right)= & \left(\sum_{a=0}^{j}(-1)^{a} h_{a}\left(t_{i}, \ldots, t_{i+j-a}\right) \varepsilon_{i+j-a}\right) \wedge \varepsilon_{i^{\prime}} \wedge \cdots \\
& +\varepsilon_{i} \wedge\left(\sum_{a=0}^{j}(-1)^{a} h_{a}\left(t_{i^{\prime}}, \ldots, t_{i^{\prime}+j-a}\right) \varepsilon_{i^{\prime}+j-a}\right) \wedge \cdots \\
& +\left(\text { terms involving } \varepsilon_{i} \wedge \varepsilon_{i^{\prime}} \wedge \cdots\right)
\end{aligned}
$$

Due to the cyclic specialization, the subscripts on the $t$ variables should be read modulo $n$, and since $i \equiv i^{\prime}(\bmod n)$, the coefficients $h_{a}\left(t_{i}, \ldots, t_{i+j-a}\right)$ and $h_{a}\left(t_{i^{\prime}}, \ldots\right.$, $t_{i^{\prime}+j-a}$ ) are equal. After applying $\psi$, the first two terms cancel, and all the others go to zero. It follows that $\xi_{t}^{j} \cdot \varepsilon_{I}$ also lies in the kernel.

The second case is similar. Recall that the $\operatorname{sign}(-1)^{\epsilon}(\lambda / \mu)$ is equal to the sign of the permutation needed to sort the residues of $I$ into increasing order. Suppose $I$ and $I^{\prime}$ have the same (distinct) residues modulo $n$, and are such that the corresponding partitions $\lambda$ and $\lambda^{\prime}$ have the same size. Then

$$
(-1)^{\epsilon(\lambda / \mu)} \varepsilon_{I}-(-1)^{\epsilon\left(\lambda^{\prime} / \mu\right)} \varepsilon_{I^{\prime}}=\left(\varepsilon_{s_{1} n+r_{1}} \wedge \cdots \wedge \varepsilon_{s_{k} n+r_{k}}\right)-\left(\varepsilon_{s_{1}^{\prime} n+r_{1}} \wedge \cdots \wedge \varepsilon_{s_{k}^{\prime} n+r_{k}}\right)
$$

for $r_{1}<\cdots<r_{k}$. A calculation analogous to the previous one shows that the product of this difference with $\xi_{t}^{j}$ remains in the kernel of $\psi$.

Proof of Theorem 13 As remarked above, it suffices to show that $\psi$ is a ring homomorphism. By the lemma, $\operatorname{ker}(\psi)$ is an ideal of $H_{T}^{*} \operatorname{Gr}(k, \infty)$. Let $A=$ $H_{T}^{*} \operatorname{Gr}(k, \infty) / \operatorname{ker}(\psi)$, and write $q \in A$ for the image of $(-1)^{k-1} \sigma_{n}$. This makes $A$ an $S[q]$-algebra, with an $S[q]$-module basis indexed by partitions fitting inside $\rho_{k, n-k}$. The homomorphism $\psi$ maps this basis onto the Schubert basis for $Q H_{T}^{*} G r(k, n)$. We must check that this is an isomorphism of algebras.

By Mihalcea's characterization of $Q H_{T}^{*} G r(k, n)$ [39, Corollary 7.1], it suffices to check that the isomorphism respects multiplication by the divisor class $\sigma_{1}$. On one hand, multiplication by $\sigma_{1}$ on $A$ is given by the action of a matrix $\xi_{q, t}$ on $\bigwedge_{S[q]}^{k}(S[q])^{n} \otimes \mathrm{e}^{-\varpi_{k}}$. (The matrix $\xi_{q, t}$ is the one of (17), with $q$ replaced by $(-1)^{k-1} q$.) That is,

$$
\begin{aligned}
\xi_{q, t} \cdot \varepsilon_{I}= & \varepsilon_{i_{1}+1} \wedge \varepsilon_{i_{2}} \wedge \cdots \wedge \varepsilon_{i_{k}}+\cdots+\varepsilon_{i_{1}} \wedge \varepsilon_{i_{2}} \wedge \cdots \wedge \varepsilon_{i_{k}+1} \\
& +\left(t_{1}+\cdots+t_{k}-t_{i_{1}}-\cdots-t_{i_{k}}\right) \varepsilon_{I} \\
& \left(+q \varepsilon_{1} \wedge \varepsilon_{i_{1}} \wedge \cdots \wedge \varepsilon_{i_{k-1}}\right),
\end{aligned}
$$

with the last term occurring only if $i_{k}=n$ and $i_{1}>1$. Writing this in terms of partitions, one recovers exactly the equivariant quantum Pieri formula [39, Theorem 1] for multiplication by $\sigma_{1}$ in $Q H_{T}^{*} G r(k, n)$, as desired. ${ }^{5}$

## 8 Closing Remarks and Other Directions

### 8.1 Relation to Quantum Integrability

In [20, 21], Gorbounov, Korff, and Stroppel indicate an alternative approach to the (equivariant) rim-hook rule, based on quantum integrability of the six-vertex model. This model can be seen as a degeneration of the Yangian Hopf algebra (a certain quantum group) which governs the cohomology of the cotangent space of the Grassmannian. This fits in a major program initiated by Maulik and Okounkov to relate representation theory of quantum groups to enumerative geometry [38, 44, 46].

The authors of $[20,21]$ consider a vector space $V=\mathbb{C} v_{0} \oplus \mathbb{C} v_{1}$. The standard basis $\left\{v_{\omega_{1}} \otimes v_{\omega_{1}} \otimes \ldots \otimes v_{\omega_{n}}\right\}$ of $V^{\otimes n}$, where $\omega_{i} \in\{0,1\}$, can be identified with set of 01 -words $\left\{\omega=\omega_{1} \omega_{2} \ldots \omega_{n}\right\}$. The set of words $\left\{\omega: \sum_{i=1}^{n} \omega_{i}=k\right\}$ is in bijection with the set of multi-indexes $I=\left\{i_{1}<\cdots<i_{k}\right\} \subseteq[n]$. The bijection is given by sending $I$ to the 01 -word with 1 's in positions $i_{1}, \ldots, i_{k}$ and 0 's elsewhere. Through

[^15]this identification, let $V_{k} \subseteq V^{\otimes n}$ be the vector subspace generated by vectors indexed by this set. In [21], two bases for $V_{k} \otimes \mathbb{C}\left[t_{1}, \ldots, t_{n}\right]$ are constructed, which are called respectively standard or spin basis ${ }^{6}$ and the Bethe vector basis. Then a connection is made with equivariant cohomology of the Grassmannian, by identifying those bases with the Schubert basis and the torus fixed point basis, respectively. The authors construct certain (generating series of) operators denoted $A+q D$ ([21, Sect. 5]) acting on the spin basis which match the action of quantum multiplication by the divisor class on the Schubert basis.

This is analogous to the setup from the previous section, where the principal nilpotent $\xi_{t}$ is corrected by the nilpotent matrix

$$
q\left[\begin{array}{llll}
0 & & & 1 \\
0 & \ddots & & \\
0 & \ddots & \ddots & \\
0 & 0 & 0 & 0
\end{array}\right]
$$

in order to match the quantum equivariant Chevalley rule after a twist.
We should also mention another closely related result by Lam and Templier, [32, Proposition 4.13] where essentially the same presentation of the quantum Chevalley operator is given.

### 8.2 Mirror Symmetry and the Gamma Conjectures

Mirror symmetry for Fano varieties roughly claims a correspondence which associate to each Fano variety $X$ a smooth quasi-projective variety $Y$ endowed with a flat map $w: Y \longrightarrow \mathbb{A}^{1}$ with quasi-projective fibres, called superpotential. Under this correspondence critical points of $w$ give rise to vanishing cycles which through homological mirror symmetry correspond to objects of the derived category $D^{b}(X)$ of coherent sheaves on $X$.

The Gamma conjectures were formulated by Galkin, Golyshev and Iritani in [13], refining previous conjectures by Dubrovin [11]. They claim that simple eigenvalues of a certain variation of the quantum Chevalley operator denoted $\star_{0}$ are related to characteristic classes $\widehat{\Gamma}_{X} C h(E)$, where $E$ is an exceptional object of $D^{b}(X)$, which in turn are constructed through flat sections of Dubrovin's quantum connection

$$
\begin{equation*}
\nabla_{z \partial_{z}}=z \frac{\partial}{\partial z}-\frac{1}{z}\left(c_{1}(X) \star_{0}\right)+\mu \tag{19}
\end{equation*}
$$

[^16]a meromorphic flat connection on the trivial bundle $H^{*}(X) \times \mathbb{P}^{1} \rightarrow \mathbb{P}^{1}$. In (19) $\mu$ is the grading operator on $H^{*}(X)$ and $z$ is a local coordinate on $\mathbb{P}^{1}$.

The Gamma conjectures were proved for Grassmannians in [13], (see also [9]), by making fundamental use of the quantum correction to the Satake correspondence.

### 8.3 Other Representations

We have focused on minuscule representations in types A and D, with their corresponding (minuscule) homogeneous spaces $G / P$ : Grassmannians, quadrics, and orthogonal Grassmannians. It is natural to ask what can be said for other spaces.

In other types ( $E_{6}$ and $E_{7}$ ), there is only one minuscule space, up to isomorphism, so the approach we used in Sects. 4 and 6 does not seem productive. Isomorphisms between quantum cohomology of these spaces and corresponding representations were worked out by Golyshev and Manivel [19]. On the other hand, it would be interesting to see an analogue of the rim-hook rule for orthogonal Grassmannians, presumably related to an infinite-dimensional spin representation.

In a different direction, one can ask for the transition matrices between the " $\varepsilon$ " and " $\bar{f}$ " bases of non-minuscule representations. This should correspond to localization formulas for MV-cycles. We have seen Schur $S$ - and $P$-functions appear naturally in the exterior and spin representations; what other symmetric functions arise this way?

Acknowledgements This note is partly based on a talk given by the first author at a conference dedicated to the memory of Dan Laksov (Mittag-Leffler, June 2014). The ideas grew out of conversations we had with Roi Docampo at IMPA. We learned about the connection between quantum Schubert calculus and the Satake correspondence from a remarkable preprint of Golyshev and Manivel [19], and the debt we owe to their work should be evident. We also thank Reimundo Heluani and Joel Kamnitzer for helping us understand the geometric Satake correspondence. Finally, we thank the referee for a very careful reading and thoughtful comments.

## Appendix: Pfaffians and Spinors

In this appendix, we prove a change-of-basis formula in the spin representation, where Pfaffians play the role analogous to determinants in the exterior algebra. This refines similar formulas of Chevalley and Manivel [10, 37].

First, for any complex vector space $V$ with symmetric bilinear form $\langle$,$\rangle , the$ associated Clifford algebra is the quotient of the tensor algebra by two-sided ideal forcing the relation $v \cdot w+v \cdot w=\langle v, w\rangle 1$ for all vectors $v, w \in V$ :

$$
C l(V):=T^{\bullet}(V) /(v \otimes w+w \otimes v-\langle v, w\rangle 1) .
$$

(Note that this definition differs slightly from the standard one, where $\langle v, w\rangle 1$ would be replaced by $2\langle v, w\rangle 1$-but it is equivalent, up to rescaling the form $\langle$,$\rangle , since we$ are not in characteristic 2.) If the bilinear form is zero, this is just the exterior algebra: $C l(V) \cong \bigwedge^{\bullet} V$. In general, $\operatorname{dim} C l(V)=2^{\operatorname{dim} V}$, with a basis consisting of products $v_{I}=v_{i_{1}} \cdots v_{i_{k}}$ of distinct basis elements of $V$. (One can prove this by degenerating to the exterior algebra.)

We note the following general formulas in $C l(V)$, which are immediate from the defining relations. First, if $x$ and $y$ are orthogonal vectors in $V$ (i.e., $\langle x, y\rangle=0$ ), then

$$
\begin{equation*}
x \cdot y=-y \cdot x \tag{A.1}
\end{equation*}
$$

in $C l(V)$. Next, suppose $x=x_{1} \cdots x_{r} \in C l(V)$ is a product of vectors in $V$ such that $X=\operatorname{span}\left\{x_{1}, \ldots, x_{r}\right\} \subseteq V$ is isotropic. Then for any vector $y \in X$,

$$
\begin{equation*}
y \cdot x=0 \tag{A.2}
\end{equation*}
$$

in $\mathrm{Cl}(\mathrm{V})$.
From now on, we assume $\operatorname{dim} V=2 n$ and its bilinear form is nondegenerate. Let $y_{\overline{n-1}}, \ldots, y_{\overline{0}}, y_{0}, \ldots, y_{n-1}$ be an "orthogonal basis", meaning $\left\langle y_{\bar{i}}, y_{j}\right\rangle=\delta_{i j}$. (We interpret the bar as a negative sign, so $\overline{\bar{l}}=i$.) Let $y=y_{0} \cdots y_{n-1}$. The spin representation is the left ideal

$$
\mathbb{S}=C l(V) \cdot y .
$$

A (pure) spinor is an element of the form $z \cdot y \in \mathbb{S}$, where $z=z_{1} \cdots z_{n}$ is a product of vectors $z_{1}, \ldots, z_{n} \in V$ which span a maximal isotropic subspace of $V$.

Let $[m]=\{0, \ldots, m-1\}$ for any integer $m \geq 1$. For a subset $I=\left\{i_{1}<\cdots<\right.$ $\left.i_{r}\right\} \subseteq[n]$, let $I^{\prime}=\left\{i_{1}^{\prime}<\cdots<i_{n-r}^{\prime}\right\}=[n] \backslash I$. Given such an $I$, the corresponding standard spinor is

$$
y_{I}:=y_{i_{1}^{\prime}} \cdots y_{t_{n-r}^{\prime}} \cdot y .
$$

These form a basis for $\mathbb{S}$, as $I$ ranges over all subsets of $[n]$. Note that $y_{\varnothing}=$ $y_{\overline{0}} \cdots y_{\overline{n-1}} \cdot y$, and

$$
\begin{equation*}
y_{i_{1}} \cdots y_{i_{r}} \cdot y_{\emptyset}=(-1)^{i_{1}+\cdots+i_{r}} y_{I} \tag{A.3}
\end{equation*}
$$

which is sometimes useful.
The spin representation becomes an $\mathfrak{5 o}_{2 n}$-module via the embedding $\mathfrak{s o}_{2 n} \cong$ $\bigwedge^{2} V \hookrightarrow C l(V)$ by

$$
v \wedge w \mapsto \frac{1}{2}(v \cdot w-w \cdot v)=v \cdot w-\frac{1}{2}\langle v, w\rangle 1
$$

So $\mathfrak{S o}_{2 n}$ preserves the parity of basis vectors of $C l(V)$, and the spin representation breaks into two irreducible half-spin representations

$$
\mathbb{S}=\mathbb{S}^{+} \oplus \mathbb{S}^{-}
$$

Our convention will be that $\mathbb{S}^{+}$is spanned by standard spinors $y_{I}$, with $I$ an even subset of [ $n$ ], i.e., it has even cardinality.

Now let $x_{\overline{n-1}}, \ldots, x_{\overline{0}}, x_{0}, \ldots, x_{n-1}$ be another orthogonal basis, and assume it is related to the $y_{i}$ 's by a unitriangular matrix:

$$
x_{\bar{\imath}}=y_{\bar{\imath}}+\sum_{j<i} \bar{c}_{j i} y_{\bar{J}}+\sum_{j=0}^{n-1} c_{j i} y_{j}
$$

and

$$
x_{i}=y_{i}+\sum_{j>i} b_{j i} y_{j} .
$$

Thus $x=x_{0} \cdots x_{n-1}$ is equal to $y=y_{0} \cdots y_{n-1}$.
We define pure spinors $x_{I}$ analogously, by

$$
x_{I}:=x_{\bar{l}_{1}^{\prime}} \cdots x_{\bar{l}_{n-r}^{\prime}} \cdot y .
$$

Our goal is to compute the expansion of $x_{I}$ in the basis $y_{K}$.
In fact, we will be free to multiply by the inverse of the matrix $\bar{C}=\left(\bar{c}_{j i}\right)$ and assume that $x_{\bar{l}}$ is written

$$
x_{\bar{\imath}}=y_{\bar{\imath}}+\sum_{j=0}^{n-1} a_{j i} y_{j} .
$$

In this case, the fact that the $x_{\bar{I}}$ span an isotropic space is equivalent to the fact that the matrix $A=\left(a_{j i}\right)$ is skew-symmetric. Up to appropriately indexing rows and columns, the $x$ basis is related to the $y$ basis by a matrix of the form

$$
\left(\begin{array}{c|c}
A & B \\
\hline w_{\circ} & 0
\end{array}\right),
$$

where $w_{\circ}$ is the matrix with 1 's on the antidiagonal and zeroes elsewhere. The top $n$ rows determine a skew-symmetric matrix (by replacing $w_{\circ}$ with $-B^{t}$ ). For a subset $I \subseteq[n]$, let $B(I)$ be the submatrix formed by taking columns $I$ of $B$, and let $A(I)$ be the skew-symmetric matrix

$$
A(I)=\left(\begin{array}{c|c}
A & B(I) \\
\hline-B(I)^{t} & 0
\end{array}\right) .
$$

Only the top $n$ rows are needed to perform operations with such a matrix. For instance, if $n=3$ and $I=\{1\}$, the top rows are

$$
\left(\begin{array}{ccc|ccc}
0 & a_{21} & a_{20} & b_{20} & b_{21} & \\
& 0 & a_{10} & b_{10} & 1 & 0 \\
& & 0 & 1 & 0 & 0
\end{array}\right)
$$

and

$$
A(I)=\left(\begin{array}{ccc|c}
0 & a_{21} & a_{20} & b_{21} \\
-a_{21} & 0 & a_{10} & 1 \\
-a_{20} & -a_{10} & 0 & 0 \\
\hline-b_{21} & -1 & 0 & 0
\end{array}\right)
$$

For even $r$, the Pfaffian of any $r \times r$ skew-symmetric matrix $A$ may be computed recursively using the Laplace-type expansion formula

$$
\operatorname{Pf}(A)=\sum_{j=1}^{r-1}(-1)^{j-1} a_{j r} \operatorname{Pf}_{\widehat{j, r}}(A)
$$

where $\operatorname{Pf}_{\widehat{j}, r}(A)$ is the Pfaffian of the submatrix of $A$ obtained by removing the $j$ th and $r$ th rows and columns. Let $\mathrm{Pf}_{K}(A)$ denote the Pfaffian of the submatrix on rows and columns $K$, and we always use the convention $\operatorname{Pf}_{\emptyset}(A)=1$.

Theorem For a subset $I \subseteq[n]$, we have

$$
\begin{equation*}
x_{I}=\sum_{\substack{K \subseteq[n] \\ K \text { even }}} \operatorname{Pf}_{K}(A(I)) y_{K}, \tag{A.4}
\end{equation*}
$$

where the sum is over subsets $K$ of even cardinality.
See also [43, Corollary 2.4] for a related Pfaffian identity.
Proof First we establish the formula for $x_{\emptyset}$. In fact, for any $m \leq n$, we have

$$
x_{\overline{0}} \cdots x_{\overline{m-1}} \cdot y=\sum_{\substack{K \subseteq[m] \\ K \text { even }}} \operatorname{Pf}_{K}(A) y_{\bar{k}_{1}^{\prime}} \cdots y_{\vec{k}_{s}^{\prime}} \cdot y
$$

where $K^{\prime}=\left\{k_{1}^{\prime}<\cdots<k_{s}^{\prime}\right\}=[m] \backslash K$. This is proved by induction on $m$, the base case $m=0$ being the tautology $y=y$. For the induction step, we compute using the Laplace expansion formula:

$$
\begin{aligned}
x_{\overline{0}} \cdots x_{\bar{m}} \cdot y & =(-1)^{m} x_{\bar{m}} \cdot x_{\overline{0}} \cdots x_{\overline{m-1}} \cdot y \\
& =(-1)^{m}\left(y_{\bar{m}}+\sum_{j=0}^{m-1} a_{j m} y_{j}\right) \sum_{K \subseteq[m]} \operatorname{Pf}_{K}(A) y_{\bar{k}_{1}^{\prime}} \cdots y_{\bar{k}_{r}^{\prime}} \cdot y \\
& =\sum_{K \subseteq[m]} \operatorname{Pf}_{K}(A) y_{y_{k_{1}^{\prime}}} \cdots y_{\bar{k}_{r}^{\prime}} \cdot y_{\bar{m}} \cdot y
\end{aligned}
$$

$$
+\sum_{j=0}^{m-1} \sum_{K \subseteq[m]}(-1)^{m} a_{j m} \operatorname{Pf}_{K}(A) y_{j} \cdot y_{\vec{k}_{1}} \cdots y_{\vec{k}_{r}^{\prime}} \cdot y .
$$

Note

$$
y_{j} \cdot y_{\vec{k}_{1}^{\prime}} \cdots y_{\vec{k}_{r}^{\prime}} \cdot y= \begin{cases}0 & \text { if } j \in K \\ (-1)^{a-1} y_{\widehat{k_{1}^{\prime}}} \cdots \widehat{y_{k_{a}^{\prime}}} \cdots y_{\vec{k}_{r}^{\prime}} \cdot y & \text { if } j=k_{a}^{\prime} \in K^{\prime} .\end{cases}
$$

So the second sum above becomes

$$
\sum_{K} \sum_{j \neq K}(-1)^{m+a-1} a_{j m} \operatorname{Pf}_{K}(A) y_{\bar{k}_{1}^{\prime \prime}} \cdots y_{\bar{k}_{r-1}^{\prime \prime}} \cdot y,
$$

where $K^{\prime \prime}=K^{\prime} \backslash j$. Combining the two and using Laplace expansion yields the claimed formula for $x_{\emptyset}$.

The general case is similar. Using (A.3), it is equivalent to show

$$
x_{i_{1}} \cdots x_{i_{r}} \cdot x_{\emptyset}=(-1)^{i_{1}+\cdots+i_{r}} \sum_{\substack{K \subseteq[n] \\ K \text { even }}} \operatorname{Pf}_{K}(A(I)) y_{K},
$$

and we do this by induction on $r$. Let $I=\left\{i_{1}<\cdots<i_{r}\right\}$ and consider $i>i_{r}$. We compute:

$$
\begin{aligned}
x_{i_{1}} \cdots x_{i_{r}} \cdot x_{i} \cdot x_{\emptyset} & =(-1)^{r} x_{i} \cdot x_{i_{1}} \cdots x_{i_{r}} \cdot x_{\emptyset} \\
& =(-1)^{r}\left(\sum_{j} b_{j i} y_{j}\right)(-1)^{i_{1}+\cdots+i_{r}} \sum_{K} \operatorname{Pf}_{K}(A(I)) y_{K} \\
& =(-1)^{r+i_{1}+\cdots+i_{r}} \sum_{j} \sum_{K \ngtr j} b_{j i} \operatorname{Pf}_{K}(A(I)) y_{j} \cdot y_{K} \\
& =\cdots \\
& =(-1)^{i+i_{1}+\cdots+i_{r}} \sum_{K} \operatorname{Pf}_{K}(A(I \cup\{i\})) y_{K},
\end{aligned}
$$

as desired.

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# Positive Level, Negative Level and Level Zero 

Finn McGlade, Arun Ram, and Yaping Yang


#### Abstract

This is a survey on the combinatorics and geometry of integrable representations of quantum affine algebras with a particular focus on level 0 . Pictures and examples are included to illustrate the affine Weyl group orbits, crystal graphs and Macdonald polynomials that provide detailed understanding of the structure of the extremal weight modules and their characters. The final section surveys the alcove walk method of working with the positive level, negative level and level zero affine flag varieties and describes the corresponding actions of the affine Hecke algebra.


Keywords Affine flag varieties • Integrable representations • Quantum affine algebras

AMS Subject Classifications Primary 17B37 • Secondary 17B67

## 1 Introduction

This paper is about positive level, negative level and level 0 . It was motivated by the striking result of [23], which establishes a Pieri-Chevalley formula for the K-theory of the semi-infinite (level 0) affine flag variety. This made us want to learn more about the level 0 integrable modules of quantum affine algebras. Our trek brought us face to face with a huge literature, including important contributions from Drinfeld, Kashiwara, Beck, Chari, Nakajima, Lenart-Schilling-Shimozono, Cherednik-Orr,

Dedicated to I. G. Macdonald and A. O. Morris.

[^17]Naito-Sagaki, Feigin-Makedonskyi, Kato, their coauthors and many others. It is a beautiful theory and we count ourselves lucky to have been drawn into it.

The main point is that the integrable modules for quantum affine algebras $\mathbf{U}$ naturally partition themselves into families: positive level, negative level and level zero. Their structure is shadowed by the orbits of the affine Weyl group on the lattice of weights for the affine Lie algebra, which take the shape of a concave up paraboloid at positive level, a concave down paraboloid at negative level and a tube at level 0 . These integrable modules have crystal bases which provide detailed control of their characters. At level 0 the characters are (up to a factor similar to a Weyl denominator) Macdonald polynomials specialised at $t=0$. The next amazing feature is that there are Borel-Weil-Bott theorems for each case: positive level, negative level and level 0 , where, respectively, the appropriate geometry is a positive level (thin) affine flag variety, a negative level (thick) affine flag variety, and a level zero (semi-infinite) affine flag variety.

This paper is a survey of the general picture of positive level, negative level, and level zero, in the context of the combinatorics of affine Weyl groups and crystals, of the representation theory of integrable modules for quantum affine algebras, and of the geometry of affine flag varieties. In recent years, the picture has become more and more rich and taken clearer focus. We hope that this paper will help to bring this story to a wider audience by providing pictures and some explicit small examples for $\widehat{\mathfrak{s l}}_{2}$ and $\widehat{\mathfrak{s l}}_{3}$.

### 1.1 Orbits of the Affine Weyl Group W Action on $\mathfrak{h}^{*}$

For $\widehat{\mathfrak{s l}}_{2}$ the vector space $\mathfrak{h}^{*}$ is three dimensional with basis $\left\{\delta, \omega_{1}, \Lambda_{0}\right\}$, and the orbits of the action of the affine Weyl group $W^{\text {ad }}$ on $\mathfrak{h}^{*}$ on different levels look like


Although informative, the picture above is misleading as it is a two dimensional projection $(\bmod \delta)$ of what is actually going on. The $W^{\text {ad }}$-action fixes the level (the coefficient of $\Lambda_{0}$ ) but it actually changes the $\delta$ coordinate significantly. Let us look at the orbits in $\delta$ and $\omega_{1}$ coordinates (i.e. $\bmod \Lambda_{0}$ ).


The positive level orbit $W^{\text {ad }}\left(\omega_{1}+2 \Lambda_{0}\right) \bmod \Lambda_{0}$

When the level (coefficient of $\Lambda_{0}$ ) is large the parabola is wide, and it gets tighter as the level decreases.


At level 0 , the parabola pops and becomes two straight lines.


At negative level the parabola forms again, but this time facing the opposite way, and getting wider as the level gets more and more negative.


The orbit $W^{\text {ad }}\left(-\omega_{1}-2 \Lambda_{0}\right)($ negative level $) \bmod \Lambda_{0}$

The three different Bruhat orders on the affine Weyl group are visible on the $W^{\text {ad }}$-orbits:

$$
\begin{array}{ll}
v \leq w & \text { if } v\left(\omega_{1}+\Lambda_{0}\right) \text { is higher than }\left(\omega_{1}+\Lambda_{0}\right) \text { in } W^{\text {ad }}\left(\omega_{1}+\Lambda_{0}\right) \\
v \leq 0 & \text { if } v \omega_{1} \text { is higher than } w\left(\omega_{1}+0 \Lambda_{0}\right) \text { in } W^{\text {ad }}\left(\omega_{1}+0 \Lambda_{0}\right) \\
v \leq w & \text { if } v\left(-\omega_{1}-\Lambda_{0}\right) \text { is higher than } w\left(-\omega_{1}-\Lambda_{0}\right) \text { in } W^{\text {ad }}\left(-\omega_{1}-\Lambda_{0}\right)
\end{array}
$$

The definitions of the Bruhat orders on $W^{\text {ad }}$ and their relation to the closure order for Schubert cells in affine flag varieties is made precise in Sect. 2.3. Indicative relations illustrating the from of the Hasse diagrams of the positive level, negative level, and level zero Bruhat orders for the Weyl group of $\widehat{\mathfrak{s l}}_{3}$ are pictured in Plate $A$ (there are
additional relations which are not displayed in the pictures-in an effort to make the periodicity pattern easily visible).

For the case of $\mathfrak{g}=\mathfrak{s l}_{3}$ the affine Weyl group orbits take a similar form, with the points sitting on a downward paraboloid at positive level, on an upward paraboloid at negative level and with the paraboloid popping and becoming a tube at level 0 (for an example tube see the picture of $B\left(\omega_{1}+\omega_{2}\right)$ for $\widehat{\mathfrak{s}}_{3}$ in Plate D).


Positive level orbit $W^{\text {ad }}\left(\omega_{1}+\omega_{2}+2 \Lambda_{0}\right)$ for $\widehat{\mathfrak{s l}}_{3}$ Negative level orbit $W^{\text {ad }}\left(-\omega_{1}-\omega_{2}-2 \Lambda_{0}\right)$ for $\widehat{\mathfrak{s}}_{3}$

### 1.2 Extremal Weight Modules L( $\boldsymbol{\Lambda}$ ) and Their Crystals B( $\Lambda$ )

For the affine Lie algebra $\mathfrak{g}=\widehat{\mathfrak{s}}_{2}$ the weights of integrable $\mathfrak{g}$-modules always lie in the set

$$
\mathfrak{h}_{\mathbb{Z}}^{*}=\mathbb{C} \delta+\mathbb{Z} \omega_{1}+\mathbb{Z} \Lambda_{0}
$$

A set of representatives for the orbits of the action of $W^{\text {ad }}$ on $\mathfrak{h}_{\mathbb{Z}}^{*}$ is

$$
\left(\mathfrak{h}^{*}\right)_{\text {int }}=\left(\mathfrak{h}^{*}\right)_{\text {int }}^{+} \cup\left(\mathfrak{h}^{*}\right)_{\text {int }}^{0} \cup\left(\mathfrak{h}^{*}\right)_{\text {int }}^{-}, \quad \begin{aligned}
& \left(\mathfrak{h}^{*}\right)_{\text {int }}^{0}=\left\{a \delta+n \omega_{1} \in \mathfrak{h}_{\mathbb{Z}}^{*} \mid 0 \leq n\right\}, \\
& \text { where }
\end{aligned} \begin{aligned}
& \left(\mathfrak{h}^{*}\right)_{\text {int }}^{+}=\left\{a \delta+m \omega_{1}+n \Lambda_{0} \in \mathfrak{h}_{\mathbb{Z}}^{*} \mid 0 \leq m \leq n\right\}, \\
& \left(\mathfrak{h}^{*}\right)_{\text {int }}^{-}=\left\{a \delta-m \omega_{1}-n \Lambda_{0} \in \mathfrak{h}_{\mathbb{Z}}^{*} \mid 0 \leq m \leq n\right\} .
\end{aligned}
$$

These sets are illustrated $(\bmod \delta)$ below.
For each of the $\Lambda \in\left(\mathfrak{h}^{*}\right)_{\mathrm{int}}$, there is a (universal) integrable extremal weight module $L(\Lambda)$, which is highest weight if $\Lambda \in\left(\mathfrak{h}^{*}\right)_{\text {int }}^{+}$, is lowest weight (and not highest weight) if $\Lambda \in\left(\mathfrak{h}^{*}\right)_{\text {int }}^{-}$, and which is neither highest or lowest weight when $\Lambda \in\left(\mathfrak{h}^{*}\right)_{\text {int }}^{0}$. The module $L(\Lambda)$ has a crystal $B(\Lambda)$.

At positive level and negative levels the crystals $B(\Lambda)$ are connected, but the crystal $B(\Lambda)$ is usually not connected in level 0 . The connected components and their structure are known explicitly from a combination of results of Kashiwara, Beck-Chari-Pressley, Nakajima, Beck-Nakajima, Fourier-Littelmann, Ion and others. These results are collected in Theorem 5 and Eq. (28) expresses the characters of the $L(\Lambda)$ in terms of Macdonald polynomials specialised at $t=0$.


### 1.3 Affine Flag Varieties $G / I^{+}, G / I^{0}$ and $G / I^{-}$

There are three kinds of affine flag varieties for the loop group $G=G(\mathbb{C}((\epsilon)))$ : the positive level (thin) affine flag variety $G / I^{+}$, the negative level (thick) affine flag variety $G / I^{-}$and the level 0 (semi-infinite) affine flag variety $G / I^{0}$. A combination of results of Kumar, Mathieu, Kashiwara, Kashiwara-Tanisaki, KashiwaraShimozono, Varagnolo-Vasserot, Lusztig and Braverman-Finkelberg have made it clear that there is a Borel-Weil-Bott theorem for each of these:

$$
\begin{aligned}
H^{0}\left(G / I^{+}, \mathcal{L}_{\Lambda}\right) \cong L(\Lambda), & \text { for positive level } \Lambda \in\left(\mathfrak{h}^{*}\right)_{\mathrm{int}}^{+}, \\
H^{0}\left(G / I^{0}, \mathcal{L}_{\lambda}\right) \cong L(\lambda), & \text { for level zero } \lambda \in\left(\mathfrak{h}^{*}\right)_{\mathrm{int}}^{0}, \\
H^{0}\left(G / I^{-}, \mathcal{L}_{-\Lambda}\right) \cong L(-\Lambda), & \text { for negative }-\Lambda \in\left(\mathfrak{h}^{*}\right)_{\mathrm{int}}^{-} .
\end{aligned}
$$

These Borel-Weil-Bott theorems tightly connect the representation theory with the geometry. In all essential aspects the combinatorics of the positive level affine flag variety and the loop Grassmannian generalizes to the negative level and the level 0 affine flag varieties.

Section 6 extends the results of [38] and displays the alcove walk combinatorics for each of the three cases (positive level, negative level and level 0 ) in parallel. In addition it describes the method for deriving the natural affine Hecke algebra actions on the function spaces $C\left(G / I^{+}\right), C\left(I^{+} \backslash G / I^{+}\right), C\left(I^{+} \backslash G / I^{+}\right)$and $C\left(I^{-} \backslash G / I^{+}\right)$.

### 1.4 References and Technicalities

Section 2.1 introduces the affine Lie algebra and the homogeneous Heisenberg subalgebra following [17] and Sect. 2.2 gives explicit matrices describing the actions of the affine Weyl group $W^{\text {ad }}$ on the affine Cartan $\mathfrak{h}$ and its dual $\mathfrak{h}^{*}$. Section 2.3 defines the Bruhat orders on the affine Weyl group and explains their relation to the corresponding affine flag varieties following [24] and [27, Sects. 7 and 11]. Sections 2.4 and 2.5 introduce the affine braid groups and Macdonald polynomials following [40]. Section 2.6 treats the specializations of (nonsymmetric) Macdonald polynomials at $t=0, t=\infty, q=0$ and $q=\infty$ and reviews the result of Ion [15] that relates Macdonald polynomials at $t=0$ to Demazure operators.

Section 2 follows [1,3, 4], introducing the quantum affine algebra $\mathbf{U}$, the conversion to its loop presentation, the PBW-type elements and the quantum homogeneous Heisenberg subalgebra. Section 3.2 defines integrable U-modules and Sect. 3.3 introduces the extremal weight modules $L(\Lambda)$ following [19, (8.2.2)] and [20, Sect.3.1]. Section 3.4 reviews the Demazure character formulas for extremal weight modules. Section 3.5 discusses the loop presentation of the level 0 extremal weight modules and the fact that these coincide with the universal standard modules of [34] and the global Weyl modules of [9]. Letting $\mathbf{U}^{\prime}$ be $\mathbf{U}$ without the element $D$, Sect. 3.6 explains how to shrink the extremal weight module to a local Weyl module and how this provides a classification of finite dimensional simple $\mathbf{U}^{\prime}$-modules by Drinfeld polynomials.

We have made a concerted effort to make a useful survey. In order to simplify the exposition we have brushed under the rug a number of technicalities which are wisely ignored when one learns the subject (particularly (a) the difference between simply laced cases and the general case requires proper attention to the diagonal matrix which symmetrizes the affine Cartan matrix [17, (2.1.1)] causing the constants $d_{i}$ which pepper the quantum group literature and (b) the machinations necessary for allowing multiple parameters $t_{i}^{\frac{1}{2}}$ in Macdonald polynomials). The reader who needs to sort out these features is advised to drink a strong double espresso to optimise clear thinking, consult the references (particularly [4, 40]) and not trust our exposition. Perhaps in the future a more complete (probably book length) version of this paper will be completed which allows us to attend more carefully to these nuances and include more detailed proofs. Having made this point, we can say that a careful effort has been made to provide specific references to the literature at every step and we hope that this will be useful for the reader that wishes to go further.

## PLATE A: Bruhat orders on the affine Weyl group (partial, indicative, relations)


postive level Bruhat order for $\widehat{\mathfrak{s l}}_{2}$ 1 is minimal

level zero Bruhat order for $\widehat{\mathfrak{s l}}_{2}$ translation invariant

postive level Bruhat order for $\widehat{\mathfrak{s l}}_{3}$ 1 is minimal

level zero Bruhat order for $\widehat{\mathfrak{s l}}_{3}$ translation invariant

negative level Bruhat order for $\widehat{\mathfrak{s l}}_{2}$
1 is maximal

negative level Bruhat order for $\widehat{\mathfrak{s}}_{3}$
1 is maximal

PLATE B: Pictures of $B\left(\omega_{1}+\Lambda_{0}\right), B\left(\omega_{1}+0 \Lambda_{0}\right)$ and $B\left(-\omega_{1}-\Lambda_{0}\right)$ for $\widehat{\mathfrak{s}}_{2}$


Initial portion of the crystal graph of $B\left(\omega_{1}+\Lambda_{0}\right)$ for $\widehat{\mathfrak{s}}_{2}$


Final portion of the crystal graph of $B\left(-\omega_{1}-\Lambda_{0}\right)$ for $\widehat{\mathfrak{s l}}_{2}$


Middle portion of the crystal graph of $B\left(\omega_{1}+0 \Lambda_{0}\right)$ for $\mathfrak{g}=\widehat{\mathfrak{s l}}_{2}$
PLATE C: Pictures for $B\left(2 \omega_{1}\right)$. Representative paths from the (first five) connected components of $B\left(2 \omega_{1}\right)$ are

and the paths in $B\left(2 \omega_{1}\right)_{0} \subseteq B\left(\omega_{1}\right) \otimes B\left(\omega_{1}\right)$ are



PLATE D: Pictures and characters of $B\left(\omega_{1}+\omega_{2}\right)$ for $\widehat{\mathfrak{s}}_{3}$. The colour red indicates change in the $\delta$-coordinate.

$\begin{aligned} \operatorname{gchar}\left(B^{\mathrm{fin}}\left(\omega_{1}+\omega_{2}\right)\right)= & X^{\rho}+X^{\alpha_{1}}+X^{\alpha_{2}}+X^{-\alpha_{1}} \\ & +X^{-\alpha_{2}}+X^{-\rho}+2+q^{-1}\end{aligned}$

$$
+X^{-\alpha_{2}}+X^{-\rho}+2+q^{-1}
$$


$B\left(\omega_{1}+\omega_{2}\right)$ crystal graph

At $t=0$ and $t=\infty$ the normalized nonsymmetric Macdonald polynomials $\tilde{E}_{s_{1} s_{2} s_{1} \rho}(q, t)$ are

$$
\begin{aligned}
\tilde{E}_{s_{1} s_{1} \rho}(q, 0) & =X^{s_{1} s_{2} s_{1} \rho}+X^{s_{1} s_{2} \rho}+X^{s_{2} s_{1} \rho}+X^{s_{2} \rho}+X^{s_{1} \rho}+X^{\rho}+2+q, \\
\tilde{E}_{s_{1} s_{2} s_{1} \rho}(q, \infty) & =X^{s_{1} s_{2} s_{1} \rho}+q^{-1}\left(X^{s_{1} s_{2} \rho}+X^{s_{2} s_{1} \rho}\right)+q^{-2}\left(X^{s_{1} \rho}+X^{s_{2} \rho}+X^{\rho}\right)+2 q^{-2}+q^{-1} .
\end{aligned}
$$

Letting $q=e^{-\delta}$ (as in [17, (12.1.9)]), the Demazure module $L\left(\omega_{1}+\omega_{2}\right)_{\leq s_{1} s_{2} s_{1}}$ has character

$$
\begin{aligned}
\operatorname{char}\left(L\left(\omega_{1}+\omega_{2}\right)_{\leq s_{1} s_{2} s_{1}}\right) & =\frac{1}{\left(1-q^{-1}\right)^{2}} \tilde{E}_{s_{1} s_{2} s_{1} \rho}\left(q^{-1}, 0\right) \text { and } \\
\operatorname{char}\left(L\left(\omega_{1}+\omega_{2}\right)\right) & =0_{q} 0_{q} \tilde{E}_{s_{1} s_{2} s_{1} \rho}\left(q^{-1}, 0\right)=0_{q} 0_{q} \tilde{E}_{s_{1} s_{2} s_{1} \rho}\left(q^{-1}, \infty\right)
\end{aligned}
$$

where $0_{q}=\cdots+q^{-3}+q^{-2}+q^{-1}+1+q+q^{2}+\cdots$ as in Remark 2.
Remark 1 The expansion

$$
\frac{1}{\left(1-q^{-1}\right)^{2}}=\left(1+q^{-1}+q^{-2}+\cdots\right)\left(1+q^{-1}+q^{-2}+\cdots\right)=1+2 q^{-1}+3 q^{-2}+4 q^{-3}+5 q^{-4}+\cdots
$$

show that the sizes of the weight spaces of $L\left(\omega_{1}+\omega_{2}\right)_{\leq s_{1} s_{2} s_{1}}$ are growing as $\delta$ increases. Similarly, in the character formula of $L\left(\omega_{1}+\omega_{2}\right)$, the factor $0_{q} 0_{q}$ has coefficient of $q^{n}$ equal to $\operatorname{Card}\left(\left\{\left(k_{1}, k_{2}\right) \in \mathbb{Z}^{2} \mid k_{1}+k_{2}=n\right\}\right)=\infty$. This shows that every weight space of the extremal weight module $L\left(\omega_{1}+\omega_{2}\right)$ is infinite dimensional.

## 2 Affine Weyl Groups, Braid Groups and Macdonald Polynomials

### 2.1 The Affine Lie Algebra $\mathfrak{g}$

Let $\mathfrak{g}$ be a finite dimensional complex semisimple Lie algebra and fix a Cartan subalgebra $\mathfrak{a} \subseteq \mathfrak{g}$ and a symmetric, ad-invariant, nondegenerate, bilinear form $\langle\rangle:, \mathfrak{g} \times \mathfrak{g} \rightarrow \mathbb{C}$. The affine Kac-Moody algebra is

$$
\begin{gather*}
\mathfrak{g}=\left(\bigoplus_{k \in \mathbb{Z}} \mathfrak{g} \epsilon^{k}\right) \oplus \mathbb{C} K \oplus \mathbb{C} d, \quad \text { with bracket given by }\left[K, x \epsilon^{k}\right]=0, \quad[K, d]=0, \\
{\left[d, x \epsilon^{k}\right]=k x \epsilon^{k}, \quad \text { and } \quad\left[x \epsilon^{k}, y \epsilon^{\ell}\right]=[x, y] \epsilon^{k+\ell}+k \delta_{k,-\ell}\langle x, y\rangle K,} \tag{1}
\end{gather*}
$$

for $x, y \in \mathfrak{g}$ and $k, \ell \in \mathbb{Z}$ (see [17, (7.2.2)]). Let $\theta$ be the highest root of $\mathfrak{g}$ (the highest weight of the adjoint representation) and define

$$
e_{0}=f_{\theta} \epsilon, \quad f_{0}=e_{\theta} \epsilon^{-1}, \quad \text { and } \quad h_{0}=\left[e_{0}, f_{0}\right]=-h_{\theta}+K
$$

The miracle is that $\mathfrak{g}$ is a Kac-Moody Lie algebra with Chevalley generators

$$
\begin{equation*}
e_{0}, \ldots, e_{n}, h_{0}, \ldots h_{n}, d, f_{0}, \ldots f_{n}, \quad \text { which satisfy Serre relations. } \tag{2}
\end{equation*}
$$

Because of (2), $\mathfrak{g}$ has a corresponding quantum enveloping algebra $\mathbf{U}=U_{q} \mathfrak{g}$.
The Cartan subalgebra of $\mathfrak{g}$ is

$$
\mathfrak{h}=\mathfrak{a} \oplus \mathbb{C} K \oplus \mathbb{C} d, \quad \text { where } \mathfrak{a} \subseteq \mathfrak{g} \text { is the Cartan subalgebra of } \mathfrak{g}
$$

Let $\stackrel{\circ}{R}^{+}$be the set of positive roots for $\dot{\mathfrak{g}}$. For $\alpha \in \stackrel{\circ}{R}^{+}, k \in \mathbb{Z}, \ell \in \mathbb{Z}_{\neq 0}$ and $i \in$ $\{1, \ldots, n\}$, let

$$
x_{\alpha+k \delta}=e_{\alpha} \epsilon^{k}, \quad x_{-\alpha+k \delta}=f_{\alpha} \epsilon^{k}, \quad h_{i, \ell}=h_{i} \epsilon^{\ell} .
$$

The homogeneous Heisenberg subalgebra (see [17, Sects. 8.4 and 14.8]) is

$$
\begin{equation*}
\mathbb{C} K \oplus \mathfrak{a}\left[\epsilon, \epsilon^{-1}\right] \quad \text { with } \quad\left[h_{i} \epsilon^{k}, h_{j} \epsilon^{\ell}\right]=k \delta_{k,-\ell} \frac{2}{\left\langle\alpha_{i}, \alpha_{i}\right\rangle} \alpha_{i}\left(h_{j}\right) K \tag{3}
\end{equation*}
$$

and $\mathfrak{a}[\epsilon]$ is a commutative Lie algebra with basis $\left\{h_{i} \epsilon^{k} \mid i \in\{1, \ldots, n\}, k \in \mathbb{Z}_{\geq 0}\right\}$.

### 2.2 The Affine Weyl Group $W^{\text {ad }}$ and Its Action on $\mathfrak{h}^{*}$ and $\mathfrak{h}$

Let $\delta, \omega_{1}, \ldots, \omega_{n}, \Lambda_{0}$ be the basis in $\mathfrak{h}^{*}$ which is the dual basis to the basis $d, h_{1}, \ldots, h_{n}, K$ of $\mathfrak{h}$. The affine Weyl group $W^{\text {ad }}$ is the subgroup of $G L\left(\mathfrak{h}^{*}\right)$ generated by the linear transformations $s_{0}, s_{1}, \ldots, s_{n}$ which, in the basis $\delta, \omega_{1}, \ldots, \omega_{n}, \Lambda_{0}$, are
and, writing $\theta=a_{1} \alpha_{1}+\cdots+a_{n} \alpha_{n}$ and $h_{\theta}=\left[e_{\theta}, f_{\theta}\right]=a_{1}^{\vee} h_{1}+\cdots+a_{n}^{\vee} h_{n}$,

$$
s_{0}=\left(\begin{array}{cccccc}
1 & a_{1}^{\vee} & a_{2}^{\vee} & \cdots & a_{n}^{\vee} & -1  \tag{5}\\
0 & 1-a_{1} a_{1}^{\vee} & -a_{1} a_{2}^{\vee} & \cdots & -a_{1} a_{n}^{\vee} & a_{1} \\
0 & -a_{2} a_{1}^{\vee} & 1-a_{2} a_{2}^{\vee} & \cdots & -a_{2} a_{n}^{\vee} & a_{2} \\
\vdots & & & \vdots & & \vdots \\
0 & -a_{n} a_{1}^{\vee} & -a_{n} a_{2}^{\vee} & \cdots & 1-a_{n} a_{n}^{\vee} & a_{n} \\
0 & 0 & 0 & \cdots & 0 & 1
\end{array}\right) .
$$

Let $\mathfrak{a}_{\mathbb{R}}^{*}=\mathbb{R}-\operatorname{span}\left\{\alpha_{1}, \ldots, \alpha_{n}\right\}$. An alcove is a fundamental region for the action of $W^{\text {ad }}$ on $\left(\mathbb{R} \delta+\mathfrak{a}_{\mathbb{R}}^{*}+\Lambda_{0}\right) / \mathbb{R} \delta$. As explained (for example) in [39, 40], there is a bijection

$$
\begin{align*}
W^{\text {ad }} & \text { \{alcoves }\} \\
1 & \longmapsto\left\{x+\Lambda_{0} \in \mathfrak{a}_{\mathbb{R}}^{*}+\Lambda_{0} \mid x\left(h_{i}\right)>0 \text { for } i \in\{0, \ldots, n\}\right\} \tag{6}
\end{align*}
$$

Let $\mathfrak{a}_{\mathbb{Z}}^{\text {ad }}=\mathbb{Z}$-span $\left\{h_{1}, \ldots, h_{n}\right\}$. The finite Weyl group $W_{\text {fin }}$ is generated by $s_{1}, \ldots, s_{n}$. The translation presentation of the affine Weyl group is

$$
W^{\text {ad }}=\mathfrak{a}_{\mathbb{Z}}^{\text {ad }} \rtimes W_{\text {fin }}=\left\{t_{\mu^{\vee}} u \mid \mu^{\vee} \in \mathfrak{a}_{\mathbb{Z}}^{\text {ad }}, u \in W_{\text {fin }}\right\} \quad \text { with } \quad \begin{align*}
& t_{\mu^{\vee}} t_{\nu^{\vee}}=t_{\mu^{\vee}+\nu^{\vee}} \text { and }  \tag{7}\\
& u t_{\mu^{\vee}}=t_{u \mu^{\vee}} u,
\end{align*}
$$

for $\mu^{\vee}, \nu^{\vee} \in \mathfrak{a}_{\mathbb{Z}}^{\text {ad }}$ and $u \in W_{\text {fin }}$.
Let $\alpha_{i}^{\vee}$ be the image of $h_{i}$ under the isomorphism $\mathfrak{a} \xrightarrow{\sim} \mathfrak{a}^{*}$ coming from the nondegenerate bilinear form on $\mathfrak{a}$ which is the restriction of the nondegerate bilinear
form on $\mathfrak{g}$. In matrix form with respect to the basis $\delta, \omega_{1}, \ldots, \omega_{n}, \Lambda_{0}$ of $\mathfrak{h}^{*}$ the action of $W^{\text {ad }}$ on $\mathfrak{h}^{*}$ is given by
so that $-\frac{1}{2}\left\langle\mu^{\vee}, \mu^{\vee}\right\rangle=-\frac{1}{2}\left(\mu_{1}^{\vee} k_{1}+\cdots \mu_{n}^{\vee} k_{n}\right)$. (In (8) $d_{1}, \ldots, d_{n}$ are the minimal positive integers such that the product of the diagonal matrix $\operatorname{diag}\left(d_{1}, \ldots, d_{n}\right)$ with the Cartan matrix is symmetric, see [17, (2.1.1)].)

The basis $\left\{d, h_{1}, \ldots, h_{n}, K\right\}$ of $\mathfrak{h}$ is the dual basis to the basis $\left\{\delta, \omega_{1}, \ldots, \omega_{n}, \Lambda_{0}\right\}$ of $\mathfrak{h}^{*}$. Using the $W^{\text {ad }}$-action on $\mathfrak{h}$ given by

$$
s_{i} \mu^{\vee}=\mu^{\vee}-\alpha_{i}\left(\mu^{\vee}\right) h_{i}, \quad \text { for } i \in\{0, \ldots, n\} \text { and } \mu^{\vee} \in \mathfrak{h},
$$

the matrices for the action of $s_{0}, s_{1}, \ldots, s_{n}$ on $\mathfrak{h}$, in the basis $\left\{d, h_{1}, \ldots, h_{n}, K\right\}$, are the transposes of the matrices in (4) and (5).

### 2.3 The Positive Level, Negative Level and Level 0 Bruhat Orders on $W^{\text {ad }}$

In the framework of Sect. 6, where $G=\dot{G}(\mathbb{C}((\epsilon)))$ is the loop group, the closure orders for the Schubert cells in the positive level (thin) affine flag variety $G / I^{+}$, the negative level (thick) affine flag variety $G / I^{-}$, and the level 0 (semi-infinite) affine flag variety $G / I^{0}$ give partial orders on the affine Weyl group $W^{\text {ad }}$ :
$\overline{I^{+} w I^{+}}=\bigsqcup_{x \leq \pm w} I^{+} x I^{+}, \quad \overline{I^{+} w I^{0}}=\bigsqcup_{x \leq 0} I^{+} x I^{0}, \quad \overline{I^{+} w I^{-}}=\bigsqcup_{x \leq w} I^{+} x I^{-}$.
These orders can be described combinatorially as follows.
An element $w \in W^{\text {ad }}$ is dominant if

$$
w\left(\rho+\Lambda_{0}\right) \in \mathbb{R}_{\geq 0}-\operatorname{span}\left\{\omega_{1}, \ldots, \omega_{n}\right\}+\Lambda_{0}, \quad \text { where } \quad \rho=\omega_{1}+\cdots+\omega_{n}
$$

In the identification (6) of elements of $W^{\text {ad }}$ with alcoves, the dominant elements of $W^{\text {ad }}$ are the alcoves in the dominant Weyl chamber.

Let $x, w \in W^{\text {ad }}$ and let $w=s_{i_{1}} \cdots s_{i_{\ell}}$ be a reduced word for $w$ in the generators $s_{0}, \ldots, s_{n}$. The positive level Bruhat order on $W^{\text {ad }}$ is defined by
$x \leq \pm w$ if $x$ has a reduced word which is a subword of $w=s_{i_{1}} \cdots s_{i_{\ell}}$
The negative level Bruhat order on $W^{\text {ad }}$ is defined by $x \leq w$ if $x \geq w$.
The level 0 Bruhat order on $W^{\text {ad }}$ is determined by
(a) $\leq 0$ for dominant elements: If $x, w$ are dominant then $x \leq 0 w$ if and only if $x \leq \pm w$,
(b) $\leq 0$ translation invariance: If $\mu^{\vee} \in \mathfrak{a}_{\mathbb{Z}}^{\text {ad }}$ and $x, w \in W$ then $x \leq 0 w$ if and only if $x t_{\mu^{\vee}} \leq 0 w t_{\mu^{\vee}}$.
The positive level length is $\ell^{+}: W^{\text {ad }} \rightarrow \mathbb{Z}_{\geq 0}$ given by $\ell^{+}(w)=$ (length of a reduced word for $w$ ).
The negative level length is $\ell^{-}: W^{\text {ad }} \rightarrow \mathbb{Z}_{\leq 0}$ given by $\ell^{-}(w)=-\ell^{+}(w)$.
The level 0 length is $\ell^{0}: W^{\text {ad }} \rightarrow \mathbb{Z}$ given by
$\ell^{0}(w)=\ell^{+}(w)$ if $w$ is dominant $\quad$ and $\quad \ell^{0}\left(x t_{\mu^{\vee}}\right)-\ell^{0}\left(y t_{\mu^{\vee}}\right)=\ell^{0}(x)-\ell^{0}(y)$,
for $x, y \in W^{\text {ad }}$ and $\mu^{\vee} \in \mathfrak{a}_{\mathbb{Z}}^{\text {ad }}$. Using the formula for $\ell^{+}$given in [30, (2.8)], gives a formula for $\ell^{0}$,

$$
\begin{equation*}
\ell^{0}\left(u t_{\mu^{\vee}}\right)=\ell^{+}(u)+2\left\langle\rho, \mu^{\vee}\right\rangle, \quad \text { for } u \in W_{\mathrm{fin}}, \mu^{\vee} \in \mathfrak{a}_{\mathbb{Z}}^{\text {ad }} \tag{9}
\end{equation*}
$$

The length functions $\ell^{+}, \ell^{-}$and $\ell^{0}$ return, respectively, the dimension, the codimension and the relative dimension of Schubert cells in the positive level, negative level and level 0 affine flag varieties.

### 2.4 The Affine Braid Groups $\mathcal{B}^{\text {sc }}$ and $\mathcal{B}^{\text {ad }}$

Let $\omega_{1}^{\vee}, \ldots, \omega_{n}^{\vee}$ be the basis of $\mathfrak{a}$ which is dual to the basis $\alpha_{1}, \ldots, \alpha_{n}$ of $\mathfrak{a}^{*}$. Let

$$
\mathfrak{a}_{\mathbb{Z}}^{\text {ad }}=\mathbb{Z}-\operatorname{span}\left\{h_{1}, \ldots, h_{n}\right\} \subseteq \mathfrak{a}_{\mathbb{Z}}^{\text {sc }}=\mathbb{Z}-\operatorname{span}\left\{\omega_{1}^{\vee}, \ldots, \omega_{n}^{\vee}\right\}
$$

The affine braid group $\mathcal{B}^{\text {ad }}$ (resp. $\mathcal{B}^{\text {sc }}$ ) is generated by $T_{1}, \ldots, T_{n}$ and $Y^{\lambda^{\vee}}, \lambda^{\vee} \in \mathfrak{a}_{\mathbb{Z}}^{\text {ad }}$ (resp. $\lambda^{\vee} \in \mathfrak{a}_{\mathbb{Z}}^{\text {sc }}$ ), with relations

$$
Y^{\lambda^{\vee}} Y^{\sigma^{\vee}}=Y^{\lambda^{\vee}+\sigma^{\vee}}, \quad \underbrace{T_{i} T_{j} \cdots}_{m_{i j} \text { factors }}=\underbrace{T_{j} T_{i} \cdots,}_{m_{i j} \text { factors }} \quad \begin{aligned}
& T_{i}^{-1} Y^{\lambda^{\vee}}=Y^{s_{i} \lambda^{\vee}} T_{i}^{-1}, \text { if }\left\langle\lambda^{\vee}, \alpha_{i}\right\rangle=0, \\
& T_{i}^{-1} Y^{\lambda^{\vee}} T_{i}^{-1}=Y^{s_{i} \lambda^{\vee}}, \text { if }\left\langle\lambda^{\vee}, \alpha_{i}\right\rangle=1,
\end{aligned}
$$

for $i \in\{1, \ldots, n\}$ and $\lambda^{\vee} \in \mathfrak{a}_{\mathbb{Z}}^{\text {ad }}$ (resp. $\left.\mathfrak{a}_{\mathbb{Z}}^{\text {sc }}\right)$ and $m_{i j}=\alpha_{i}\left(h_{j}\right) \alpha_{j}\left(h_{i}\right)$ for $i, j \in$ $\{1, \ldots, n\}$ with $i \neq j$.

### 2.5 Macdonald Polynomials

Let

$$
\mathfrak{h}_{\mathbb{Z}}^{*}=\mathbb{Z}-\operatorname{span}\left\{\delta, \omega_{1}, \ldots, \omega_{n}, \Lambda_{0}\right\} \quad \text { and } \quad \mathfrak{a}_{\mathbb{Z}}^{*}=\mathbb{Z}-\operatorname{span}\left\{\omega_{1}, \ldots, \omega_{n}\right\} .
$$

The double affine Hecke algebra $\tilde{H}$ is presented by generators $T_{0}, \ldots, T_{n}$ and $X^{\mu}$, $\mu \in \mathfrak{h}_{\mathbb{Z}}^{*}$, with relations

$$
\begin{equation*}
X^{\lambda} X^{\mu}=X^{\lambda+\mu}, \quad \underbrace{T_{i} T_{j} \cdots}_{m_{i j} \text { factors }}=\underbrace{T_{j} T_{i} \cdots}_{m_{i j} \text { factors }}, \quad T_{i}^{2}=\left(t^{\frac{1}{2}}-t^{-\frac{1}{2}}\right) T_{i}+1, \tag{10}
\end{equation*}
$$

$$
T_{i} X^{\mu}=X^{s_{i} \mu} T_{i}+\left(t^{\frac{1}{2}}-t^{-\frac{1}{2}}\right) \frac{X^{\mu}-X^{s_{i} \mu}}{1-X^{\alpha_{i}}}, \quad T_{i}^{-1} X^{\mu}=X^{s_{i} \mu} T_{i}^{-1}-\left(t^{\frac{1}{2}}-t^{-\frac{1}{2}}\right) \frac{X^{\mu}-X^{s_{i} \mu}}{1-X^{-\alpha_{i}}} .
$$

for $i \in\{0, \ldots, n\}$ and $\mu \in \mathfrak{h}_{\mathbb{Z}}^{*}$. For $w \in W^{\text {ad }}$ put

$$
Y^{w}=\left\{\begin{array}{ll}
Y^{w s_{i}} T_{i}^{-1}, & \text { if } w<0 w s_{i},  \tag{11}\\
Y^{w s_{i}} T_{i}, & \text { if } w o>w s_{i},
\end{array} \quad \text { and let } \quad Y^{\lambda^{\vee}}=Y^{t_{\lambda \vee}} \text { for } \lambda^{\vee} \in \mathfrak{a}_{\mathbb{Z}}^{\text {ad }}\right.
$$

Putting $q=X^{\delta}=Y^{-K}$ then, as an algebra over $\mathbb{C}\left[q^{ \pm 1}, t^{ \pm \frac{1}{2}}\right]$,

$$
\tilde{H} \text { has basis } \quad\left\{X^{\mu} T_{u} Y^{\lambda^{\vee}} \mid \mu \in \mathfrak{a}_{\mathbb{Z}}^{*}+\mathbb{Z} \Lambda_{0}, u \in W_{\text {fin }}, \lambda^{\vee} \in \mathfrak{a}_{\mathbb{Z}}^{\mathrm{ad}}\right\}, \quad \text { where } \quad T_{u}=T_{i_{1}} \cdots T_{i_{\ell}},
$$

for a reduced word $u=s_{i_{1}} \cdots s_{i_{\ell}}$. The affine Hecke algebra is the subalgebra $\underset{\tilde{H}}{H}$ of $\tilde{H}$ with basis $\left\{T_{u} Y^{\lambda^{\vee}} \mid u \in W_{\text {fin }}, \lambda^{\vee} \in \mathfrak{a}_{\mathbb{Z}}^{\text {ad }}\right\}$. The polynomial representation of $\tilde{H}$ is

$$
\begin{equation*}
\mathbb{C}[X]=\operatorname{Ind}_{H}^{\tilde{H}}(\mathbf{1}) \quad \text { with basis } \quad\left\{X^{\mu} \mathbf{1} \mid \mu \in \mathfrak{a}_{\mathbb{Z}}^{*}\right\} \tag{12}
\end{equation*}
$$

and $\quad Y^{K} \mathbf{1}=q^{-1} \mathbf{1}, \quad Y^{-\alpha_{i}^{\vee}} \mathbf{1}=t \mathbf{1}, \quad$ and $\quad T_{i} \mathbf{1}=t^{\frac{1}{2}} \mathbf{1}$ for $i \in\{1, \ldots, n\}$.
Let

$$
\begin{equation*}
T_{0}^{\vee}=Y^{\alpha_{0}^{\vee}} X^{-\Lambda_{0}} T_{0}^{-1} X^{\Lambda_{0}} \quad \text { and } \quad T_{i}^{\vee}=T_{i} \text { for } i \in\{1, \ldots, n\} \tag{13}
\end{equation*}
$$

The automorphism of $\tilde{H}$ given by conjugation by $X^{-\Lambda_{0}}$ is the automorphism $\tau: \tilde{H} \rightarrow$ $\tilde{H}$ of $[10,(2.8)]$. Extend $\tilde{H}$ to allow rational functions in the $Y^{\lambda^{\nu}}$. For each $i \in$ $\{0,1, \ldots, n\}$, the intertwiner $\tau_{i}^{\vee} \in \widetilde{H}$ is
$\tau_{i}^{\vee}=T_{i}^{\vee}+\frac{t^{-\frac{1}{2}}(1-t)}{1-Y^{-\alpha_{i}^{\vee}}}=\left(T_{i}^{\vee}\right)^{-1}+\frac{t^{-\frac{1}{2}}(1-t) Y^{-\alpha_{i}^{\vee}}}{1-Y^{-\alpha_{i}^{\vee}}}$ so that $Y^{\lambda^{\vee}} \tau_{i}^{\vee}=\tau_{i}^{\vee} Y^{s_{i} \lambda^{\vee}}$.
Let, for simplicity, $\mu \in \mathbb{Z}$-span $\left\{\alpha_{1}, \ldots, \alpha_{n}\right\}$ (the general case $\mu \in \mathfrak{a}_{\mathbb{Z}}^{*}$ requires consideration of the group $\Omega^{\vee}$, the quotient of the weight lattice by the root lattice, and is
treated in detail in [40]). The nonsymmetric Macdonald polynomial $E_{\mu}=E_{\mu}(q, t)$ is

$$
\begin{equation*}
E_{\mu}=E_{\mu}(q, t)=\tau_{i_{1}}^{\vee} \ldots \tau_{i_{\ell}}^{\vee} \mathbf{1}, \quad \text { where } m_{\mu}=s_{i_{1}} \ldots s_{i_{\ell}} \text { is a reduced word } \tag{15}
\end{equation*}
$$

for the minimal length element in the coset $t_{\mu} W_{\text {fin }}$. The $E_{\mu}$ form a basis of $\mathbb{C}[X]$ consisting of eigenvectors for the $Y^{\lambda^{\vee}}$ (the Cherednik-Dunkl operators).

Fix a reduced word $m_{\mu}=s_{i_{1}} \ldots s_{i_{\ell}}$ as in (15). Identifying the elements of $W^{\text {ad }}$ with alcoves as in (6), an alcove walk of type $\vec{m}_{\mu}=\left(i_{1}, \ldots, i_{\ell}\right)$ beginning at 1 (the fundamental alcove) is a sequence of steps, of types $i_{1}, \ldots, i_{\ell}$, where a step of type $j$ is (the signs - and + indicate that $z s_{j}{ }^{\circ}>z$ )


Let $\mathcal{B}\left(1, \vec{m}_{\mu}\right)$ be the set of alcove walks of type $\vec{m}_{\mu}=\left(i_{1}, \ldots, i_{\ell}\right)$ beginning at 1 . For a walk $p \in \mathcal{B}\left(1, \vec{m}_{\mu}\right)$ let

$$
\begin{aligned}
& f^{+}(p)=\{k \in\{1, \ldots, \ell\} \mid \text { the } k \text { th step of } p \text { is a positive fold }\}, \\
& f^{-}(p)=\{k \in\{1, \ldots, \ell\} \mid \text { the } k \text { th step of } p \text { is a negative fold }\}, \\
& f(p)=f^{+}(p) \cup f^{-}(p)=\{k \in\{1, \ldots, \ell\} \mid \text { the } k \text { th step of } p \text { is a fold }\} .
\end{aligned}
$$

For $p \in \mathcal{B}\left(1, \vec{m}_{\mu}\right)$ let end $(p)$ be the endpoint of $p$ (an element of $W^{\text {ad }}$ ) and define the weight $\operatorname{wt}(p)$ and the final direction $\varphi(p)$ of $p$ by

$$
X^{\mathrm{end}(p)}=X^{\mathrm{wt}(p)} T_{\varphi(p)}^{\vee}, \quad \text { with wt }(p) \in \mathfrak{a}_{\mathbb{Z}}^{*} \text { and } \varphi(p) \in W_{\text {fin }}
$$

Using (14) and doing a left to right expansion of the terms of $\tau_{i_{1}}^{\vee} \cdots \tau_{i_{\ell}}^{\vee} \mathbf{1}$ produces the monomial expansion of $E_{\mu}$ as sum over alcove walks as given in the following theorem. For simplicity we state the following theorem for $\mu \in \mathbb{Z}$-span $\left\{\alpha_{1}, \ldots, \alpha_{n}\right\}$. It holds, after a small technical adjustment to the statement, for all $\mu \in \mathfrak{a}_{\mathbb{Z}}^{*}$, see [40] for details.

Theorem 1 ([40, Theorem 3.1 and Remark 3.3]) Let $\mu \in \mathbb{Z}$-span $\left\{\alpha_{1}, \ldots, \alpha_{n}\right\}$ and let $m_{\mu}$ be the minimal length element in the coset $t_{\mu} W_{\text {fin }}$. Fix a reduced word $\vec{m}_{\mu}=$ $s_{i_{1}} \cdots s_{i_{\ell}}$, let

$$
\beta_{1}^{\vee}=s_{i_{\ell}} \cdots s_{i_{2}} \alpha_{i_{1}}^{\vee}, \quad \beta_{2}^{\vee}=s_{i_{\ell}} \cdots s_{i_{3}} \alpha_{i_{2}}^{\vee}, \quad \ldots, \quad \beta_{\ell}^{\vee}=\alpha_{i_{\ell}}^{\vee}
$$

and let $\operatorname{sh}\left(\beta_{k}^{\vee}\right)$ and $\operatorname{ht}\left(\beta_{k}^{\vee}\right)$ be defined by $Y^{\beta_{k}^{\vee}} \mathbf{1}=q^{\operatorname{sh}\left(\beta_{k}^{\vee}\right)} t^{\mathrm{ht}\left(\beta_{k}^{\vee}\right)} \mathbf{1}$ for $k \in\{1, \ldots, \ell\}$. Then

$$
E_{\mu}(q, t)=\sum_{p \in \mathcal{B}\left(1, \vec{m}_{\mu}\right)} X^{w t(p)} t^{\frac{1}{2}(\ell(\varphi(p))} \prod_{k \in f^{+}(p)} \frac{t^{-\frac{1}{2}}(1-t)}{1-q^{\operatorname{sh}\left(\beta_{k}^{\vee}\right)} t^{\mathrm{ht}\left(\beta_{k}^{\vee}\right)}} \prod_{k \in f^{-}(p)} \frac{t^{-\frac{1}{2}}(1-t) q^{\operatorname{sh}\left(\beta_{k}^{\vee}\right)} t^{\mathrm{ht}\left(\beta_{k}^{\vee}\right)}}{1-q^{\operatorname{sh}\left(\beta_{k}^{\vee}\right)} t^{\mathrm{ht}\left(\beta_{k}^{\vee}\right)}}
$$

### 2.6 Specializations of the Normalized Macdonald Polynomials $\tilde{E}_{\mu}(\boldsymbol{q}, \boldsymbol{t})$

If $m_{\mu}=t_{\mu} m$ with $m \in W_{\text {fin }}$ then $E_{\mu}(q, t)$ has top term $t^{\frac{1}{2} \ell(m)} X^{\mu}$ (this term is the term corresponding to the unique alcove walk in $\mathcal{B}\left(\vec{m}_{\mu}\right)$ with no folds). The normalized nonsymmetric Macdonald polynomial is

$$
\tilde{E}_{\mu}(q, t)=t^{-\frac{1}{2} \ell(m)} E_{\mu}(q, t) \text { so that } \tilde{E}_{\mu}(q, t) \text { has top term } X^{\mu}
$$

A path $p \in \mathcal{B}\left(1, \vec{m}_{\mu}\right)$ is positively folded if there are no negative folds, i.e. $\# f^{-}(p)=$ 0.

A path $p \in \mathcal{B}\left(1, \vec{m}_{\mu}\right)$ is negatively folded if there are no positive folds, i.e. $\# f^{+}(p)=$ 0 .

A path $p \in \mathcal{B}\left(1, \vec{m}_{\mu}\right)$ is positive semi-infinite if $\ell(\varphi(p))-\ell(m)-\# f(p)+2$ $\sum_{k \in f^{-}(p)} h t\left(\beta_{k}^{\vee}\right)=0$. A path $p \in \mathcal{B}\left(1, \vec{m}_{\mu}\right)$ is negative semi-infinite if $\ell(m)-\ell(\varphi(p))+$
$\# f(p)+2 \sum_{k \in f^{+}(p)} h t\left(\beta_{k}^{\vee}\right)=0$.
Proposition 1 Let $\mu \in \mathfrak{a}_{\mathbb{Z}}^{*}$. The specializations $q=0, t=0, q^{-1}=0$ and $t^{-1}=0$ are well defined and given, respectively, by

$$
\begin{aligned}
& \tilde{E}_{\mu}(0, t)=\sum_{\substack{p \in \mathcal{E}\left(1, \bar{m}^{\prime} \mu\right) \\
p \text { pos folded }}} t^{\frac{1}{2}(\ell(\varphi(p))-\ell(m)-\# f(p))}(1-t)^{\# f(p)} X^{\mathrm{wt}(p)}, \\
& \tilde{E}_{\mu}(q, 0)=\sum_{\substack{p \in \mathcal{B}\left(1, \bar{m}_{\mu}\right) \\
p \text { neg semi }-\mathrm{inf}}} q^{\sum_{k \in f-(p)} \operatorname{sh}\left(\beta_{k}^{\vee}\right)} X^{\mathrm{wt}(p)}, \\
& \tilde{E}_{\mu}(\infty, t)=\sum_{\substack{p \in \mathcal{B}\left(1, m_{\mu}\right) \\
p \text { neg folded }}} t^{-\frac{1}{2}(\ell(\varphi(p))-\# f(p))}\left(1-t^{-1}\right)^{\# f(p)} X^{\mathrm{wt}(p)}, \\
& \tilde{E}_{\mu}(q, \infty)=\sum_{\substack{p \in \mathcal{B}\left(1, \tilde{m}_{\mu}\right) \\
p \text { pos semi }- \text { inf }}} q^{-\sum_{k \in f+(p)} \operatorname{sh}\left(\beta_{k}^{\vee}\right)} X^{\operatorname{wt}(p)} .
\end{aligned}
$$

For $i \in\{0,1, \ldots, n\}$ and $f \in \mathbb{C}\left[\mathfrak{h}_{\mathbb{Z}}^{*}\right]$ define

$$
\begin{equation*}
\Delta_{i} f=\frac{f-s_{i} f}{1-X^{-\alpha_{i}}} \quad \text { and } \quad D_{i} f=\left(1+s_{i}\right) \frac{1}{1-X^{-\alpha_{i}}} f=\frac{f-X^{-\alpha_{i}} s_{i} f}{1-X^{-\alpha_{i}}} \tag{16}
\end{equation*}
$$

Equations (14), (13) and the last relation in (10) give that, as operators on the polynomial representation,

$$
\begin{aligned}
t^{\frac{1}{2}} \tau_{i}^{\vee} & =X^{-\Lambda_{0}}\left(D_{i}-t \Delta_{i}\right) X^{\Lambda_{0}}+\frac{(1-t) Y^{-\alpha_{i}^{\vee}}}{1-Y^{-\alpha_{i}^{\vee}}} \text { for } i \in\{1, \ldots, n\}, \quad \text { and } \\
t^{\frac{1}{2}} \tau_{0}^{\vee} Y^{\alpha_{0}^{\vee}} & =Y^{-\alpha_{0}^{\vee}} t^{\frac{1}{2}} \tau_{0}^{\vee}=X^{-\Lambda_{0}}\left(D_{0}-t \Delta_{0}\right) X^{\Lambda_{0}}+\frac{(1-t) Y^{-\alpha_{0}^{\vee}}}{1-Y^{-\alpha_{0}^{\vee}}}
\end{aligned}
$$

When applied in the formula of (15), these formulas for (normalized) intertwiners are specializable at $t=0$, giving the following result (see the examples computed in Sect.4.2).

Theorem 2 ([15, Sect.4.1]) Let $\mu \in \mathfrak{a}_{\mathbb{Z}}^{*}$. There are unique $\nu \in \mathfrak{a}_{\mathbb{Z}}^{*}$ and $j \in \mathbb{Z}$ such that

$$
\Lambda_{0}+\nu \in\left(\mathfrak{h}^{*}\right)_{\mathrm{int}}^{+} \text {and }-j \delta+\mu+\Lambda_{0} \in W^{\mathrm{ad}}\left(\Lambda_{0}+\nu\right)
$$

Let $w \in W^{\text {ad }}$ be minimal length such that $-j+\mu+\Lambda_{0}=w\left(\Lambda_{0}+\nu\right)$ and let $w=$ $s_{i_{1}} \cdots s_{i_{\ell}}$ be a reduced word. Letting $D_{0}, \ldots, D_{n}$ be the Demazure operators given in (16),

$$
\tilde{E}_{\mu}(q, 0)=q^{-j} X^{-\Lambda_{0}} D_{i_{1}} \cdots D_{i_{\ell}} X^{\nu} X^{\Lambda_{0}} \mathbf{1}
$$

## 3 Quantum Affine Algebras U and Integrable Modules

### 3.1 The Quantum Affine Algebra U

The quantum affine algebra $\mathbf{U}$ is the $\mathbb{C}(q)$-algebra generated by

$$
E_{0}, \ldots, E_{n}, F_{0}, \ldots, F_{n},, K_{0}^{ \pm 1}, \ldots, K_{n}^{ \pm 1}, C^{ \pm \frac{1}{2}}, D^{ \pm 1}
$$

with Chevalley-Serre type relations corresponding to the affine Dynkin diagram. Following [12, (11)], [28, Sect. 39], [1, end of Sect. 1], there is an action of the affine braid group $\mathcal{B}^{\text {sc }}$ on $\mathbf{U}$ by automorphisms.

Let $\mathbf{U}^{+}$be the subalgebra of $\mathbf{U}$ generated by $E_{0}, E_{1}, \ldots, E_{n}$. As explained in [3, Lemma 1.1 (iv)] and [4, (3.1)], there is a doubly infinite "longest element" for the affine Weyl group with a favourite reduced expression $w_{\infty}=\cdots s_{i_{-1}} s_{i_{0}} s_{i_{1}} \cdots$. This reduced word is used, with the braid group action, to define root vectors in $\mathbf{U}^{+}$by

$$
\begin{equation*}
E_{\beta_{0}}=E_{i_{0}}, \quad \text { and } \quad E_{\beta_{-k}}=T_{i_{0}}^{-1} T_{i_{-1}}^{-1} \cdots T_{i_{-(k-1)}}^{-1} E_{i_{-k}} \quad \text { and } \quad E_{\beta_{k}}=T_{i_{1}} T_{i_{2}} \ldots T_{i_{k-1}} E_{i_{k}} \tag{17}
\end{equation*}
$$

for $k \in \mathbb{Z}_{>0}$. For $i \in\{1, \ldots, n\}$ and $r, s \in \mathbb{Z}$ define the loop generators in $\mathbf{U}^{+}$

$$
\begin{equation*}
\mathbf{x}_{i, r}^{+}=E_{\alpha_{i}+r \delta}=Y^{-r \omega_{i}^{\vee}} E_{i} \quad \text { and } \quad \mathbf{x}_{i, s}^{-}=E_{-\alpha_{i}+s \delta}=Y^{s \omega_{i}^{\vee}} F_{i}, \tag{18}
\end{equation*}
$$

where $Y^{r \omega_{i}^{\vee}}$ and $Y^{s \omega_{i}^{\vee}}$ are elements of the braid group $\mathcal{B}^{\text {sc }}$ as defined in Sect.2.4. For $r \in \mathbb{Z}_{\geq 0}$ and $s \in \mathbb{Z}_{>0}$ these are special cases of the root vectors in (17). For $i \in\{1, \ldots, n\}$ and $r \in \mathbb{Z}_{>0}$ define $\mathbf{q}_{r}^{(i)}$ by

$$
\begin{equation*}
\mathbf{q}_{r}^{(i)}=\mathbf{x}_{i, r}^{-} \mathbf{x}_{i, 0}^{+}-q^{-2} \mathbf{x}_{i, 0}^{+} \mathbf{x}_{i, r}^{-}, \quad \text { let } \quad \mathbf{q}_{+}^{(i)}(z)=1+\left(q-q^{-1}\right) \sum_{s \in \mathbb{Z}_{>0}} \mathbf{q}_{s}^{(i)} z^{s}, \tag{19}
\end{equation*}
$$

and define $\mathbf{p}_{r}^{(i)}$ and $\mathbf{e}_{r}^{(i)}$ by

$$
\begin{equation*}
\mathbf{q}_{+}^{(i)}(z)=\exp \left(\sum_{r \in \mathbb{Z}_{>0}}\left(q-q^{-1}\right) \mathbf{p}_{r}^{(i)} z^{r}\right) \text { and } \exp \left(\sum_{r \in \mathbb{Z}_{>0}} \frac{\mathbf{p}_{r}^{(i)}}{[r]} z^{r}\right)=1+\sum_{k \in \mathbb{Z}_{>0}} \mathbf{e}_{k}^{(i)} z^{k} \tag{20}
\end{equation*}
$$

For a sequence of partitions $\vec{\kappa}=\left(\kappa^{(1)}, \ldots, \kappa^{(n)}\right)$ define

$$
\begin{equation*}
\mathbf{s}_{\vec{k}}=\mathbf{s}_{\kappa^{(1)}} \cdots \mathbf{s}_{\kappa^{(n)}}, \quad \text { where } \quad \mathbf{s}_{\kappa^{(i)}}=\operatorname{det}\left(\mathbf{e}_{\left(\kappa^{(i)}\right)_{r}^{\prime}-r+s}^{(i)}\right)_{1 \leq r, s \leq m_{i}}, \tag{21}
\end{equation*}
$$

where $\left(\kappa^{(i)}\right)_{r}^{\prime}$ is the length of the $r$ th column of $\kappa^{(i)}$ and $m_{i}=\ell\left(\kappa^{(i)}\right)$ (see [29, Chap. I (3.5)]). For a sequence $\mathbf{c}=\left(\cdots, c_{-3}, c_{-2}, c_{-1}, \vec{\kappa}, c_{1}, c_{2}, c_{3}, \ldots\right)$ with $c_{i} \in \mathbb{Z}_{\geq 0}$ and all but a finite number of $c_{i}$ equal to 0 . The corresponding $P B W$-type element of $\mathbf{U}^{+}$ is

$$
\begin{equation*}
E_{\mathbf{c}}=\left(E_{\beta_{0}}^{\left(c_{0}\right)} E_{\beta_{-1}}^{\left(c_{-1}\right)} E_{\beta_{-2}}^{\left(c_{-2}\right)} \cdots\right)\left(\mathbf{s}_{\kappa^{(1)}} \cdots \mathbf{s}_{\kappa^{(n)}}\right)\left(\cdots E_{\beta_{3}}^{\left(c_{3}\right)} E_{\beta_{2}}^{\left(c_{2}\right)} E_{\beta_{1}}^{\left(c_{1}\right)}\right) . \tag{22}
\end{equation*}
$$

The Cartan involution is the $\mathbb{C}$-linear anti-automorphism $\Omega: \mathbf{U} \rightarrow \mathbf{U}$ given by

$$
\Omega\left(E_{i}\right)=F_{i}, \quad \Omega\left(F_{i}\right)=E_{i}, \quad \Omega\left(K_{i}\right)=K_{i}^{-1}, \quad \Omega(D)=D^{-1}, \quad \Omega(q)=q^{-1}
$$

and $\mathbf{U}^{-}=\Omega\left(\mathbf{U}^{+}\right)$. Putting $\mathbf{p}_{-r}^{(i)}=\Omega\left(\mathbf{p}_{r}^{(i)}\right)$, then (see [12, Th. 2 (6)] and [1, Th. 4.7 (2)])

$$
\begin{equation*}
\left[\mathbf{p}_{k}^{(i)}, \mathbf{p}_{l}^{(i)}\right]=\delta_{k,-l} \frac{1}{k}\left(\frac{q^{k \alpha_{i}\left(h_{j}\right)}-q^{-k \alpha_{i}\left(h_{j}\right)}}{q-q^{-1}}\right) \frac{C^{k}-C^{-k}}{q-q^{-1}} \tag{23}
\end{equation*}
$$

Define $q_{-s}^{(i)}=\Omega\left(q_{s}^{(i)}\right)$ and

$$
\begin{equation*}
\mathbf{q}_{-}^{(i)}\left(z^{-1}\right)=1+\left(q-q^{-1}\right) \sum_{s \in \mathbb{Z}_{>0}} \mathbf{q}_{-s}^{(i)} z^{-s} \tag{24}
\end{equation*}
$$

The Heisenberg subalgebra $\mathbf{H}$ is the subalgebra of $\mathbf{U}$ generated by $\left\{\mathbf{p}_{k}^{(i)} \mid i \in\right.$ $\left.\{1, \ldots, n\}, k \in \mathbb{Z}_{\neq 0}\right\}$. In $\mathbf{H} \cap \mathbf{U}^{+}$the $\mathbf{p}_{r}^{(i)}\left(r \in \mathbb{Z}_{>0}\right)$ are the power sums, the $\mathbf{q}_{r}^{(i)}(r \in$ $\left.\mathbb{Z}_{>0}\right)$ are the Hall-Littlewoods and the $\mathbf{e}_{r}^{(i)}\left(r \in \mathbb{Z}_{>0}\right)$ are the elementary symmetric functions and the $\mathbf{s}_{\vec{k}}$ are the Schur functions.

### 3.2 Integrable U-Modules

As in (5), let $h_{\theta}=a_{1}^{\vee} h_{1}+\cdots+a_{n}^{\vee} h_{n}$ be the highest root of $\mathfrak{g}$ and let

$$
\Lambda_{i}=\omega_{i}+a_{i}^{\vee} \Lambda_{0}, \quad \text { for } i \in\{1, \ldots, n\},
$$

so that $\left\{\delta, \Lambda_{1}, \ldots, \Lambda_{n}, \Lambda_{0}\right\}$ is the dual basis in $\mathfrak{h}^{*}$ to the basis $\left\{d, h_{1}, \ldots, h_{n}, h_{0}\right\}$ of $\mathfrak{h}$. Let

$$
\mathfrak{h}_{\mathbb{Z}}^{*}=\left\{\Lambda \in \mathfrak{h}^{*} \mid\left\langle\Lambda, \alpha_{i}^{\vee}\right\rangle \in \mathbb{Z} \text { for } i \in\{0,1, \ldots, n\}\right\}=\mathbb{C} \delta+\mathbb{Z}-\operatorname{span}\left\{\Lambda_{0}, \ldots, \Lambda_{n}\right\} .
$$

A set of representatives for the $W^{\text {ad }}$-orbits on $\mathfrak{h}_{\mathbb{Z}}^{*}$ is

$$
\begin{array}{ll}
\left(\mathfrak{h}^{*}\right)_{\text {int }}^{*}=\left(\mathfrak{h}^{*}\right)_{\text {int }}^{+} \cup\left(\mathfrak{h}^{*}\right)_{\text {int }}^{0} \cup\left(\mathfrak{h}^{*}\right)_{\text {int }}^{-}, \quad \text { where } \delta \\
\quad\left(\mathfrak{Z}_{\geq 0}^{*}\right)_{\text {int }}^{0}-\operatorname{span}\left\{\Lambda_{0}, \ldots, \Lambda_{n}\right\}, \\
& \left(\mathfrak{h}^{*}\right)_{\text {int }}^{-}=\mathbb{C} \delta+0 \Lambda_{0}+\mathbb{Z}_{\geq 0} \leq-\operatorname{span}\left\{\Lambda_{0}, \ldots, \Lambda_{n}\right\} . \tag{25}
\end{array}
$$

For $\widehat{\mathfrak{s l}}_{2}$ these sets are pictured $(\bmod \delta)$ in (0.2).
For $i \in\{0,1, \ldots, n\}$ let $\mathbf{U}_{(i)}$ be the subalgebra of $\mathbf{U}$ generated by $\left\{E_{i}, F_{i}, K_{i}^{ \pm 1}\right\}$. An integrable $\mathbf{U}$-module is a $\mathbf{U}$-module $M$ such that if $i \in\{0, \ldots, n\}$ then
$\operatorname{Res}_{\mathbf{U}_{(i)}}^{\mathbf{U}}(M) \quad$ is a direct sum of finite dimensional $\mathbf{U}_{(i)}$-modules,
where $\operatorname{Res}_{\mathbf{U}_{(i)}}^{\mathbf{U}}(M)$ denotes the restriction of the $\mathbf{U}$-module $M$ to a $\mathbf{U}_{(i)}$-module.
Let $M$ be an integrable U-module. Following [28, Sect. 5], for each $w \in W^{\text {ad }}$ there is a linear map

$$
T_{w}: M \rightarrow M \text { such that } T_{w}(u m)=T_{w}(u) T_{w}(m),
$$

for $u \in \mathbf{U}$ and $m \in M$ (here $T_{w}(u)$ refers to the braid group action on $\mathbf{U}$ ). Thus, every integrable module $M$ is a module for the semidirect product $\mathcal{B}^{\text {ad }} \ltimes \mathbf{U}$ where $\mathcal{B}^{\text {ad }}$ is the braid group of $W^{\text {ad }}$.

### 3.3 Extremal Weight Modules L( $\boldsymbol{\Lambda}$ )

Let $\Lambda \in\left(\mathfrak{h}^{*}\right)_{\text {int }}$. Following [19, (8.2.2)] and [20, Sect.3.1], the extremal weight module $L(\Lambda)$ is the $\mathbf{U}$-module
generated by $\left\{u_{w \Lambda} \mid w \in W\right\} \quad$ with relations $\quad K_{i}\left(u_{w \Lambda}\right)=q^{\left\langle w \Lambda, \alpha_{i}^{\vee}\right\rangle} u_{w \Lambda}$,

$$
\begin{array}{cl}
E_{i} u_{w \Lambda}=0, \quad \text { and } \quad F_{i}^{\left\langle w \Lambda, \alpha_{i}^{\vee}\right\rangle} u_{w \Lambda}=u_{s_{i} w \Lambda}, & \text { if }\left\langle w \Lambda, \alpha_{i}^{\vee}\right\rangle \in \mathbb{Z}_{\geq 0},  \tag{26}\\
F_{i} u_{w \Lambda}=0, \quad \text { and } \quad E_{i}^{-\left\langle w \Lambda, \alpha_{i}^{\vee}\right\rangle} u_{w \Lambda}=u_{s_{i} w \Lambda}, & \text { if }\left\langle w \Lambda, \alpha_{i}^{\vee}\right\rangle \in \mathbb{Z}_{\leq 0}
\end{array}
$$

for $i \in\{0, \cdots, n\}$. The module $L(\Lambda)$ has a crystal $B(\Lambda)$ ([19, Prop. 8.2.2(ii)], [20, Sect. 3.1]).

- If $\Lambda \in\left(\mathfrak{h}^{*}\right)_{\text {int }}^{+}$then $L(\Lambda)$ is the simple $\mathbf{U}$-module of highest weight $\Lambda$ (see [17, (10.4.6)]).
- If $\Lambda \notin\left(\mathfrak{h}^{*}\right)_{\text {int }}^{+}$then $L(\Lambda)$ is not a highest weight module.
- If $\Lambda \in\left(\mathfrak{h}^{*}\right)_{\text {int }}^{-}$then $L(\Lambda)$ is the simple $\mathbf{U}$-module of lowest weight $\Lambda$.

The finite dimensional simple modules, denoted $L^{\text {fin }}(a(u))$ ), are integrable weight modules which are not extremal weight modules. The connection between the $L^{\text {fin }}(a(u))$ and the $L(\lambda)$ for $\lambda \in\left(\mathfrak{h}^{*}\right)_{\text {int }}^{0}$ is given by Theorem 4 below.

The module $L(\Lambda)$ is universal (see [21, Sect. 2.6], [2, Sect. 2.1], [36, Sect. 2.5]). One way to formulate this universality is to let $\mathbf{U}_{0}$ be the subalgebra generated by $K_{1}^{ \pm 1}, \ldots, K_{n}^{ \pm 1}, C^{ \pm \frac{1}{2}}, D^{ \pm 1}$, let intInd be an induction functor in the category of integrable $\mathbf{U}$-modules and write

$$
L(\Lambda)=\operatorname{intInd}_{\mathbf{U}_{0} \rtimes \mathcal{B}^{\text {ad }}}^{\mathrm{U}}(S(\Lambda)), \quad \text { where } \quad S(\Lambda)=\operatorname{span}\left\{u_{w \Lambda} \mid w \in W^{\text {ad }}\right\}
$$

is the $\mathbf{U}_{0} \rtimes \mathcal{B}^{\text {ad }}$-module with action given by $T_{i} u_{w \Lambda}=(-q)^{\left\langle w \Lambda_{i}, \alpha_{i}^{\rangle}\right\rangle} u_{s_{i} w \Lambda}$ and $K_{i} u_{w \Lambda}$ $=q^{\left\langle w \Lambda, \alpha_{i}^{\vee}\right\rangle} u_{w \Lambda}$, for $i \in\{0,1, \ldots, n\}$ and $w \in W^{\text {ad }}$.

### 3.4 Demazure Submodules $L(\Lambda)_{\leq w}$

Let $w \in W^{\text {ad }}$. The Demazure module $L(\Lambda)_{\leq w}$ is the $\mathbf{U}^{+}$-submodule of $L(\Lambda)$ given by

$$
L(\Lambda)_{\leq w}=\mathbf{U}^{+} u_{w \Lambda} \quad \text { and } \quad \operatorname{char}\left(L(\Lambda)_{\leq w}\right)=\sum_{p \in B(\Lambda)_{\leq w}} e^{\operatorname{wt}(p)}
$$

since $L(\Lambda)_{\leq w}$ has a crystal $B(\Lambda)_{\leq w}$. The $B G G$-Demazure operator on $\mathbb{C}\left[\mathfrak{h}_{\mathbb{Z}}^{*}\right]=$ $\mathbb{C}$-span $\left\{X^{\lambda} \mid \lambda \in \mathfrak{h}_{\mathbb{Z}}^{*}\right\}$ is given by

$$
D_{i}=\left(1+s_{i}\right) \frac{1}{1-X^{-\alpha_{i}}}, \quad \text { for } i \in\{0,1, \ldots, n\}
$$

Let $\Lambda \in\left(\mathfrak{h}^{*}\right)_{\text {int }}, w \in W$ and $i \in\{0,1, \ldots, n\}$.
If $\Lambda \in\left(\mathfrak{h}^{*}\right)_{\text {int }}^{+}$then $\left.\quad D_{i} \operatorname{char}\left(L(\Lambda)_{\leq w}\right)\right)= \begin{cases}\operatorname{char}\left(L(\Lambda)_{\leq s_{i} w}\right), & \text { if } s_{i} w \geq w, \\ \operatorname{char}\left(L(\Lambda)_{\leq w}\right), & \text { if } s_{i} w \leq w ;\end{cases}$
if $\lambda \in\left(\mathfrak{h}^{*}\right)_{\text {int }}^{0} \quad$ then $\quad D_{i} \operatorname{char}\left(L(\lambda)_{\leq w}\right)= \begin{cases}\operatorname{char}\left(L(\lambda)_{\leq s_{i} w}\right), & \text { if } s_{i} w 0 \geq w, \\ \operatorname{char}\left(L(\lambda)_{\leq w}\right), & \text { if } s_{i} w \leq 0 w ;\end{cases}$
if $\Lambda \in\left(\mathfrak{h}^{*}\right)_{\text {int }}^{-}$then $\quad D_{i} \operatorname{char}\left(L(\Lambda)_{\leq w}\right)= \begin{cases}\operatorname{char}\left(L(\Lambda)_{\leq s_{i} w}\right), & \text { if } s_{i} w \geq w, \\ \operatorname{char}\left(L(\Lambda)_{\leq w}\right), & \text { if } s_{i} w \leq w ;\end{cases}$
(see [24, Theorem 8.2.9], [18], [21, Sect. 2.8] and [22, Theorems 4.7 and 4.11]).

### 3.5 An Alternate Presentation for Level 0 Extremal Weight Modules

For $\lambda=m_{1} \omega_{1}+\cdots+m_{n} \omega_{n} \in\left(\mathfrak{h}^{*}\right)_{\text {int }}$ and let $x_{1,1}, \ldots, x_{m_{1}, 1}, x_{1,2}, \ldots, x_{m_{2}, 2}, \ldots$ $x_{1, n}, \ldots, x_{m_{n}, n}$ be $n$ sets of formal variables. Letting $e_{i}^{(j)}=e_{i}\left(x_{1, j}, \ldots, x_{m_{j}, j}\right)$ denote the elementary symmetric function in the variables $x_{1, j}, \ldots, x_{m_{j}, j}$ define

$$
\begin{aligned}
R G_{\lambda} & =\mathbb{C}\left[x_{1,1}^{ \pm 1}, \ldots, x_{m_{1}, 1}^{ \pm 1}\right]^{S_{m_{1}}} \otimes \cdots \otimes \mathbb{C}\left[x_{1, n}^{ \pm 1}, \ldots, x_{m_{n}, n}^{ \pm 1} \int^{S_{m_{n}}}\right. \\
& =\mathbb{C}\left[e_{1}^{(1)}, \ldots, e_{m_{1}-1}^{(1)},\left(e_{m_{1}}^{(1)}\right)^{ \pm 1}\right] \otimes \cdots \otimes \mathbb{C}\left[e_{1}^{(n)}, \ldots, e_{m_{n}-1}^{(n)},\left(e_{m_{n}}^{(n)}\right)^{ \pm 1}\right] \\
R G_{\lambda}^{+} & =\mathbb{C}\left[x_{1,1}, \ldots, x_{m_{1}, 1}\right]^{S_{m_{1}}} \otimes \cdots \otimes \mathbb{C}\left[x_{1, n}, \ldots, x_{m_{n}, n}\right]^{S_{m_{n}}}, \quad \text { and } \\
R G_{\lambda}^{-} & =\mathbb{C}\left[x_{1,1}^{-1}, \ldots, x_{m_{1}, 1}^{-1}\right]^{S_{m_{1}}} \otimes \cdots \otimes \mathbb{C}\left[x_{1, n}^{-1}, \ldots, x_{m_{n}, n}^{-1}\right]^{S_{m_{n}}}
\end{aligned}
$$

Let

$$
\begin{aligned}
e_{+}^{(i)}(u) & =\left(1-x_{1, i} u\right)\left(1-x_{2, i} u\right) \cdots\left(1-x_{m_{i}, i} u\right) \quad \text { and } \\
e_{-}^{(i)}\left(u^{-1}\right) & =\left(1-x_{1, i}^{-1} u^{-1}\right)\left(1-x_{2, i}^{-1} u^{-1}\right) \cdots\left(1-x_{m_{i}, i}^{-1} u^{-1}\right) .
\end{aligned}
$$

Let $\mathbf{U}^{\prime}$ be the subalgebra of $\mathbf{U}$ without the generator $D$.
Theorem 3 (see [36, Sect.3.4]) The extremal weight module $L(\lambda)$ is the $\left(\mathbf{U}^{\prime} \otimes_{\mathbb{Z}}\right.$ $R G_{\lambda}$ )-module generated by a single vector $m_{\lambda}$ with relations

$$
\begin{gathered}
\mathbf{x}_{i, r}^{+} m_{\lambda}=0, \quad K_{i} m_{\lambda}=q^{m_{i}} m_{\lambda}, \quad C m_{\lambda}=m_{\lambda} \\
\mathbf{q}_{+}^{(i)}(u) m_{\lambda}=K_{i} \frac{e_{+}^{(i)}\left(q^{-1} u\right)}{e_{+}^{(i)}(q u)} m_{\lambda} \quad \text { and } \quad \mathbf{q}_{-}^{(i)}\left(u^{-1}\right) m_{\lambda}=K_{i} \frac{e_{-}^{(i)}\left(q u^{-1}\right)}{e_{-}^{(i)}\left(q^{-1} u^{-1}\right)} m_{\lambda}
\end{gathered}
$$

where $\mathbf{q}_{+}^{(i)}(u)$ and $\mathbf{q}_{-}^{(i)}\left(u^{-1}\right)$ are as defined in (19) and (24).
In this form $L(\lambda)$ has been termed the universal standard module [36, Sect. 3.4]) or the global Weyl module [9, Sect. 2]. See [36, Theorem 2] and [36, Remark 2.15]
for discussion of how to see that the extremal weight module, the universal standard module and the global Weyl module coincide.

Remark 2 Let

$$
0_{q}=\frac{1}{1-q}+\frac{q^{-1}}{1-q^{-1}}=\cdots+q^{-3}+q^{-2}+q^{-1}+1+q+q^{2}+\cdots
$$

(although $\frac{q^{-1}}{1-q^{-1}}=\frac{1}{q-1}=\frac{-1}{1-q}$, it is important to note that $0_{q}$ is not equal to 0 , it is a doubly infinite formal series in $q$ and $\left.q^{-1}\right)$. Since $\operatorname{deg}\left(e_{j}^{(i)}\right)=j$,

$$
\begin{aligned}
& \operatorname{gchar}\left(R G_{\lambda}^{+}\right)=\left(\prod_{i=1}^{n} \prod_{k=1}^{m_{i}} \frac{1}{1-q^{k}}\right) \quad \operatorname{gchar}\left(R G_{\lambda}^{-}\right)=\left(\prod_{i=1}^{n} \prod_{k=1}^{m_{i}} \frac{1}{1-q^{-k}}\right) \quad \text { and } \\
& \operatorname{gchar}\left(R G_{\lambda}\right)=\left(0_{q^{m_{1}}} \prod_{k=1}^{m_{1}-1} \frac{1}{1-q^{k}}\right)\left(0_{q^{m_{2}}} \prod_{k=1}^{m_{2}-1} \frac{1}{1-q^{k}}\right) \cdots\left(0_{q^{m_{n}}} \prod_{k=1}^{m_{n}-1} \frac{1}{1-q^{k}}\right)
\end{aligned}
$$

### 3.6 Level 0 L $(\lambda)$ and Finite Dimensional Simple U-Modules $L^{\text {fin }}(\boldsymbol{a}(u))$

The loop presentation provides a triangular decomposition of $\mathbf{U}$ (different form the usual triangular decomposition coming from the Kac-Moody presentation). The extremal weight module $L(\lambda)$ is the standard (Verma type) module for the loop triangular decomposition (see [8, Theorem 2.3(b)], [36, Lemma 2.14] and [7, outline of proof of Theorem 12.2.6]).

A Drinfeld polynomial is an $n$-tuple of polynomials $a(u)=\left(a^{(1)}(u), \ldots, a^{(n)}(u)\right)$ with $a^{(i)}(u) \in \mathbb{C}[u]$, represented as

$$
a(u)=a^{(1)}(u) \omega_{1}+\cdots+a^{(n)}(u) \omega_{n}, \quad \text { with } \quad a^{(i)}(u)=\left(u-a_{1, i}\right) \cdots\left(u-a_{m_{i}, i}\right)
$$

so that the coefficient of $u^{j}$ in $a^{(i)}(u)$ is $e_{m_{i}-j}^{(i)}\left(a_{1, i}, \ldots, a_{m_{i}, i}\right)$, the $\left(m_{i}-j\right)$ th elementary symmetric function evaluated at the values $a_{1, i}, \ldots, a_{m_{i}, i}$. The local Weyl module (a finite dimensional standard module) is defined by

$$
M^{\operatorname{fin}}(a(u))=L(\lambda) \otimes_{R G_{\lambda}} m_{a(u)}, \quad \text { where } e_{k}^{(i)}\left(x_{1, i}, x_{2, i}, \ldots\right) m_{a(u)}=e_{k}^{(i)}\left(a_{1, i}, \ldots, a_{m_{i}, i}\right) m_{a(u)}
$$

specifies the $R G_{\lambda}$-action on $m_{a(u)}$. In other words, the module $M^{\mathrm{fin}}(a(u))$ is $L(\lambda)$ except that variables $x_{j, i}$ from Sect. 3.5 have been specialised to the values $a_{j, i}$. As in Theorem 3, let $\mathbf{U}^{\prime}$ be the subalgebra of $\mathbf{U}$ without the generator $D$.

Theorem 4 (see [12, Theorem 2] and [8, Theorem 3.3]) The standard module $M^{\mathrm{fin}}(a(u))$ has a unique simple quotient $L^{\mathrm{fin}}(a(u))$ and

$$
\begin{array}{cc}
\begin{array}{c}
\{\text { Drinfeld polynomials }\} \\
a(u)=a^{(1)}(u) \omega_{1}+\cdots+a^{(n)}(u) \omega_{n}
\end{array} \longrightarrow & \left.\longrightarrow \text { finite dimensional simple } \mathbf{U}^{\prime} \text {-modules }\right\} \\
L^{\text {fin }}(a(u))
\end{array}
$$

is a bijection.

### 3.7 Path Models for the Crystals $B(\Lambda)$

The work of Littelmann [25, 26] provided a particularly convenient model for the crystals $B(\Lambda)$ when $\Lambda$ is positive or negative level. This model realizes the crystal as a set of paths $p: \mathbb{R}_{[0,1]} \rightarrow \mathfrak{h}^{*}$ with combinatorially defined Kashiwara operators $\tilde{e}_{0}, \ldots, \tilde{e}_{n}, \tilde{f}_{0}, \ldots, \tilde{f}_{n}$. In the LS (Lakshmibai-Seshadri) model the generator of the crystal $B(\Lambda)$ is the straight line path to $\Lambda$.

When $\lambda \in\left(\mathfrak{h}^{*}\right)_{\text {int }}^{0}$ using

$$
\begin{aligned}
p_{\lambda}: \mathbb{R}_{[0,1]} & \rightarrow \mathfrak{h}^{*} \\
t & \rightarrow t \lambda, \quad \text { the straight line path from } 0 \text { to } \lambda,
\end{aligned}
$$

as a generator for $B(\lambda)$ may not be the optimal choice. Remarkably, Naito and Sagaki (see [16, Definition 3.1.4 and Theorem 3.2.1] and [32, Theorem 4.6.1(b)]), have shown that $B(\lambda)$ can be constructed with sequences of Weyl group elements and rational numbers as in [25, Sect. 1.2, 1.3 and 2.2] but with the positive level length $\ell^{+}$and Bruhat order $\leq \pm$replaced by the level zero length $\ell^{0}$ and Bruhat order $\leq 0$. However, when working with the Naito-Sagaki construction one must be very careful not to identify the Naito-Sagaki sequences with actual paths (piecewise linear maps from $\mathbb{R}_{[0,1]}$ to $\mathfrak{h}^{*}$ ) because the natural map from Naito-Sagaki sequences to paths is not always injective (an example is provided by [20, Remark 5.10]).

If $\Lambda \in\left(\mathfrak{h}^{*}\right)_{\text {int }}^{+}$then $-\Lambda \in\left(\mathfrak{h}^{*}\right)_{\text {int }}^{-}$and

$$
\begin{aligned}
B(\Lambda) & =\left\{\tilde{f}_{i_{1}} \cdots \tilde{f}_{i_{k}} p_{\Lambda} \mid k \in \mathbb{Z}_{\geq 0} \text { and } i_{1}, \ldots, i_{k} \in\{0,1, \ldots, n\}\right\} \\
B(-\Lambda) & =\left\{\tilde{e}_{i_{1}} \cdots \tilde{e}_{i_{k}} p_{-\Lambda} \mid k \in \mathbb{Z}_{\geq 0} \text { and } i_{1}, \ldots, i_{k} \in\{0,1, \ldots, n\}\right\} .
\end{aligned}
$$

are each a single connected component and their characters are determined by the Weyl-Kac character formula [17, Theorem 11.13.3].

### 3.8 Crystals for Level 0 Extremal Weight Modules L( $\lambda$ )

For general $\lambda \in\left(\mathfrak{h}^{*}\right)_{\text {int }}^{0}$ the crystal $B(\lambda)$ is not connected (as a graph with edges determined by the Kashiwara operators $\tilde{e}_{0}, \ldots, \tilde{e}_{n}, \tilde{f}_{0}, \ldots \tilde{f}_{n}$ ). Let $\lambda=m_{1} \omega_{1}+\cdots+$ $m_{n} \omega_{n}$, with $m_{1}, \ldots, m_{n} \in \mathbb{Z}_{\geq 0}$. By [4, Corollary 4.15], the map

$$
\begin{array}{rlr}
\Phi_{\lambda}: L(\lambda) & \longrightarrow L\left(\omega_{1}\right)^{\otimes m_{1}} \otimes \cdots \otimes L\left(\omega_{n}\right)^{\otimes m_{n}} & \text { is injective }  \tag{27}\\
u_{\lambda} & \longmapsto u_{\omega_{1}}^{\otimes m_{1}} \otimes \cdots \otimes u_{\omega_{n}}^{\otimes m_{n}} &
\end{array}
$$

and gives rise to an injection of crystals

$$
B(\lambda) \hookrightarrow B\left(\omega_{1}\right)^{\otimes m_{1}} \otimes \cdots \otimes B\left(\omega_{n}\right)^{\otimes m_{n}}
$$

which takes the connected component of $B(\lambda)$ containing $b_{\lambda}$ to the connected component of $B\left(\omega_{1}\right)^{\otimes m_{1}} \otimes \cdots \otimes B\left(\omega_{n}\right)^{\otimes m_{n}}$ containing $b_{\omega_{1}}^{\otimes m_{1}} \otimes \cdots \otimes b_{\omega_{n}}^{\otimes m_{n}}$. Kashiwara [20, Theorem 5.15] fully described the structure of $L\left(\omega_{i}\right)$ (see [4, Theorem 2.16]). BeckNakajima analyzed the PBW basis of (22) and use (27) to show that the connected components of $B(\lambda)$ are labeled by $n$-tuples of partitions $\kappa=\left(\kappa^{(1)}, \ldots, \kappa^{(n)}\right)$ such that $\ell\left(\kappa^{(i)}\right)<m_{i}$. Together with a result of Fourier-Littelmann which shows that the crystal of the level 0 module $M^{\text {fin }}(a(u))$ is isomorphic to a level one Demazure crystal $B\left(\nu+\Lambda_{0}\right)_{\leq w}$, the full result is as detailed in Theorem 5 below.

A labeling set for a basis of $R G_{\lambda} /\left\langle e_{m_{i}}^{(i)}=1\right\rangle$ is

$$
S^{\lambda}=\left\{\vec{\kappa}=\left(\kappa^{(1)}, \ldots, \kappa^{(n)}\right) \mid \kappa^{(i)} \text { is a partition with } \ell\left(\kappa^{(i)}\right)<m_{i} \text { for } i \in\{1, \ldots, n\}\right\}
$$

The connected component of $b_{\lambda}$ in $B(\lambda)$ is

$$
B(\lambda)_{0}=\left\{\tilde{r}_{i_{1}} \cdots \tilde{r}_{i_{k}} b_{\lambda} \mid k \in \mathbb{Z}_{\geq 0} \text { and } \tilde{r}_{i_{1}}, \ldots, \tilde{r}_{i_{k}} \in\left\{\tilde{e}_{0}, \ldots, \tilde{e}_{n}, \tilde{f}_{0}, \ldots, \tilde{f}_{n}\right\}\right\}
$$

Define $B^{\text {fin }}(\lambda)$ to be the "crystal" which has a crystal graph which is the "quotient" of the crystal graph of $L(\lambda)$ obtained by identifying the vertices $b$ and $b^{\prime}$ if there is an element $\mathbf{s} \in R G_{\lambda}$ such that $\mathbf{s} G(b)=G\left(b^{\prime}\right)$, where $G(b)$ denotes the canonical basis element of $L(\lambda)$ corresponding to $b$.

$$
B^{\mathrm{fin}}(\lambda) \text { is the crystal of the finite dimensional standard module } M^{\mathrm{fin}}(a(u))
$$

Theorem 5 (see [4, Theorem 4.16], [36, Sect.3.4]) and [14, Proposition 3]) Let $\lambda=$ $m_{1} \omega_{1}+\cdots+m_{n} \omega_{n} \in\left(\mathfrak{h}^{*}\right)_{\mathrm{int}}^{0}$. As in Theorem 2 , let $\nu \in \mathfrak{a}_{\mathbb{Z}}^{*}, j \in \mathbb{Z}_{\geq 0}$ and $w \in W^{\text {ad }}$ such that $w\left(\nu+\Lambda_{0}\right)=-j \delta+\lambda+\Lambda_{0}$ and $w$ is minimal length. Then

$$
B(\lambda) \simeq B(\lambda)_{0} \times S^{\lambda} \quad \text { and } \quad B(\lambda)_{0} \simeq \mathbb{Z}^{k} \times B^{\text {fin }}(\lambda) \quad \text { and } \quad B^{\mathrm{fin}}(\lambda) \simeq B\left(\nu+\Lambda_{0}\right)_{\leq w}
$$

where $k$ is the number of elements of $m_{1}, \ldots, m_{n}$ which are nonzero.
Additional useful references for Theorem 5 are [36, Theorem 1] and [2, Theorem 1]. The first two statements in Theorem 5 are reflections of the very important fact that $L(\lambda)$ is free as an $R G_{\lambda}$-module. This fact that $L(\lambda)$ is free as an $R G_{\lambda}$-module was deduced geometrically, via the quiver variety, in [34, Theorem 7.3.5] (note that property $\left(T_{G_{\mathbf{w}} \times \mathbb{C}^{\times}}\right)$there includes the freeness, see the definition of property $\left(T_{G}\right)$ after $[34,(7.1 .1)])$. This freeness was understood more algebraically in the work of Fourier-Littelmann [14] and Chari-Ion [6, Cor. 2.10]. Further understanding of
the $R G_{\lambda}$-action in terms of the geometry of the semi-infinite flag variety is in [5, Sect. 5.1].

The last statement of Theorem 5 is proved by considering the map

$$
\begin{array}{ccc}
B(\lambda) & \longrightarrow B\left(\Lambda_{0}\right) \otimes B(\lambda) \\
\cup \cup & \text { Ul } & \text { given by } \\
B^{\mathrm{fin}}(\lambda) \xrightarrow{\sim} B\left(\nu+\Lambda_{0}\right)_{\leq w} & & b b_{\Lambda_{0}} \otimes b .
\end{array}
$$

where $b_{\Lambda_{0}}$ is the highest weight of the crystal $B\left(\Lambda_{0}\right)$. Combining the isomorphism $B^{\text {fin }}(\lambda) \simeq B\left(\nu+\Lambda_{0}\right)_{\leq w}$ with Theorem 2 and the positive level formula in Sect.3.4 gives

$$
\operatorname{char}(M(a(u)))=\operatorname{char}\left(L\left(\nu+\Lambda_{0}\right)_{\leq w}\right)=q^{-j} X^{-\Lambda_{0}} \tilde{E}_{\lambda}(q, 0) \quad \text { and }
$$

$$
\begin{equation*}
\operatorname{char}\left(L(\lambda)_{\leq w_{0}}\right)=\operatorname{gchar}\left(R G_{\lambda}^{+}\right) \tilde{E}_{w_{0} \lambda}(q, 0), \quad \operatorname{char}(L(\lambda))=\operatorname{gchar}\left(R G_{\lambda}\right) \tilde{E}_{w_{0} \lambda}(q, 0) \tag{28}
\end{equation*}
$$

## 4 Examples for $\mathfrak{g}=\widehat{\mathfrak{s l}}_{2}$

Let $\mathfrak{g}=\widehat{\mathfrak{s l}}_{2}$ with $\mathfrak{h}^{*}=\mathbb{C} \omega_{1} \oplus \mathbb{C} \Lambda_{0} \oplus \mathbb{C} \delta$ with affine Cartan matrix

$$
\binom{\alpha_{0}\left(h_{0}\right) \alpha_{0}\left(h_{1}\right)}{\alpha_{1}\left(h_{0}\right) \alpha_{1}\left(h_{1}\right)}=\left(\begin{array}{cc}
2 & -2 \\
-2 & 2
\end{array}\right) \quad \text { and } \quad \begin{aligned}
& \theta=\alpha_{1}=2 \omega_{1}, \theta^{\vee}=\alpha_{1}^{\vee}=h_{1} \\
& \Lambda_{1}=\omega_{1}+\Lambda_{0}, \alpha_{0}=-\alpha_{1}+\delta
\end{aligned}
$$

Using that $\left\langle\alpha_{1}, \alpha_{1}\right\rangle=2$, if $k \in \mathbb{Z}$ and $\mu^{\vee}=k \alpha_{1}^{\vee}=k \alpha_{1}=k 2 \omega_{1}=2 k \omega_{1}$ so that $\mu_{1}^{\vee}=2 k$ and $-\frac{1}{2}\left\langle\mu^{\vee}, \mu^{\vee}\right\rangle=-\frac{1}{2}(k 2 k)=-k^{2}$. Thus, following (8), (4) and (5), in the basis $\left\{\delta, \omega_{1}, \Lambda_{0}\right\}$ of $\mathfrak{h}^{*}$,

$$
t_{k \alpha^{\vee}}=\left(\begin{array}{ccc}
1 & -k & -k^{2} \\
0 & 1 & 2 k \\
0 & 0 & 1
\end{array}\right), \quad s_{1}=\left(\begin{array}{ccc}
1 & 0 & 0 \\
0 & -1 & 0 \\
0 & 0 & 1
\end{array}\right), \quad s_{0}=\left(\begin{array}{ccc}
1 & 1 & -1 \\
0 & -1 & 2 \\
0 & 0 & 1
\end{array}\right)
$$

These matrices are used to compute the $W$-orbits pictured in Sect. 1.1.

### 4.1 Macdonald Polynomials

The Demazure operators are given by

$$
D_{1} f=\frac{f-X^{-2 \omega_{1}}\left(s_{1} f\right)}{1-X^{-2 \omega_{1}}} \quad \text { and } \quad D_{0} f=\frac{f-X^{-\alpha_{0}}\left(s_{0} f\right)}{1-X^{-\alpha_{0}}}=\frac{f-q^{-1} X^{2 \omega_{1}}\left(s_{0} f\right)}{1-q^{-1} X^{2 \omega_{1}}} .
$$

The normalized Macdonald polynomials $\tilde{E}_{\omega_{1}}(q, t)$ and $\tilde{E}_{-\omega_{1}}(q, t)$ are

$$
\begin{aligned}
& \tilde{E}_{\omega_{1}}(q, t) \text { and } \tilde{E}_{-\omega_{1}}(q, t) \text { are } \\
& +\quad \leftarrow+ \\
& \tilde{E}_{\omega_{1}}(q, t)=X^{\omega_{1}} \quad \text { and } \quad \tilde{E}_{-\omega_{1}}(q, t)=X^{-\omega_{1}} \quad+\frac{(1-t)}{1-q t} X^{\omega_{1}} \\
& =X^{-\omega_{1}} \quad+\frac{\left(1-t^{-1}\right) q^{-1}}{1-q^{-1} t^{-1}} X^{\omega_{1}}
\end{aligned}
$$

giving $\tilde{E}_{\omega_{1}}(0, t)=\tilde{E}_{\omega_{1}}(\infty, t)=\tilde{E}_{\omega_{1}}(q, 0)=\tilde{E}_{\omega_{1}}(q, \infty)=X^{\omega_{1}}$ and

$$
\begin{array}{ll}
\tilde{E}_{-\omega_{1}}(0, t)=X^{-\omega_{1}}+(1-t) X^{\omega_{1}}, & \tilde{E}_{-\omega_{1}}(\infty, t)=X^{-\omega_{1}} \\
\tilde{E}_{-\omega_{1}}(q, 0)=X^{-\omega_{1}}+X^{\omega_{1}}, & \tilde{E}_{-\omega_{1}}(q, \infty)=X^{-\omega_{1}}+q^{-1} X^{\omega_{1}}
\end{array}
$$

The normalized Macdonald polynomials $\tilde{E}_{2 \omega_{1}}(q, t)$ and $\tilde{E}_{-2 \omega_{1}}(q, t)$ are


$$
\begin{aligned}
\tilde{E}_{2 \omega_{1}}(q, t) & =X^{2 \omega_{1}} \\
& =X^{2 \omega_{1}}
\end{aligned}
$$

$$
+\frac{(1-t) q}{1-q t}
$$

$$
+\frac{1-t^{-1}}{1 q^{-1} t^{-1}}
$$


$\tilde{E}_{-2 \omega_{1}}(q, t)=X^{-2 \omega_{1}}$


$$
=X^{-2 \omega_{1}} \quad+\frac{\left(1-t^{-1}\right) q^{-1}}{1-q^{-1} t^{-1}} \quad+\frac{\left(1-t^{-1}\right) q^{-2}}{1-q^{-2} t^{-1}} X^{2 \omega_{1}} \quad+\frac{\left(1-t^{-1}\right) q^{-2}}{\left(1-q^{-2} t^{-1}\right)} \frac{\left(1-t^{-1}\right)}{\left(1-q^{-1} t^{-1}\right)}
$$

giving

$$
\begin{aligned}
& \tilde{E}_{2 \omega_{1}}(0, t)=X^{2 \omega_{1}}, \quad \tilde{E}_{2 \omega_{1}}(\infty, t)=X^{2 \omega_{1}}+\left(1-t^{-1}\right) \\
& \tilde{E}_{2 \omega_{1}}(q, 0)=X^{2 \omega_{1}}+q, \tilde{E}_{2 \omega_{1}}(q, \infty)=X^{2 \omega_{1}}+1
\end{aligned}
$$

and

$$
\begin{array}{ll}
\tilde{E}_{-2 \omega_{1}}(0, t)=X^{-2 \omega_{1}}+(1-t) X^{2 \omega_{1}}+(1-t), & \tilde{E}_{-2 \omega_{1}}(\infty, t)=X^{-2 \omega_{1}}, \\
\tilde{E}_{-2 \omega_{1}}(q, 0)=X^{-2 \omega_{1}}+X^{2 \omega_{1}}+1+q, & \tilde{E}_{-2 \omega_{1}}(q, \infty)=X^{-2 \omega_{1}}+q^{-2} X^{2 \omega_{1}}+q^{-1}+q^{-2} .
\end{array}
$$

### 4.2 The Crystal B( $\left.\Lambda_{0}\right)$ and $B\left(\omega_{1}+\Lambda_{0}\right)$



Initial portion of the crystal graph of $B\left(\Lambda_{0}\right)$ for $\widehat{\mathfrak{s l}}_{2}$
The characters of the first few Demazure modules in $L\left(\Lambda_{0}\right)$ are

$$
\begin{aligned}
\operatorname{char}\left(L\left(\Lambda_{0}\right)_{\leq 1}\right) & =X^{\Lambda_{0}}=X^{\Lambda_{0}} \tilde{E}_{0}\left(q^{-1}, 0\right) \\
\operatorname{char}\left(L\left(\Lambda_{0}\right)_{\leq s_{0}}\right) & =D_{0} X^{\Lambda_{0}}=X^{\Lambda_{0}}\left(1+q X^{2 \omega_{1}}\right)=q X^{\Lambda_{0}}\left(X^{2 \omega_{1}}+q^{-1}\right)=q X^{\Lambda_{0}} \tilde{E}_{2 \omega_{1}}\left(q^{-1}, 0\right), \\
\operatorname{char}\left(L\left(\Lambda_{0}\right)_{\leq s_{1} s_{0}}\right) & =D_{1} D_{0} X^{\Lambda_{0}}=X^{\Lambda_{0}}\left(1+q X^{2 \omega_{1}}+q+q X^{-2 \omega_{1}}\right) \\
& =q X^{\Lambda_{0}}\left(X^{-2 \omega_{1}}+X^{2 \omega_{1}}+1+q^{-1}\right)=q X^{\Lambda_{0}} \tilde{E}_{-2 \omega_{1}}\left(q^{-1}, 0\right) .
\end{aligned}
$$

The crystal graph of $B\left(\omega_{1}+\Lambda_{0}\right)$ is pictured in Plate B and

$$
\begin{aligned}
\operatorname{char}\left(L\left(\omega_{1}+\Lambda_{0}\right)_{\leq 1}\right) & =X^{\Lambda_{0}} X^{\omega_{1}}=X^{\Lambda_{0}} \tilde{E}_{\omega_{1}}\left(q^{-1}, 0\right) \\
\operatorname{char}\left(L\left(\omega_{1}+\Lambda_{0}\right)_{\leq s_{1}}\right) & =D_{1} X^{\omega_{1}+\Lambda_{0}}=X^{\Lambda_{0}}\left(X^{\omega_{1}}+X^{-\omega_{1}}\right)=X^{\Lambda_{0}} \tilde{E}_{-\omega_{1}}\left(q^{-1}, 0\right) \\
\operatorname{char}\left(L\left(\omega_{1}+\Lambda_{0}\right)_{\leq s_{0} s_{1}}\right) & =D_{0} D_{1} X^{\omega_{1}+\Lambda_{0}}=X^{\Lambda_{0}}\left(X^{\omega_{1}}+X^{-\omega_{1}}+q X^{\omega_{1}}+q^{2} X^{3 \omega_{1}}\right) \\
& =q^{2} X^{\Lambda_{0}} \tilde{E}_{3 \omega_{1}}\left(q^{-1}, 0\right)
\end{aligned}
$$

### 4.3 The Crystal B( $\omega_{1}$ )

The crystal $B\left(\omega_{1}\right)=\left\{p_{-\omega_{1}+k \delta}, p_{\omega_{1}+k \delta} \mid k \in \mathbb{Z}\right\}$ is a single connected component. The crystal graph of $B\left(\omega_{1}\right)$ is pictured in Plate B. Following [17, (12.1.9)] and putting $q=$ $X^{-\delta}$ noting that $\tilde{E}_{-\omega_{1}}\left(q^{-1}, 0\right)=X^{\omega_{1}}+X^{-\omega_{1}}$ and $\tilde{E}_{-\omega_{1}}\left(q^{-1}, \infty\right)=X^{-\omega_{1}}+q X^{\omega_{1}}$, then

$$
\begin{aligned}
\operatorname{char}\left(L\left(\omega_{1}\right)_{\leq s_{1}}\right) & =\frac{1}{1-q^{-1}} \tilde{E}_{-\omega_{1}}\left(q^{-1}, 0\right) \quad \text { and } \\
\operatorname{char}\left(L\left(\omega_{1}\right)\right) & =0_{q} \tilde{E}_{-\omega_{1}}\left(q^{-1}, 0\right)=0_{q} \tilde{E}_{-\omega_{1}}\left(q^{-1}, \infty\right)
\end{aligned}
$$

where $0_{q}=\cdots+q^{-3}+q^{-2}+q^{-1}+1+q+q^{2}+\cdots$ as in Remark 2 .

### 4.4 The Crystal B(2 $\left.\omega_{1}\right)$

On crystals, the injective $\mathbf{U}$-module homomorphism $L\left(2 \omega_{1}\right) \hookrightarrow L\left(\omega_{1}\right) \otimes L\left(\omega_{1}\right)$ given in (27) is the inclusion

$$
\begin{aligned}
& B\left(\omega_{1}\right) \otimes B\left(\omega_{1}\right)=\left\{p_{w_{1} \omega_{1}} \otimes p_{w_{2} \omega_{1}} \mid w_{1}, w_{2} \in W^{\text {ad }}\right\} \\
& \cup \mid \\
& B\left(2 \omega_{1}\right)=\left\{p_{w_{1} \omega_{1}} \otimes p_{w_{2} \omega_{1}} \mid w_{1}, w_{2} \in W^{\text {ad }} \text { with } w_{1} \unrhd w_{2}\right\}
\end{aligned}
$$

The connected components of $B\left(2 \omega_{1}\right)$ are determined by

$$
B\left(2 \omega_{1}\right)=B\left(2 \omega_{1}\right)_{0} \times S^{2 \omega_{1}}, \quad \text { where } \quad S^{2 \omega_{1}}=\{\text { partitions } \kappa \text { with } \ell(\kappa)<2\}=\mathbb{Z}_{\geq 0}
$$

For $\kappa \in \mathbb{Z}_{\geq 0}$, the connected component corresponding to $\kappa$, as a subset of $B\left(\omega_{1}\right) \otimes$ $B\left(\omega_{1}\right)$, is

$$
B\left(2 \omega_{1}\right)_{\kappa}=\left\{p_{w_{1} \omega_{1}} \otimes p_{w_{2} \omega_{1}} \mid w_{1} \xrightarrow{0} w_{2} \text { and } \ell^{0}\left(w_{1}\right)-\ell^{0}\left(w_{2}\right) \in\{\kappa, \kappa+1\}\right\}
$$

Representatives of the components and the crystal graph of the connected component $B\left(2 \omega_{1}\right)_{0}$ are pictured in Plate C. Inspection of the crystal graphs gives

$$
\begin{aligned}
\operatorname{char}\left(L\left(2 \omega_{1}\right)_{\leq s_{1}}\right) & =\frac{1}{1-q^{-2}} \frac{1}{1-q^{-1}} \tilde{E}_{-2 \omega_{1}}\left(q^{-1}, 0\right), \quad \text { and } \\
\operatorname{char}\left(L\left(2 \omega_{1}\right)\right) & =0_{q} \frac{1}{1-q} \tilde{E}_{-2 \omega_{1}}\left(q^{-1}, 0\right)=0_{q} \frac{1}{1-q} \tilde{E}_{-2 \omega_{1}}\left(q^{-1}, \infty\right)
\end{aligned}
$$

## 5 Examples for $\mathfrak{g}=\widehat{\mathfrak{s l}}_{3}$

Let $\mathfrak{g}=\widehat{\mathfrak{s l}}_{3}$. The affine Dynkin diagram is $\bigcirc$ and the affine Cartan matrix is

$$
\left(\begin{array}{lll}
\alpha_{0}\left(h_{0}\right) & \alpha_{0}\left(h_{1}\right) & \alpha_{0}\left(h_{2}\right) \\
\alpha_{1}\left(h_{0}\right) & \alpha_{1}\left(h_{1}\right) & \alpha_{1}\left(h_{2}\right) \\
\alpha_{2}\left(h_{0}\right) & \alpha_{2}\left(h_{1}\right) & \alpha_{2}\left(h_{2}\right)
\end{array}\right)=\left(\begin{array}{ccc}
2 & -1 & -1 \\
-1 & 2 & -1 \\
-1 & -1 & 2
\end{array}\right)
$$

and $h_{0}=-\left(h_{1}+h_{2}\right)+K, \quad \alpha_{0}=-\left(\alpha_{1}+\alpha_{2}\right)+\delta, \quad \Lambda_{1}=\omega_{1}+\Lambda_{0}, \quad \Lambda_{2}=\omega_{2}$ $+\Lambda_{0}$.

In the basis $\left\{\delta, \omega_{1}, \omega_{2}, \Lambda_{0}\right\}$ of $\mathfrak{h}^{*}$ the action of the affine Weyl group $W^{\text {ad }}$ is given by

$$
\begin{gathered}
s_{1}=\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & -1 & 0 & 0 \\
0 & 1 & 1 & 0 \\
0 & 0 & 0 & 1
\end{array}\right) \quad s_{2}=\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & 1 & 1 & 0 \\
0 & 0 & -1 & 0 \\
0 & 0 & 0 & 1
\end{array}\right) \quad s_{0}=\left(\begin{array}{cccc}
1 & 1 & 1 & -1 \\
0 & 0 & -1 & 1 \\
0 & -1 & 0 & 1 \\
0 & 0 & 0 & 1
\end{array}\right), \quad \text { and } \\
t_{k_{1} h_{1}+k_{2} h_{2}}=\left(\begin{array}{cccc}
1 & -k_{1}-k_{2}-k_{1}^{2}-k_{2}^{2}+k_{1} k_{2} \\
0 & 1 & 0 & 2 k_{1}-k_{2} \\
0 & 0 & 1 & 2 k_{2}-k_{1} \\
0 & 0 & 0 & 1
\end{array}\right) \quad \text { for } k_{1}, k_{2} \in \mathbb{Z} .
\end{gathered}
$$

These matrices are used to compute the orbits pictured at the end of Sect. 1.1.

### 5.1 The Extremal Weight Modules $L\left(\omega_{1}\right)$ and $L\left(\omega_{2}\right)$

Letting $\mathbb{C}^{3}=\mathbb{C}$-span $\left\{v_{1}, v_{2}, v_{3}\right\}$ be the standard representation of $\mathfrak{g}=\mathfrak{s l}_{3}$, the extremal weight modules $L\left(\omega_{1}\right)$ and $L\left(\omega_{2}\right)$ for $\mathfrak{g}=\widehat{\mathfrak{s l}}_{3}$ are

$$
L\left(\omega_{1}\right)=\mathbb{C}^{3} \otimes_{\mathbb{C}} \mathbb{C}\left[\epsilon, \epsilon^{-1}\right] \quad \text { and } \quad L\left(\omega_{2}\right)=\left(\Lambda^{2} \mathbb{C}^{3}\right) \otimes_{\mathbb{C}} \mathbb{C}\left[\epsilon, \epsilon^{-1}\right]
$$

with $u_{\omega_{1}}=v_{1}$ and $u_{\omega_{2}}=v_{1} \wedge v_{2}$, respectively. The crystals $B\left(\omega_{1}\right)$ and $B\left(\omega_{2}\right)$ have realizations as sets of straight line paths:

$$
B\left(\omega_{1}\right)=\left\{p_{\omega_{1}+k \delta}, p_{s_{1} \omega_{1}+k \delta}, p_{-\omega_{2}+k \delta} \mid k \in \mathbb{Z}\right\}, \quad B\left(\omega_{2}\right) \quad=\left\{p_{\omega_{2}+k \delta}, p_{s_{2} \omega_{2}+k \delta}, p_{-\omega_{1}+k \delta} \mid k \in \mathbb{Z}\right\} .
$$


$B\left(\omega_{1}\right)$ crystal graph

$B\left(\omega_{2}\right)$ crystal graph

Each of the crystals $B\left(\omega_{1}\right)$ and $B\left(\omega_{2}\right)$ has a single connected component, all weight spaces are one dimensional and

$$
\begin{aligned}
\operatorname{char}\left(L\left(\omega_{1}\right)_{\leq s_{2} s_{1}}\right) & =\frac{1}{1-q^{-1}}\left(X^{\omega_{1}}+X^{s_{1} \omega_{1}}+X^{-\omega_{2}}\right)=\frac{1}{1-q^{-1}} \tilde{E}_{-\omega_{2}}\left(q^{-1}, 0\right), \quad \text { and } \\
\operatorname{char}\left(L\left(\omega_{1}\right)\right) & =0_{q}\left(X^{\omega_{1}}+X^{-\omega_{2}}+X^{s_{1} \omega_{1}}\right)=0_{q} \tilde{E}_{-\omega_{2}}\left(q^{-1}, 0\right) \\
& =0_{q}\left(q^{-1} X^{\omega_{1}}+q^{-1} X^{\omega_{2}-\omega_{1}}+X^{-\omega_{2}}\right)=0_{q} \tilde{E}_{-\omega_{2}}\left(q^{-1}, \infty\right)
\end{aligned}
$$

where $0_{q}=\cdots+q^{-3}+q^{-2}+q^{-1}+1+q+q^{2}+\cdots$ as in Remark 2 .
For $a \in \mathbb{C}$, the crystals of $M^{\text {fin }}\left((u-a) \omega_{1}\right) \cong \mathbb{C}^{3}$ and $M^{\text {fin }}\left((u-a) \omega_{2}\right) \cong \Lambda^{2}\left(\mathbb{C}^{3}\right)$ have crystal graphs $B^{\text {fin }}\left(\omega_{1}\right)$ and $B^{\text {fin }}\left(\omega_{2}\right)$.

$B^{\mathrm{fin}}\left(\omega_{1}\right)$

$B^{\mathrm{fin}}\left(\omega_{2}\right)$

### 5.2 The Extremal Weight Module $L\left(\omega_{1}+\omega_{2}\right)$

To construct the crystal $B\left(\omega_{1}+\omega_{2}\right)$ use

$$
B\left(\omega_{1}\right) \otimes B\left(\omega_{2}\right)=\left\{p_{v \omega_{1}+k \delta} \otimes p_{w \omega_{2}+\ell \delta} \mid v \in\left\{1, s_{1}, s_{2} s_{1}\right\}, w \in\left\{1, s_{2}, s_{1} s_{2}\right\}, k \in \mathbb{Z}, \ell \in \mathbb{Z}\right\}
$$

with the tensor product action for crystals given by (see, for example, [39, Prop. 5.7])

$$
\tilde{f}_{i}\left(p_{1} \otimes p_{2}\right)=\left\{\begin{array}{ll}
\tilde{f}_{i} p_{1} \otimes p_{2}, & \text { if } d_{i}^{+}\left(p_{1}\right)>d_{i}^{-}\left(p_{2}\right), \\
p_{1} \otimes \tilde{f}_{i} p_{2}, & \text { if } d_{i}^{+}\left(p_{1}\right) \leq d_{i}^{-}\left(p_{2}\right),
\end{array} \quad \tilde{e}_{i}\left(p_{1} \otimes p_{2}\right)= \begin{cases}\tilde{e}_{i} p_{1} \otimes p_{2}, & \text { if } d_{i}^{+}\left(p_{1}\right) \geq d_{i}^{-}\left(p_{2}\right), \\
p_{1} \otimes \tilde{e}_{i} p_{2}, & \text { if } d_{i}^{+}\left(p_{1}\right)<d_{i}^{-}\left(p_{2}\right) .\end{cases}\right.
$$

where $d_{i}^{ \pm}(p)$ are determined by

$$
\tilde{f}_{i}^{d_{i}^{+}(p)} p \neq 0 \text { and } \tilde{f}_{i}^{d_{i}^{+}(p)+1} p=0, \quad \text { and } \quad \tilde{e}_{i}^{d_{i}^{-}(p)} p \neq 0 \text { and } \tilde{e}_{i}^{d_{i}^{-}(p)+1} p=0
$$

The crystal $B\left(\omega_{1}+\omega_{2}\right)$ is realized as a subset of $B\left(\omega_{1}\right) \otimes B\left(\omega_{2}\right)$ via the crystal embedding

$$
\begin{aligned}
B\left(\omega_{1}+\omega_{2}\right) & \hookrightarrow B\left(\omega_{1}\right) \otimes B\left(\omega_{2}\right) \\
b_{\omega_{1}+\omega_{2}} & \longmapsto p_{\omega_{1}} \otimes p_{\omega_{2}}
\end{aligned}
$$

By Theorem 5, $B\left(\omega_{1}+\omega_{2}\right)$ is connected and is generated by $b_{\omega_{1}+\omega_{2}}$. The crystal $B\left(\omega_{1}+\omega_{2}\right)$ is pictured and its character is computed in Plate D .

## 6 Alcove Walks for Affine Flag Varieties

### 6.1 The Affine Kac-Moody Group G

The most visible form of the loop group is $G=\stackrel{\circ}{G}(\mathbb{C}((\epsilon)))$, where $\mathbb{C}((\epsilon))$ is the field of formal power series in a variable $\epsilon$ and $\stackrel{\circ}{G}$ is a reductive algebraic group. The favourite example is when

$$
\text { when } \quad \stackrel{\circ}{G}=G L_{n} \quad \text { and the loop group is } \quad G=G L_{n}(\mathbb{C}((\epsilon))) \text {, }
$$

the group of $n \times n$ invertible matrices with entries in $\mathbb{C}((\epsilon))$. A slightly more extended, and extremely powerful, point of view is to let $G$ be the Kac-Moody group whose Lie algebra is the affine Lie algebra $\mathfrak{g}$ of Sect. 2.1, so that $G$ is a central extension of a semidirect product (where the semidirect product comes from the action of $\mathbb{C}^{\times}$on $\dot{G}(\mathbb{C}((\epsilon))$ by "loop rotations")

$$
\begin{equation*}
\{1\} \rightarrow \mathbb{C}^{\times} \rightarrow G \rightarrow \dot{G}(\mathbb{C}((\epsilon))) \rtimes \mathbb{C}^{\times} \rightarrow\{1\} \quad \text { so that } \quad G=\exp (\mathfrak{g}) \tag{29}
\end{equation*}
$$

In this section we wish to work with $G$ via generators and relations. Fortunately, presentations of $G$ are well established in the work of Steinberg [41] and Tits [42] and others (see [38, Sect. 3] for a survey). More specifically, up to an extra (commutative) torus $T \cong\left(\mathbb{C}^{\times}\right)^{k}$ for some $k$, the group $G$ is generated by subgroups isomorphic to $S L_{2}(\mathbb{C})$, the images of homomorphisms

$$
\begin{array}{rlll}
\varphi_{i}: & \longrightarrow G & \\
\left(\begin{array}{cc}
S L_{2}(\mathbb{C}) & \\
1 & c \\
0 & 1
\end{array}\right) & \longmapsto & x_{\alpha_{i}}(c) & \\
\left(\begin{array}{cc}
1 & 0 \\
c & 1
\end{array}\right) & \longmapsto & x_{-\alpha_{i}}(c) & \\
\text { for each vertex } & \quad \\
\left(\begin{array}{cc}
d & 0 \\
0 & d^{-1}
\end{array}\right) & \longmapsto & h_{\alpha_{i}^{\vee}}(d) &  \tag{30}\\
\text { of the affine Dynkin diagram. } \\
\left(\begin{array}{cc}
c & 1 \\
-1 & 0
\end{array}\right) & \longmapsto & y_{i}(c)
\end{array}
$$

For each $\alpha \in \stackrel{\circ}{R}^{+}$there is a homomorphism $\varphi_{\alpha}: S L_{2}(\mathbb{C}((\epsilon))) \rightarrow G$ and we let

$$
x_{\alpha}(f)=\varphi_{\alpha}\left(\begin{array}{ll}
1 & f  \tag{31}\\
0 & 1
\end{array}\right) \quad \text { and } \quad x_{-\alpha}(f)=\varphi_{\alpha}\left(\begin{array}{ll}
1 & 0 \\
f & 1
\end{array}\right) \quad \text { for } f \in \mathbb{C}((\epsilon)) .
$$

For each $z \in W^{\text {ad }}$ fix a reduced word $z=s_{j_{1}} \cdots s_{j_{r}}$, define
$n_{z}=y_{j_{1}}(0) \cdots y_{j_{r}}(0)$ and define $\quad x_{ \pm \alpha+r \delta}(c)=x_{ \pm \alpha}\left(c \epsilon^{r}\right)$ for $\alpha \in \AA^{+}$and $r \in \mathbb{Z}$.
In the following sections, for simplicity, we will work only with the loop group $G=\stackrel{\circ}{G}(\mathbb{C}((\epsilon)))$ rather than the full affine Kac-Moody group as in (29). We shall also use a slightly more general setting where $\mathbb{C}$ is replaced by an arbitrary field $\mathbb{k}$ so that, in Proposition 4 , we may let $\mathbb{k}=\mathbb{F}_{q}$ be the finite field with $q$ elements.

### 6.2 The Affine Flag Varieties $G / I^{+}, G / I^{0}$ and $G / I^{-}$

Let $\mathbb{k}$ be a field. Define subgroups of $\stackrel{\circ}{G}(\mathbb{k})$ by

$$
\begin{aligned}
\text { "unipotent upper triangular matrices" } & U^{+}(\mathbb{k})=\left\langle x_{\alpha}(c) \mid \alpha \in \stackrel{\circ}{R}^{+}, c \in \mathbb{k}\right\rangle, \\
\text { "diagonal matrices" } & H\left(\mathbb{k}^{\times}\right)=\left\langle h_{\alpha_{i}^{\vee}}(c) \mid i \in\{1, \ldots, n\}, c \in \mathbb{k}^{\times}\right\rangle, \\
\text {"unipotent lower triangular matrices" } & \left.\left.U^{-}(\mathbb{k})=\left\langle x_{-\alpha}(c)\right| \alpha \in \stackrel{\circ}{R}^{+}, c \in \mathbb{k}\right\}\right\rangle .
\end{aligned}
$$

Let
$\mathbb{F}=\mathbb{k}((\epsilon))$, which has $\mathfrak{o}=\mathbb{k}[[\epsilon]]$ and $\mathfrak{o}^{\times}=\{p \in \mathbb{k}[[\epsilon]] \mid p(0) \neq 0\}$, or let $\mathbb{F}=\mathbb{k}\left[\epsilon, \epsilon^{-1}\right], \quad$ which has $\mathfrak{o}=\mathbb{k}[\epsilon] \quad$ and $\quad \mathfrak{o}^{\times}=\mathbb{k}^{\times}$.

Define subgroups of $G=\stackrel{\circ}{G}(\mathbb{F})$ by

$$
U^{-}(\mathbb{F})=\left\langle x_{-\alpha+k \delta}(c) \mid c \in \mathbb{k}, \alpha \in \stackrel{R}{R}^{+}, k \in \mathbb{Z}\right\rangle \quad \text { and } \quad H\left(\mathfrak{o}^{\times}\right)=\left\langle h_{\alpha_{i}^{\vee}}(c) \mid c \in \mathfrak{o}^{\times}\right\rangle
$$

Let $g(0)$ denote $g$ evaluated at $\epsilon=0$ and let $g(\infty)$ denote $g$ evaluated at $\epsilon^{-1}=$ 0 . Following, for example, [37, Lect. 11 Theorem (B)] and [27, Sect. 11], define subgroups of $G(\mathbb{F})$ by

$$
\begin{array}{rlrl}
\text { (positive Iwahori) } & & I^{+} & =\left\{g \in G \mid g(0) \text { exists and } g(0) \in U^{+}(\mathbb{k}) H\left(\mathbb{k}^{\times}\right)\right\}, \\
\text {(level 0 Iwahori) } & I^{0} & =U^{-}(\mathbb{F}) H\left(\mathfrak{o}^{\times}\right) \\
\text {(negative Iwahori) } & I^{-} & =\left\{g \in G \mid g(\infty) \text { exists and } g(\infty) \in U^{-}(\mathbb{k}) H\left(\mathbb{k}^{\times}\right)\right\},
\end{array}
$$

Then

$$
G / I^{+} \text {is the positive level (thin) affine flag variety, }
$$

$G / I^{0}$ is the level 0 (semi-infinite) affine flag variety,
$G / I^{-}$is the negative level (thick) affine flag variety.
These are studied with the aid of the decompositions (relax notation and write $z I^{+}$ for $n_{z} I^{+}$)

$$
G=\bigsqcup_{x \in W^{\mathrm{ad}}} I^{+} x I^{+}, \quad G=\bigsqcup_{y \in W^{\mathrm{ad}}} I^{0} y I^{+}, \quad G=\bigsqcup_{z \in W^{\mathrm{ad}}} I^{-} z I^{+}
$$

### 6.3 Labeling Points of $I^{+} w I^{+}, I^{+} w I^{+} \cap I^{0} v I^{+}$and $I^{+} w I^{+} \cap I^{-} z I^{+}$

Recall the bijection between elements of $W^{\text {ad }}$ and alcoves given in (6) and note that if $v \in W^{\text {ad }}$ and $i \in\{0, \ldots, n\}$ then the hyperplane separating the alcoves corresponding to $v$ and $s_{i} v$ is

$$
\mathfrak{h}^{v \alpha_{i}^{\vee}}=\left\{x+\Lambda_{0} \mid x \in \mathfrak{a}_{\mathbb{R}}^{*} \text { and }\left\langle x+\Lambda_{0}, v \alpha_{i}^{\vee}\right\rangle=0\right\}
$$

A blue labeled step of type $i$ is

$$
\xrightarrow{\mathfrak{h}^{v \alpha_{i}^{\vee}}} \underset{c}{v s_{i}} \quad \text { with } \quad v s_{i} \Perp v . \quad \text { Let } \quad \Phi^{+}\left(\begin{array}{l|l}
v^{h^{v \alpha_{i}^{\vee}}} \\
v & v s_{i} \\
\hline & c
\end{array}\right)=y_{i}(c),
$$

where $y_{i}(c)$ is as in (30). A blue labeled walk of type $\vec{w}=\left(i_{1}, \ldots, i_{\ell}\right)$ is a sequence $\left(p_{1}, \ldots, p_{\ell}\right)$ where $p_{k}$ is a blue labeled step of type $i_{k}$ and which begins at the alcove 1.

A red labeled step of type $i$ is


With notations as in (32), let

A red labeled walk of type $\vec{w}=\left(i_{1}, \ldots, i_{\ell}\right)$ is a sequence $\left(p_{1}, \ldots, p_{\ell}\right)$ where $p_{k}$ is a red labeled step of type $i_{k}$ and which begins at the alcove 1.
A green labeled step of type $i$ is


Let

A green labeled walk of type $\vec{w}=\left(i_{1}, \ldots, i_{\ell}\right)$ is a sequence $\left(p_{1}, \ldots, p_{\ell}\right)$ where $p_{k}$ is a green labeled step of type $i_{k}$ and which begins at the alcove 1.

Theorem 6 Let $v, w \in W^{\text {ad }}$ and fix a reduced expression $w=s_{i_{1}} \ldots s_{i_{\ell}}$ for $w$. The maps $\Psi_{\vec{w}}^{+}, \Psi_{\vec{w}, v}^{0}, \Psi_{\vec{w}, v}^{-}$are bijections.

$$
\left.\begin{array}{c}
\Phi_{\vec{w}}^{+}:\left\{\begin{array}{ccc}
\left.\begin{array}{l}
\text { blue labeled paths of type } \\
\vec{w}=\left(i_{1}, \ldots, i_{\ell}\right)
\end{array}\right\} & \xrightarrow{\sim} & \left(I^{+} w I^{+}\right) / I^{+} \\
\left(p_{1}, \ldots, p_{\ell}\right)
\end{array}\right. \\
\\
\Phi_{\vec{w}, v}^{0}:\left\{\begin{array}{l}
\text { red labeled paths of type } \\
\vec{w}=\left(i_{1}, \ldots, i_{\ell}\right) \text { ending in } v
\end{array}\right\} \\
\left(p_{1}, \ldots, p_{\ell}\right)
\end{array}\right) \stackrel{\sim}{\longrightarrow}\left(I^{0} v I^{+} \cap I^{+} w I^{+}\right) / I^{+} .
$$

The proof of Theorem 6 is by induction on the length of $w$ following [38, Theorem 4.1 and $\S 7$ ] where (a) and (b) are proved. The induction step can be formulated as the following proposition.

Proposition 2 Let $v, w \in W^{\text {ad }}$ and fix a reduced expression $\vec{w}=s_{i_{1}} \ldots s_{i_{\ell}}$ for $w$. Let $j \in\{0, \ldots, n\}$ and $c \in \mathbb{C}$. If $w s_{j}<+w$ then assume that the reduced word for $w$ is chosen with $i_{\ell}=j$. Let
$\tilde{c} \in \mathbb{C}$ and $\tilde{b}_{1} \in I^{+}$be the unique elements such that $b_{1} y_{j}(c)=y_{j}(\tilde{c}) \tilde{b}_{1}$.
(a) Let $p=\left(p_{1}, \ldots, p_{\ell}\right)$ be a blue labeled path of type $\left(i_{1}, \ldots, i_{\ell}\right)$ and let $\Phi_{\vec{w}}^{+}(p)=$ $y_{i_{1}}\left(c_{1}\right) \ldots y_{i_{\ell}}\left(c_{\ell}\right)$. Then

$$
\begin{aligned}
& \left(y_{i_{1}}\left(c_{1}\right) \cdots y_{i_{\ell}}\left(c_{\ell}\right) b_{1}\right)\left(y_{j}(c) b_{2}\right) \\
& = \begin{cases}y_{i_{1}}\left(c_{1}\right) \cdots y_{i_{\ell}}\left(c_{\ell}\right) y_{j}(\tilde{c}) \tilde{b}_{1} b_{2}, & \text { if } w \ll w s_{j}, \\
y_{i_{1}}\left(c_{1}\right) \cdots y_{i_{\ell-1}}\left(c_{\ell-1}\right) y_{i_{\ell}}\left(c_{\ell}-\tilde{c}^{-1}\right) h_{h_{1}^{v}}(\tilde{c}) x_{\alpha_{i_{\ell}}}\left(-\tilde{c}^{-1}\right) \tilde{b}_{1} b_{2}, & \text { if } w s_{j} \ll w \text { and } \tilde{c} \neq 0, \\
y_{i_{1}}\left(c_{1}\right) \cdots y_{i_{\ell-1}}\left(c_{\ell-1}\right) h_{\alpha_{i_{\ell}}}(-1) x_{\alpha_{i_{\ell}}}(c) \tilde{b}_{1} b_{2}, & \text { if } w s_{j}<w \text { and } \tilde{c}=0,\end{cases}
\end{aligned}
$$

(b) Let $p=\left(p_{1}, \ldots, p_{\ell}\right)$ be a red labeled path of type $\left(i_{1}, \ldots, i_{\ell}\right)$ ending in $v$ and let $\Phi_{\vec{w}, v}^{0}(p)=x_{\gamma_{1}}\left(c_{1}\right) \cdots x_{\gamma_{\ell}}\left(c_{\ell}\right) n_{v} I^{+}$where the notation is as in (32). Then

$$
\begin{aligned}
& \left(x_{\gamma_{1}}\left(c_{1}\right) \cdots x_{\gamma_{\ell}}\left(c_{\ell}\right) n_{v} b_{1}\right)\left(y_{j}(c) b_{2}\right) \\
& = \begin{cases}x_{\gamma_{1}}\left(c_{1}\right) \cdots x_{\gamma_{\ell}}\left(c_{\ell}\right) x_{v \alpha_{j}}( \pm \tilde{c}) n_{v s_{j}} \tilde{b}_{1} b_{2} . & \text { if } w<0 w s_{j}, \\
\left.x_{\gamma_{1}}\left(c_{1}\right) \cdots x_{\gamma_{\ell}}\left(c_{\ell}\right) x_{-v \alpha_{j}}\left(\tilde{c}^{-1}\right)\right)_{v} x_{\alpha_{j}}(-\tilde{c}) h_{\alpha_{j}^{\vee}}(\tilde{c}) \tilde{b}_{1} b_{2}, & \text { if } w s_{j}<0 w \text { and } \tilde{c} \neq 0, \\
x_{\gamma_{1}}\left(c_{1}\right) \cdots x_{\gamma_{\ell}}\left(c_{\ell}\right) x_{-v \alpha_{j}}(0) n_{v s_{j}} \tilde{b}_{1} b_{2}, & \text { if } w s_{j}<0 w \text { and } \tilde{c}=0,\end{cases}
\end{aligned}
$$

(c) Let $p=\left(p_{1}, \ldots, p_{\ell}\right)$ be a green labeled path of type ( $i_{1}, \ldots, i_{\ell}$ ) ending in $v$ and let $\Phi_{\vec{w}, v}^{-}(p)=x_{\gamma_{1}}\left(c_{1}\right) \cdots x_{\gamma_{\ell}}\left(c_{\ell}\right) n_{v} I^{+}$where the notation is as in (32). Then

$$
\begin{aligned}
& \left(x_{\gamma_{1}}\left(c_{1}\right) \cdots x_{\gamma_{\ell}}\left(c_{\ell}\right) n_{v} b_{1}\right)\left(y_{j}(c) b_{2}\right) \\
& = \begin{cases}x_{\gamma_{1}}\left(c_{1}\right) \cdots x_{\gamma_{\ell}}\left(c_{\ell}\right) x_{v \alpha_{j}}( \pm \tilde{c}) n_{v s_{j}} \tilde{b}_{1} b_{2} . & \text { if } w<w s_{j}, \\
x_{\gamma_{1}}\left(c_{1}\right) \cdots x_{\gamma_{\ell}}\left(c_{\ell}\right) x_{-v \alpha_{j}}\left(\tilde{c}^{-1}\right) n_{v} x_{\alpha_{j}}(-\tilde{c}) h_{\alpha_{j}^{\vee}}(\tilde{c}) \tilde{b}_{1} b_{2}, & \text { if } w s_{j}<w \text { and } \tilde{c} \neq 0, \\
x_{\gamma_{1}}\left(c_{1}\right) \cdots x_{\gamma_{\ell}}\left(c_{\ell}\right) x_{-v \alpha_{j}}(0) n_{v s_{j}} \tilde{b}_{1} b_{2}, & \text { if } w s_{j}<w \text { and } \tilde{c}=0,\end{cases}
\end{aligned}
$$

### 6.4 Actions of the Affine Hecke Algebra

The affine flag representation is

$$
C\left(G / I^{+}\right)=\mathbb{C}-\operatorname{span}\left\{y_{\vec{w}}(\vec{c}) I^{+} \mid w \in W^{\text {ad }}, \vec{c}=\left(c_{1}, \ldots, c_{\ell}\right) \in \mathbb{C}^{\ell(w)}\right\}
$$

where, for a fixed reduced word $w=s_{i_{1}} \cdots s_{i_{\ell}}$

$$
y_{\vec{w}}(\vec{c})=y_{i_{1}}\left(c_{1}\right) \cdots y_{i_{\ell}}\left(c_{\ell}\right) I^{+}=\sum_{g \in y_{i_{1}}\left(c_{1}\right) \cdots y_{i_{\ell}}\left(c_{\ell}\right) I^{+}} g \quad \text { (a formal sum) }
$$

Let

$$
\begin{aligned}
& C\left(I^{+} \backslash G / I^{+}\right)=\mathbb{C}-\operatorname{span}\left\{T_{w} \mid w \in W^{\text {ad }}\right\}, \text { where } T_{w}=I^{+} w I^{+}, \\
& C\left(I^{0} \backslash G / I^{+}\right)=\mathbb{C}-\operatorname{span}\left\{X^{w} \mid w \in W^{\text {ad }}\right\}, \text { where } X^{w}=I^{0} w I^{+}, \\
& C\left(I^{-} \backslash G / I^{+}\right)=\mathbb{C}-\operatorname{span}\left\{L^{w} \mid w \in W^{\text {ad }}\right\}, \text { where } L^{w}=I^{-} w I^{+},
\end{aligned}
$$

The affine Hecke algebra is $h_{I^{+}}\left(G / I^{+}\right)=\mathbb{C}$-span $\left\{T_{w} \mid w \in W\right\}$,

$$
\text { where } T_{w}=I^{+} n_{w} I^{+}=\sum_{x \in I^{+} n_{w} I^{+}} x \quad \text { (a formal sum). }
$$

The formal sums allow us to view all of these elements as elements of the group algebra of $G$, where infinite formal sums are allowed (to do this precisely one should use Haar measure and a convolution product). Proposition 3 computes the (right) action of $h_{I^{+}}\left(G / I^{+}\right)$on $C\left(G / I^{+}\right)$, and Proposition 4 computes the (right) action of $h_{I^{+}}\left(G / I^{+}\right)$on $C\left(I^{+} \backslash G / I^{+}\right)$on $C\left(I^{0} \backslash G / I^{+}\right)$, and on $C\left(I^{-} \backslash G / I^{+}\right)$. Proposition 3 follows from Proposition 2(a) by summing over $c$ and Proposition 4 follows from Proposition 2 by summing over the appropriate double cosets. We use the convention that the normalization (Haar measure) is such that $I^{+} \cdot I^{+}=I^{+}$.

Proposition 3 Let $w \in W^{\text {ad }}$, let $\vec{w}=s_{i_{1}} \cdots s_{i_{\ell}}$ be a reduced word for $w$ and let $j \in\{0, \ldots, n\}$. Assume that if $w s_{j} \ll w$ then $i_{\ell}=j$ and let

$$
y_{\overrightarrow{w s_{j}}}(\vec{c})=y_{i_{1}}\left(c_{1}\right) \cdots y_{i_{\ell-1}}\left(c_{\ell-1}\right) \quad \text { and } \quad y_{\vec{w}}(\vec{c})=y_{i_{1}}\left(c_{1}\right) \cdots y_{i_{\ell}}\left(c_{\ell}\right)=y_{\overrightarrow{w s_{j}}}(\vec{c}) y_{j}\left(c_{\ell}\right) .
$$

Then
$y_{\vec{w}}(\vec{c}) I^{+} \cdot T_{s_{j}}= \begin{cases}y_{i_{1}}\left(c_{1}\right) \cdots y_{i_{\ell-1}}\left(c_{\ell-1}\right) I^{+}+\sum_{\tilde{c} \in \mathbb{k}^{\times}} y_{i_{1}}\left(c_{1}\right) \cdots y_{i_{\ell}}\left(c_{\ell}-\tilde{c}^{-1}\right) I^{+}, & \text {if } w s_{j}<+w, \\ \sum_{\tilde{c} \in \mathbb{k}} y_{i_{1}}\left(c_{1}\right) \cdots y_{i_{\ell}}\left(c_{\ell}\right) y_{j}(\tilde{c}) I^{+}, & \text {if } w<w s_{j} .\end{cases}$
Proposition 4 Let $\mathbb{k}=\mathbb{F}_{q}$ the finite field with $q$ elements. Let $w \in W^{\text {ad }}$ and $j \in$ $\{0, \ldots, n\}$. Then
$T_{w} T_{s_{j}}=\left\{\begin{array}{ll}T_{w s_{j}}, & \text { if } w<+w s_{j}, \\ (q-1) T_{w}+q T_{w s_{j}} . & \text { if } w s_{j}<+w,\end{array} \quad L^{w} T_{s_{j}}= \begin{cases}L^{w s_{j}}, & \text { if } w<w s_{j}, \\ q L^{w s_{j}}+(q-1) L^{w}, & \text { if } w s_{j}<w,\end{cases}\right.$

$$
\text { and } \quad X^{w} T_{s_{j}}= \begin{cases}X^{w s_{j}}, & \text { if } w<0 w s_{j} \\ q X^{w s_{j}}+(q-1) X^{w}, & \text { if } w s_{j}<0 w\end{cases}
$$

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# Stable Bases of the Springer Resolution and Representation Theory 

Changjian Su and Changlong Zhong


#### Abstract

In this expository note, we review the recent developments about Maulik and Okounkov's stable bases for the Springer resolution $T^{*}(G / B)$. In the cohomology case, we compute the action of the graded affine Hecke algebra on the stable basis, which is used to obtain the localization formulae. We further identify the stable bases with the Chern-Schwartz-MacPherson classes of the Schubert cells. This relation is used to prove the positivity conjecture of Aluffi and Mihalcea. For the K theory stable basis, we first compute the action of the affine Hecke algebra on it, which is used to deduce the localization formulae via root polynomial method. Similar as the cohomology case, they are also identified with the motivic Chern classes of the Schubert cells. This identification is used to prove the Bump-Nakasuji-Naruse conjecture about the unramified principal series of the Langlands dual group over non-Archimedean local fields. In the end, we study the wall R-matrices, which relate stable bases for different alcoves. As an application, we give a categorification of the stable bases via the localization of Lie algebras over positive characteristic fields.


Keywords Flag variety • Springer resolution • Stable bases • Hecke algebra

## 1 Introduction

Schubert calculus studies the cohomology and K-theory of the (partial) flag varieties $\mathcal{B}$. In this note, we study certain basis elements of Maulik and Okounkov, called the stable basis, in the equivariant cohomology and K-theory of the cotangent bundle $T^{*} \mathcal{B}$ (the Springer resolution). Pulling back to the zero section $\mathcal{B}$, the stable basis

[^18]becomes some natural classes in the equivariant cohomology and K-theory of the flag varieties (see Sects. 2.5 and 4.2).

Maulik and Okounkov introduced the stable basis in their study of quantum cohomology of quiver varieties [31]. Later, Okounkov and his collaborators introduced K-theory and elliptic cohomology versions of these bases [1, 34, 37]. They turn out to be very useful both in enumerative geometry and geometric representation theory [35, 36].

The Maulik-Okounkov stable bases are constructed for a class of varieties called symplectic resolutions [26]. We will focus on the Springer resolution, which is one of the classical examples of symplectic resolutions. Both the cohomological and Ktheoretic stable bases for the Springer resolution are related to standard objects for certain representation categories.

In the cohomology case, the stable basis elements are just some rational combination of the conormal bundles of the Schubert cells, which were shown to coincide with characteristic cycles of certain $\mathcal{D}$-modules on the flag variety. Via the localization theorem of Beilinson and Bernstein [6], these $\mathcal{D}$-modules correspond to the Verma modules of the Lie algebra. Furthermore, the action of the convolution algebra of the Steinberg variety [21], which is isomorphic to the graded affine Hecke algebra [29], is computed under the stable basis [43]. From this, the first author deduced the localization formula for the stable basis, which is a direct generalizaiton of the well-known AJS/Billey formula [13] for the localization of Schubert classes in the equivariant cohomology of flag varieties. The restriction formulae also play a crucial role in determining the quantum connection of the cotangent bundle of partial flag varieties [43].

On the other hand, the graded affine Hecke algebra also appears in the work of Aluffi and Mihalcea [2] about the Chern-Schwartz-MacPherson (CSM) classes [30, 41, 42] of Schubert cells. Comparing the actions, we identify the pullbacks of the stable bases with the CSM classes [3] (see also [40]). The effectivity of the characteristic cycles of $\mathcal{D}$-modules enables the authors in [3] to prove the non-equivariant positivity conjecture of Aluffi and Mihalcea [2] for CSM classes of the Schubert cells.

The K-theoretic case is similar, but with richer structure. By the famous theorem of Kazhdan, Lusztig and Ginzburg [21, 23, 27], the affine Hecke algebra is isomorphic to the convolution algebra of the equivariant K-theory of the Steinberg variety. Hence, the affine Hecke algebra acts on the equivariant K-theory of $T^{*} \mathcal{B}$. This action on the stable basis is computed in [44]. It roughly says that the stable basis is the standard basis for the regular representation of the finite Hecke algebra. With this action, we can compute the localization of the stable basis via the root polynomial method.

The K-theretic generalization of the CSM classes are motivic Chern classes defined by Brasselet, Schürmann and Yokura [14]. Pulling back the K-theory stable bases to the zero section, we get the motivic Chern classes of the Schubert cells [4, 22]. By definition, the motivic Chern classes enjoy good functorial properties. These facts enable the authors in [4] to use Schubert calculus to prove some conjectures of Bump, Nakasuji and Naruse [17,32] about unramified principal series of the $p$-adic Langland dual group.

Recall that the affine Hecke algebra can be realized either as the equivariant K-group of the Steinberg variety, or as the double Iwahori-invariant functions on the $p$-adic Langland dual group [21, Introduction], [7]. The bridge connecting the Schubert calculus over complex numbers and the representation theory of the $p$-adic Langland dual group is the shadow of these two geometric realizations. To be more precise, the (specialized) equivariant K-theory of the flag variety and the Iwahoriinvariants of an unramified principal series of the $p$-adic group are two geometric realizations of the regular representation of the finite Hecke algebra. Identifying these two regular representations, the stable basis (or the motivic Chern classes of the Schubert cells) gets identified with the standard basis on the $p$-adic side, while the fixed point basis maps to the Casselman basis [18]. With this isomorphism, we give an equivariant K-theory interpretation of Macdonald's formula for the spherical function [18] and the Casselman-Shalika formula [19] for the spherical Whittaker function in [44].

The stable basis is also related to representations of Lie algebras over a field of positive characteristic, via the localization theorem of Bezrukavnikov, Mirković and Rumynin [9, 10], which generalizes the famous localization theorem of Beilinson and Bernstein [6] over the complex numbers. This is achieved as follows. The definition of the K-theory stable basis depends on a choice of the alcove. Changing alcoves defines the so-called wall R-matrices [37]. We first recall these wall R-matrices for the Springer resolution, which are computed in [45]. The formulae can be nicely packed using the Hecke algebra actions. It turns out the wall R-matrices coincide with the monodromy matrices of the quantum connection. Bezrukavnikov and Okounkov conjecture that for general symplectic resolutions $X$, the monodromy representation is isomorphic to the representation coming from derived equivalences [11, 35]. The Springer resolution case is established by Braverman, Maulik and Okounkov in [16]. Thus, the wall R-matrices is also related to the derived equivalences. Finally, we give a categorification of the stable basis via the affine braid group action, constructed by Bezrukavnikov and Riche, on the derived category of coherent sheaves on the Springer resolution [12, 38]. Moreover, the categorified stable bases are identified with the Verma modules for the Lie algebras over positive characteristic fields under the localization equivalence.

This note is structured as follows. In Sect.2, we review the following aspects of the cohomological stable basis: the action of the affine graded Hecke algebra, restriction formula, and identification with the Chern-Schwartz-MacPherson classes of the Schubert cells and characteristic cycles of holonomic $\mathcal{D}$-modules on the flag variety. The rest of the note is devoted to the study of the K theory stable basis. In Sect. 3, we introduce the K theory stable basis, compute the action of the affine Hecke algebra on it, and obtain the restriction formula. In Sect. 4, the K theory stable bases are identified with the motivic Chern classes of the Schubert cells. In Sect. 5, we further relate the stable basis to the Casselman basis in the Iwahori invariants of an unramified principal series of the Langlands dual group over a non-Archimedean local field, and use it to prove conjectures of Bump, Nakasuji and Naruse. Finally, in Sects. 6 and 7, we review the wall crossing matrix and the categorification of the stable bases.

## 2 Cohomological Stable Basis

In this section, we first recall the definition of the stable basis of Maulik and Okounkov [31]. Then we compute the action of the graded Hecke algebra on the stable basis, which will have two applications: the localization formulae and the identification of the stable bases with the Chern-Schwartz-MacPherson classes of the Schubert cells. We further relate the stable bases to characteristic cycles of $\mathcal{D}$-modules, which plays an important role in proving the positivity conjecture of Aluffi and Mihalcea.

### 2.1 Notations

Let $G$ be a complex, semisimple, simply connected linear algebraic group with a Borel subgroup $B$ and a maximal torus $A$. Let $B^{-}$be the opposite Borel subgroup. Let $\Lambda$ (resp. $\Lambda^{\vee}$ ) be the group of characters (resp. cocharacters) of $A$. Let $R^{+}$denote the roots in $B$. Let $\mathfrak{g}$ (resp. $\mathfrak{h}$ ) be the Lie algebra of $G$ (resp. $A$ ). Let $\rho$ be the half-sum of the positive roots. Let $\mathfrak{C}_{ \pm} \subset \mathfrak{h}$ be the dominant/anti-dominant Weyl chamber in $\mathfrak{h}$. Let $\mathcal{B}=G / B$ be the variety of all Borel subgroups of $G$. Denote by $T^{*} \mathcal{B}$ (resp. $T \mathcal{B}$ ) the cotangent (tangent) bundle. The cotangent bundle $T^{*} \mathcal{B}$ is a resolution of the nilpotent cone $\mathcal{N} \subset \mathfrak{g}$, which is the so-called Springer resolution. The group $\mathbb{C}^{*}$ acts on $T^{*} \mathcal{B}$ by $z \cdot\left(B^{\prime}, x\right)=\left(B^{\prime}, z^{-2} x\right)$ for any $z \in \mathbb{C}^{*},\left(B^{\prime}, x\right) \in T^{*} \mathcal{B}$. Let $-\hbar$ be the Lie $\mathbb{C}^{*}$-weight of the cotangent fiber. Let $T=A \times \mathbb{C}^{*}$. Then $H_{T}^{*}(p t)=\mathbb{C}[\mathfrak{h}][\hbar]$, where $\mathbb{C}[\mathfrak{h}]$ denotes the functions on $\mathfrak{h}$. For any $\lambda \in \Lambda$, let $\mathcal{L}_{\lambda}=G \times_{B} \mathbb{C}_{\lambda}$ be the line bundle on $\mathcal{B}$. Pulling back to $T^{*} \mathcal{B}$, we get a line bundle on $T^{*} \mathcal{B}$, which is still denoted by $\mathcal{L}_{\lambda}$. Let $W$ be the Weyl group, and let $\leq$ denote the Bruhat order on $W$. The $A$-fixed points in $T^{*} \mathcal{B}$ are indexed by the Weyl group $W$. For each $w \in W$, the corresponding fixed point is $w B \in \mathcal{B} \subset T^{*} \mathcal{B}$. Let $\iota_{w}: w B \hookrightarrow T^{*} \mathcal{B}$ denote the embedding. For any $\gamma \in H_{T}^{*}\left(T^{*} \mathcal{B}\right)$, let $\left.\gamma\right|_{w}=\iota_{w}^{*} \gamma \in H_{T}^{*}(w B)=H_{T}^{*}(p t)$.

### 2.2 Definition of the Cohomological Stable Basis

The definition of the stable basis depends on a choice of a Weyl chamber $\mathfrak{C}$ in $\mathfrak{h}_{\mathbb{R}}$. Pick one parameter $\sigma: \mathbb{C}^{*} \rightarrow A$, such that $d \sigma \in \mathfrak{C}$. Recall the attracting set of the torus fixed point $w B \in T^{*} \mathcal{B}$

$$
\operatorname{Attr}_{\mathfrak{C}}(w)=\left\{\left.x \in T^{*} \mathcal{B}\right|_{z \rightarrow 0} \sigma(z) \cdot x=w B\right\}
$$

A partial order on $W$ is defined as

$$
v \preceq_{\mathfrak{C}} w \text { if } \overline{\operatorname{Attr}_{\mathfrak{C}}(w)} \cap v B \neq \emptyset
$$

The order determined by the dominant chamber (resp. anti-dominant chamber) is the usual Bruhat order (resp., the opposite Bruhat order). The full attracting set is

$$
\operatorname{FAttr}_{\mathfrak{C}}(w)=\cup_{v \leq \mathfrak{c} w} \operatorname{Attr}_{\mathfrak{C}}(v)
$$

For each $w \in W$, let $\epsilon_{w}=e^{A}\left(T_{w B}^{*} \mathcal{B}\right) \in H_{A}^{*}(p t)=\mathbb{C}[\mathfrak{h}]$, the $A$-equivariant Euler class of the $A$-vector space $T_{w B}^{*} \mathcal{B}$. That is, $\epsilon_{w}$ is the product of $A$-weights in the vector space $T_{w B}^{*} \mathcal{B}$. Let $N_{w}=T_{w B}\left(T^{*} \mathcal{B}\right)$. The chamber $\mathfrak{C}$ gives a decomposition of the tangent bundle

$$
N_{w}=N_{w,+} \oplus N_{w,-}
$$

into $A$-weights which are positive and negative when paired with $\sigma$, respectively. Let $e\left(N_{w,-}\right) \in H_{T}^{*}(p t)$ denote the $T$-equivariant Euler class, that is, the product of $T$-weights in the $T$-equivariant vector space $N_{w,-}$. Since $T=A \times \mathbb{C}^{*}, H_{T}^{*}(p t)=$ $H_{A}^{*}(p t)[\hbar]$. If we specialize the equivariant parameter $\hbar$ to zero, the $T$-equivariant Euler class $e\left(N_{w,-}\right) \in H_{T}^{*}(p t)$ equals $a_{w} \epsilon_{w}$, where $a_{w} \in\{1,-1\}$. See Example 2.3.

Recall the fixed point set $\left(T^{*} \mathcal{B}\right)^{A}$ is a discrete set indexed by the Weyl group $W$. Hence, $H_{T}^{*}\left(\left(T^{*} \mathcal{B}\right)^{A}\right)=\sqcup_{w \in W} H_{T}^{*}(p t)$. Let $1_{w}$ denote the unit element in the $w$-th copy $H_{T}^{*}(p t)$. The cohomological stable basis is defined as follows.

Theorem 2.1 ([31, Theorem 3.3.4]) There exists a unique map of $H_{T}^{*}(p t)$-modules

$$
\operatorname{stab}_{\mathfrak{C}}: H_{T}^{*}\left(\left(T^{*} \mathcal{B}\right)^{A}\right) \rightarrow H_{T}^{*}\left(T^{*} \mathcal{B}\right)
$$

satisfying the following properties. For any $w \in W$, denote $\operatorname{stab}_{\mathfrak{C}}(w)=\operatorname{stab}_{\mathfrak{C}}\left(1_{w}\right)$. Then
(1) (Support) $\operatorname{supp} \operatorname{stab}_{\mathfrak{C}}(w) \subset \operatorname{FAttr}_{\mathfrak{C}}(w)$,
(2) (Normalization) $\left.\operatorname{stab}_{\mathfrak{C}}(w)\right|_{w}=a_{w} e\left(N_{-, w}\right)$,
(3) (Degree) $\left.\operatorname{stab}_{\mathfrak{C}}(w)\right|_{v}$ is divisible by $\hbar$, for any $v \prec_{\mathfrak{C}} w$.

Remark 2.2 (1) The map is defined by a Lagrangian correspondence between $\left(T^{*} \mathcal{B}\right)^{A} \times T^{*} \mathcal{B}$, hence mapping middle degree to middle degree. Therefore, the last condition is equivalent to

$$
\left.\operatorname{deg}_{A} \operatorname{stab}_{\mathfrak{C}}(w)\right|_{v}<\left.\operatorname{deg}_{A} \operatorname{stab}_{\mathfrak{C}}(v)\right|_{v}
$$

Here, for any polynomial $f \in H_{T}(p t)=\mathbb{C}[\mathfrak{h}][\hbar]=\operatorname{Sym}\left(\mathfrak{h}^{*}\right)[\hbar], \operatorname{deg}_{A} f$ denotes the degree of $f$ in the $\mathfrak{h}^{*}$-variables.
(2) By the first and second conditions above, $\left\{\operatorname{stab}_{\mathfrak{C}}(w) \mid w \in W\right\}$ is a basis for the localized equivariant cohomology $H_{T}^{*}\left(T^{*} \mathcal{B}\right)_{l o c}:=H_{T}^{*}\left(T^{*} \mathcal{B}\right) \otimes_{H_{T}(p t)}$ Frac $H_{T}(p t)$. It is the so-called stable basis.
(3) By [31, Theorem 4.4.1], the stable bases for opposite chambers are dual bases. That is,

$$
\left\langle\operatorname{stab}_{\mathfrak{C}}(v), \operatorname{stab}_{-\mathfrak{C}}(w)\right\rangle=(-1)^{\operatorname{dim} G / B} \delta_{v, w}
$$

where the non-degenerate pairing $\langle-,-\rangle$ on $H_{T}^{*}\left(T^{*} \mathcal{B}\right)$ is defined via localization as follows

$$
\left\langle\gamma_{1}, \gamma_{2}\right\rangle:=\sum_{w} \frac{\left.\left.\gamma_{1}\right|_{w} \cdot \gamma_{2}\right|_{w}}{e\left(N_{w}\right)} \in \operatorname{Frac} H_{T}^{*}(p t)
$$

Let $\operatorname{stab}_{+}(w)=\operatorname{stab}_{\mathfrak{C}_{+}}(w)$ and $\operatorname{stab}_{-}(w)=\operatorname{stab}_{\mathfrak{C}_{-}}(w)$.
Example 2.3 Let $G=\operatorname{SL}(2, \mathbb{C}), B$ be the upper triangular matrices, and $\alpha$ be the simple root. In this case, the flag variety $\mathcal{B}=\mathbb{P}^{1}$, the identity element $1 \in W$ corresponds to the fixed point $0:=[1: 0] \in \mathbb{P}^{1}$, and $s_{\alpha} \in W$ corresponds to the other fixed point $\infty:=[0: 1] \in \mathbb{P}^{1}$. The torus $T=A \times \mathbb{C}^{*}$ weights of the vector spaces are as follows

$$
\text { weight }\left(T_{0} \mathbb{P}^{1}\right)=-\alpha, \quad \text { weight }\left(T_{0}^{*} \mathbb{P}^{1}\right)=\alpha-\hbar,
$$

and

$$
\text { weight }\left(T_{\infty} \mathbb{P}^{1}\right)=\alpha, \quad \text { weight }\left(T_{\infty}^{*} \mathbb{P}^{1}\right)=-\alpha-\hbar
$$

Let us pick the positive Weyl chamber $\mathfrak{C}_{+}$so that $\alpha$ takes non-negative values on it. Then $\operatorname{Attr}_{\mathfrak{C}_{+}}(0)=T_{0}^{*} \mathbb{P}^{1}$, and $\operatorname{Attr}_{\mathfrak{C}_{+}}(\infty)=\mathbb{P}^{1} \backslash\{0\}$. Hence, $0<\mathfrak{C}_{+} \infty$. The Euler classes are $\epsilon_{0}=e^{A}\left(T_{0}^{*} \mathbb{P}^{1}\right)=\alpha, e\left(N_{0,-}\right)=-\alpha, \epsilon_{\infty}=e^{A}\left(T_{\infty}^{*} \mathbb{P}^{1}\right)=-\alpha$, and $e\left(N_{\infty,-}\right)=-\alpha-\hbar$. Therefore, $a_{1}=-1$ and $a_{s_{\alpha}}=1$. By the support and normalization conditions, we get

$$
\operatorname{stab}_{+}(0)=-\left[T_{0}^{*} \mathbb{P}^{1}\right], \quad \text { and } \quad \operatorname{stab}_{+}(\infty)=\left[\mathbb{P}^{1}\right]+c\left[T_{0}^{*} \mathbb{P}^{1}\right]
$$

where $c \in H_{T}^{*}(p t)$. Restricting $\operatorname{stab}_{\mathfrak{C}_{+}}(\infty)$ to the smaller fixed point 0 , we get

$$
\left.\operatorname{stab}_{+}(\infty)\right|_{0}=(-\hbar+\alpha)+c(-\alpha)
$$

By the degree condition, the above is divisible by $\hbar$. Hence, $c=1$. So

$$
\begin{equation*}
\operatorname{stab}_{+}(0)=-\left[T_{0}^{*} \mathbb{P}^{1}\right], \quad \operatorname{stab}_{+}(\infty)=\left[\mathbb{P}^{1}\right]+\left[T_{0}^{*} \mathbb{P}^{1}\right] \tag{1}
\end{equation*}
$$

Similarly, for the negative chamber $\mathfrak{C}_{-}$, we have

$$
\begin{equation*}
\operatorname{stab}_{-}(0)=\left[T_{\infty}^{*} \mathbb{P}^{1}\right]+\left[\mathbb{P}^{1}\right], \quad \text { stab }-(\infty)=-\left[T_{\infty}^{*} \mathbb{P}^{1}\right] \tag{2}
\end{equation*}
$$

### 2.3 The Graded Affine Hecke Algebra Action

The graded affine Hecke algebra $\mathcal{H}_{\hbar}$ is generated by the elements $x_{\lambda}$ of degree 1 for $\lambda \in \mathfrak{h}^{*}, s_{i} \in W$ of degree 0 and a central element $\hbar$ of degree 1 such that
(1) $x_{\lambda+\mu}=x_{\lambda}+x_{\mu}$ for any $\lambda, \mu \in \mathfrak{h}^{*}$;
(2) $x_{\lambda} x_{\mu}=x_{\mu} x_{\lambda}$ for any $\lambda, \mu \in \mathfrak{h}^{*}$;
(3) The $s_{i}$ 's generated the Weyl group inside $\mathcal{H}_{\hbar}$;
(4) For any simple root $\alpha_{i}$ and $\lambda \in \mathfrak{h}^{*}$, we have

$$
s_{i} x_{\lambda}-x_{s_{i} \lambda} s_{i}=\hbar\left(\lambda, \alpha_{i}^{\vee}\right)
$$

According to [29], there is a natural algebra isomorphism

$$
\mathcal{H}_{\hbar} \simeq H_{*}^{G \times \mathbb{C}^{*}}(Z),
$$

where $Z=T^{*} \mathcal{B} \times_{\mathcal{N}} T^{*} \mathcal{B}$ is the Steinberg variety, $H_{*}^{G \times \mathbb{C}^{*}}(Z)$ is the $G \times \mathbb{C}^{*}$ equivariant Borel-Moore homology of $Z$, and it is endowed with the convolution algebra structure [21]. The isomorphism is constructed as follows. Recall that for any $\lambda \in \Lambda$, we have the line bundle $\mathcal{L}_{\lambda} \in \operatorname{Pic}_{A \times \mathbb{C}^{*}}\left(T^{*} \mathcal{B}\right)$. Then the above isomorphism sends $x_{\lambda}$ to the push-forward of the first Chern class $c_{1}\left(\mathcal{L}_{\lambda}\right)$ under the diagonal embedding $T^{*} \mathcal{B} \hookrightarrow Z$. The image of a simple reflection $s_{\alpha}$ in the Weyl group is constructed as follows. Let $P_{\alpha}$ denote the corresponding minimal parabolic subgroup containing $B$. Let $\mathcal{P}_{\alpha}=G / P_{\alpha}$ and $Y_{\alpha}=\mathcal{B} \times_{\mathcal{P}_{\alpha}} \mathcal{B} \subset \mathcal{B} \times \mathcal{B}$. The conormal bundle $T_{Y_{\alpha}}^{*}(\mathcal{B} \times \mathcal{B})$ is denoted by $T_{Y_{\alpha}}^{*}$, which is a smooth closed $G \times \mathbb{C}^{*}$-invariant subvariety of $Z$. Then $s_{\alpha}-1$ is sent to the cohomology class $\left[T_{Y_{\alpha}}^{*}\right] \in H_{*}^{G \times \mathbb{C}^{*}}(Z)$.

Via this isomorphism, $\mathcal{H}_{\hbar}$ acts on $H_{T}^{*}\left(T^{*} \mathcal{B}\right)$ by convolution [21, Chap.2]. Let $\pi$ denote this action. The action on the stable basis is given by

Theorem 2.4 ([43, Lemma 3.2]) For any $w \in W$ and simple root $\alpha$,

$$
\pi\left(s_{\alpha}\right)\left(\operatorname{stab}_{ \pm}(w)\right)=-\operatorname{stab}_{ \pm}\left(w s_{\alpha}\right)
$$

### 2.4 The Restriction Formula

One important corollary of the above theorem is the restriction formula for the stable basis.

Theorem 2.5 ([43, Theorem 1.1]) Let $y=s_{1} s_{2} \cdots s_{l}$ be a reduced expression for $y \in W$. Then

$$
\begin{equation*}
\left.\operatorname{stab}_{-}(w)\right|_{y}=(-1)^{l(y)} \prod_{\alpha \in R^{+} \backslash R(y)}(\alpha-\hbar) \sum_{\substack{1 \leq i_{1}<i_{2}<\cdots<i_{j} \leq l \\ w=i_{1} s_{i_{2}} \ldots i_{i_{k}}}} \hbar^{l-k} \prod_{j=1}^{k} \beta_{i_{j}}, \tag{3}
\end{equation*}
$$

where $s_{i}$ is the simple reflection associated to a simple root $\alpha_{i}, \beta_{i}=s_{1} \cdots s_{i-1} \alpha_{i}$, $R^{+}$is the set of positive roots, and $R(y)=\left\{\beta_{i} \mid 1 \leq i \leq l\right\}$. Furthermore, the sum in Eq.(3) does not depend on the reduced expression for $y$.

From this formula, we can recover the AJS/Billey formula [13]

$$
\left.\left[\overline{B^{-} w B / B}\right]\right|_{y}=\sum_{\substack{1 \leq i_{1}<i_{2}<\cdots<i_{k} \leq l \\ w=s_{i_{1}} s_{2} \cdots s_{i_{k}} \text { reduced }}} \beta_{i_{1}} \cdots \beta_{i_{k}},
$$

via the following limit formula ([43, Theorem 4.8])

$$
\left.\left[\overline{B^{-} w B / B}\right]\right|_{y}=(-1)^{\ell(w)} \lim _{\hbar \rightarrow \infty} \frac{\left.\operatorname{stab}_{-}(w)\right|_{y}}{(-\hbar)^{\operatorname{dim} \overline{B^{-} w B / B}}}
$$

Example 2.6 Let us consider the case $G=\operatorname{SL}(3, \mathbb{C}), y=s_{1} s_{2} s_{1}$ and $w=s_{1}$. Then $\left.\operatorname{stab}_{-}(w)\right|_{y}=(-1)^{3} \hbar^{2}\left(\alpha_{1}+\alpha_{2}\right)$, and $\left.\left[\overline{B^{-} w B / B}\right]\right|_{y}=\alpha_{1}+\alpha_{2}$. It is obvious that the above limit formula holds.

We can also compute restriction formula for stab_(w) for all $w \in S_{3}$ (following notations in Sect. 3.4, with $\alpha_{1+2}=\alpha_{1}+\alpha_{2}, \tilde{\alpha}=\alpha-\hbar, i j i=s_{i} s_{j} s_{i}$, ect):

$$
\begin{aligned}
\text { stab_(121) } & =-\alpha_{1} \alpha_{2} \alpha_{1+2} f_{121}, \\
\text { stab_(12) } & =-\hbar \alpha_{1} \alpha_{1+2} f_{121}+\alpha_{1} \alpha_{1+2} \tilde{\alpha}_{2} f_{12}, \\
\text { stab_(21) } & =-\hbar \alpha_{2} \alpha_{1+2} f_{121}+\alpha_{2} \alpha_{1+2} \tilde{\alpha}_{1} f_{21}, \\
\text { stab_(1) } & =-\hbar^{2} \alpha_{1+2} f_{121}+\hbar \alpha_{1} \tilde{\alpha}_{2} f_{12}+\hbar \alpha_{1+2} \tilde{\alpha}_{1} f_{21}-\alpha_{1} \tilde{\alpha}_{2} \tilde{\alpha}_{1+2} f_{1}, \\
\text { stab_( } \left._{1}\right) & =-\hbar^{2} \alpha_{1+2} f_{121}+\hbar \alpha_{2} \tilde{\alpha}_{1} f_{21}+\hbar \alpha_{1+2} \tilde{\alpha}_{2} f_{12}-\alpha_{2} \tilde{\alpha}_{1} \tilde{\alpha}_{1+2} f_{2}, \\
\text { stab_(e) } & =-\left(\hbar^{3}+\hbar \alpha_{1} \alpha_{2}\right) f_{121}+\hbar^{2} \tilde{\alpha}_{2} f_{12}+\hbar^{2} \tilde{\alpha}_{1} f_{21}-\hbar \tilde{\alpha}_{2} \tilde{\alpha}_{1+2} f_{1}-\hbar \tilde{\alpha}_{1} \tilde{\alpha}_{1+2} f_{2}+\tilde{\alpha}_{1} \tilde{\alpha}_{2} \tilde{\alpha}_{1+2} f_{e} .
\end{aligned}
$$

In the above formulae, the restriction $\left.\operatorname{stab}_{-}(w)\right|_{y}$ is given by the coefficient of $f_{y}$ in the identity of stab_( $w$ ).

### 2.5 Relation with Chern-Schwartz-MacPherson (CSM) Classes

Let us first recall the definition of Chern-Schwartz-MacPherson classes. For any algebraic variety $Y$ over $\mathbb{C}$, let $\mathcal{F}(Y)$ denote the group of constructible functions on $Y$. If $f: Y \rightarrow X$ is a proper morphism, we can define a pushforward $f_{*}: \mathcal{F}(Y) \rightarrow \mathcal{F}(X)$ by setting $f_{*}\left(1_{W}\right)(p)=\chi\left(f^{-1}(p) \cap W\right)$, where $W \subset Y$ is a locally closed subvariety, $1_{W}$ is the characteristic function of $W, p \in X$, and $\chi$ is the topological Euler characteristic. According to a conjecture attributed to Deligne and Grothendieck, there is a unique natural transformation $c_{*}: \mathcal{F} \rightarrow H_{*}$ from the functor $\mathcal{F}$ of constructible functions on a complex algebraic variety to the homology functor, such that if $X$ is smooth then $c_{*}\left(1_{X}\right)=c(T X) \cap[X]$, where $c(T X)$ is the total Chern class. The naturality of $c_{*}$ means that it commutes with proper pushforward. This conjecture was proved by MacPherson [30]. The class $c_{*}\left(1_{X}\right)$ for possibly singular
$X$ was shown to coincide with a class defined earlier by M.-H. Schwartz [41, 42]. For any constructible subset $W \subset X$, we call the class $c_{S M}(W):=c_{*}\left(1_{W}\right) \in H_{*}(X)$ the Chern-Schwartz-MacPherson (CSM) class of $W$ in $X$. The theory of CSM classes was later extended to the equivariant setting by Ohmoto [33].

In this section, we identity the equivariant homology $H_{*}^{A}(\mathcal{B})$ with the equivairant cohomology $H_{A}^{*}(\mathcal{B})$ via the Poincaré duality. Let $X(w)^{\circ}=B w B / B \subset \mathcal{B}$ and $Y(w)^{\circ}=B^{-} w B / B$ be the Schubert cells in $\mathcal{B}$. Let $X(w)=\overline{X(w)^{\circ}}$ be the Schubert variety. A similar formula to that given in Theorem 2.4 was also obtained by Aluffi and Mihalcea for the CSM classes for the Schubert cells $c_{S M}\left(X(w)^{\circ}\right)$ [2]. Comparing these formulae, it is easy to get the following relation between the stable bases and the CSM classes [3, Corollary 1.2] and [40].

Theorem 2.7 Let $i: \mathcal{B} \hookrightarrow T^{*} \mathcal{B}$ be the inclusion. For any $w \in W$,

$$
\left.i^{*}\left(\operatorname{stab}_{+}(w)\right)\right|_{\hbar=1}=(-1)^{\operatorname{dim} G / B} c_{S M}\left(X(w)^{\circ}\right) \in H_{A}^{*}(\mathcal{B})
$$

and

$$
\left.i^{*}\left(\operatorname{stab}_{-}(w)\right)\right|_{\hbar=1}=(-1)^{\operatorname{dim} G / B} c_{S M}\left(Y(w)^{\circ}\right) \in H_{A}^{*}(\mathcal{B})
$$

Example 2.8 Let us consider the case $G=\operatorname{SL}(2, \mathbb{C})$. We follow the notations in Example 2.3. By definition, $c_{S M}\left(X(i d)^{\circ}\right)=[X(i d)]=[0]$, and $c_{S M}\left(X\left(s_{\alpha}\right)^{\circ}\right)=$ $c_{S M}\left(X\left(s_{\alpha}\right)\right)-c_{S M}(X(i d))=c\left(T \mathbb{P}^{1}\right)-[0]=\left[\mathbb{P}^{1}\right]+[\infty]$. Let us check

$$
\begin{equation*}
\left.i^{*}\left(\operatorname{stab}_{+}\left(s_{\alpha}\right)\right)\right|_{\hbar=1}=-c_{S M}\left(X\left(s_{\alpha}\right)^{\circ}\right) \in H_{A}^{*}(\mathcal{B}) \tag{4}
\end{equation*}
$$

By localization, it suffices to check the equality after restricting both sides to the fixed points. For the fixed point $0,\left.i^{*}\left(\operatorname{stab}_{+}\left(s_{\alpha}\right)\right)\right|_{0, \hbar=1}=-1$, and $\left.c_{S M}\left(X\left(s_{\alpha}\right)^{\circ}\right)\right|_{0}=1$. For the fixed point $\infty,\left.i^{*}\left(\operatorname{stab}_{+}\left(s_{\alpha}\right)\right)\right|_{\infty, \hbar=1}=-1-\alpha$, and $\left.c_{S M}\left(X\left(s_{\alpha}\right)^{\circ}\right)\right|_{\infty}=1+\alpha$. Thus, Eq. (4) holds.

The equivariant cohomology of the flag variety $\mathcal{B}$ has a natural basis, namely, the Schubert basis $\{[X(w)] \mid w \in W\}$. Thus we can expand the CSM classes in terms of this basis

$$
c_{S M}\left(X(w)^{\circ}\right)=\sum_{u} c(w, u)[X(u)] \in H_{A}^{*}(\mathcal{B})
$$

where $c(w, u) \in H_{A}^{*}(p t)$. It is conjectured by Aluffi and Mihalcea that [2]

$$
c(w, u) \in \mathbb{Z}_{\geq 0}[\alpha \mid \alpha>0] .
$$

The non-equivariant case of this conjecture is proved in [3], in which the relation between the stable basis and the characteristic cycles of holonomic $\mathcal{D}_{\mathcal{B}}$-modules plays an important role.

### 2.6 Characteristic Cycles of D-modules

For $w \in W$, let $M_{w}$ be the Verma module of highest weight $-w \rho-\rho$, a module over the universal enveloping algebra $U(\mathfrak{g})$. Let $\mathfrak{M}_{w}$ denote the holonomic $\mathcal{D}_{\mathcal{B}}$-module

$$
\mathfrak{M}_{w}=\mathcal{D}_{\mathcal{B}} \otimes_{U(\mathfrak{g})} M_{w} .
$$

The image of this regular holonomic $\mathcal{D}_{\mathcal{B}}$-module under the Riemann-Hilbert correspondence is the constructible complex $\mathbb{C}_{X(w)^{\circ}}[\ell(w)][6,15,24]$. The characteristic cycle $\operatorname{Char}\left(\mathfrak{M}_{w}\right)$ of $\mathfrak{M}_{w}$ is a linear combination of the conormal bundles of the Schubert cells. The relation with the stable basis is given by the following formula. It is claimed in [31, Remark 3.5.3], and later proved in [3, Lemma 6.5].

Theorem 2.9 For any $w \in W$,

$$
\operatorname{stab}_{+}(w)=(-1)^{\operatorname{dim} G / B-\ell(w)}\left[\operatorname{Char}\left(\mathfrak{M}_{w}\right)\right] \in H_{T}^{*}\left(T^{*} \mathcal{B}\right) .
$$

## 3 K-Theoretic Stable Bases and the Affine Hecke Algebra Action

From now on, we focus on the K theory stable basis. In this section, we introduce the definitions and compute the action of the affine Hecke algebras on the K theory stable bases. As an application, we compute the restriction formula of stable bases by using the root polynomial method.

### 3.1 Notations Continued

We will follow the notations in Sect. 2. Let us introduce more notations. Let $H_{\alpha^{\vee}, n}=$ $\left\{\lambda \in \mathfrak{h}_{\mathbb{R}}^{*} \mid\left(\lambda, \alpha^{\vee}\right)=n\right\}$ be the hyperplanes determined by the coroot $\alpha^{\vee}$ and integers $n$. The union of the hyperplanes is a closed subset of $\mathfrak{h}_{\mathbb{R}}^{*}$, whose complement has connected components called alcoves. The fundamental alcove is $\nabla_{+}=\left\{\lambda \in \mathfrak{h}_{\mathbb{R}}^{*} \mid 0<\right.$ $\left\langle\lambda, \alpha^{\vee}\right\rangle<1$, for any positive coroot $\left.\alpha^{\vee}\right\}$. Denote $\nabla_{-}=-\nabla_{+}$.

In the remaining parts of this note, we will consider the equivariant K-theory, a good introduction of which can be found in [21]. If a group $H$ acts on an algebraic variety $X$, then the $H$-equivariant K-theory of $X$, which is denoted by $K_{H}(X)$, is defined to be the Grothendieck group of the $H$-equivariant coherent sheaves on $X$. Namely, $K_{H}(X):=K^{0}\left(\operatorname{Coh}_{H}(X)\right)$. By definition, $K_{H}(X)$ is a module over $K_{H}(p t)=K^{0}(\operatorname{Rep}(H))$.

Recall the group $\mathbb{C}^{*}$ acts on $T^{*} \mathcal{B}$ by $z \cdot\left(B^{\prime}, x\right)=\left(B^{\prime}, z^{-2} x\right)$ for any $z \in \mathbb{C}^{*}$, $\left(B^{\prime}, x\right) \in T^{*} \mathcal{B}$. Let $q^{-1}$ be the $\mathbb{C}^{*}$-character of the cotangent fiber under this action, i.e., $q=e^{\hbar}$. Therefore, the $\mathbb{C}^{*}$-equivariant K-group of a point is $K_{\mathbb{C}^{*}}(p t)=$
$K^{0}\left(\operatorname{Rep}\left(\mathbb{C}^{*}\right)\right)=\mathbb{Z}\left[q^{1 / 2}, q^{-1 / 2}\right]$. Recall $T=A \times \mathbb{C}^{*}$. Denote $S=K_{T}(\mathrm{pt})=\mathbb{Z}\left[q^{1 / 2}\right.$, $\left.q^{-1 / 2}\right][\Lambda]$. We consider the $T$-equivariant K-theory of $T^{*} \mathcal{B}$. Recall $\iota_{w}$ denotes the inclusion of the fixed point $w B$ into $T^{*} \mathcal{B}$. For any $\mathcal{F} \in K_{T}\left(T^{*} \mathcal{B}\right)$, let $\left.\mathcal{F}\right|_{w}$ denote $\iota_{w}^{*} \mathcal{F} \in K_{T}(p t)$. By localization theorem [21, Chapter 5], the localized K-group $K_{T}\left(T^{*} \mathcal{B}\right)_{l o c}:=K_{T}\left(T^{*} \mathcal{B}\right) \otimes_{K_{T}(p t)}$ Frac $K_{T}(p t)$ has a fixed point basis $\left\{\iota_{w *} 1 \mid w \in\right.$ $W\}$.

The non-degenerate pairing on $K_{T}\left(T^{*} \mathcal{B}\right)$ can be defined using localization as follows

$$
\langle\mathcal{F}, \mathcal{G}\rangle=\sum_{w \in W} \frac{\left.\left.\mathcal{F}\right|_{w} \mathcal{G}\right|_{w}}{\Lambda^{\bullet} T_{w}\left(T^{*} \mathcal{B}\right)}=\sum_{w \in W} \frac{\left.\left.\mathcal{F}\right|_{w} \mathcal{G}\right|_{w}}{\prod_{\alpha>0}\left(1-e^{w \alpha}\right)\left(1-q e^{-w \alpha}\right)}, \quad \mathcal{F}, \mathcal{G} \in K_{T}\left(T^{*} \mathcal{B}\right)
$$

Here for each $T$-space $V, \bigwedge^{\bullet} V=\sum_{k \geq 0}(-1)^{k} \bigwedge^{k} V^{\vee}=\Pi\left(1-e^{-\alpha}\right) \in K_{T}(\mathrm{pt})$, where the product runs through all $T$-weights in $V$, counted with multiplicities.

### 3.2 Definition of the Stable Basis

The definition of the K theory stable basis depends on choices of a chamber, an alcove and a polarization.

A polarization $T^{1 / 2}$ is a solution of the following equation

$$
T^{1 / 2}+q^{-1}\left(T^{1 / 2}\right)^{\vee}=T\left(T^{*} \mathcal{B}\right) \quad \text { in } K_{T}\left(T^{*} \mathcal{B}\right)
$$

Denote $T_{\text {opp }}^{1 / 2}=q^{-1}\left(T^{1 / 2}\right)^{\vee}$. We will frequently focus on the two mutually opposite polarizations: $T \mathcal{B}$ and $T^{*} \mathcal{B}$. Recall $N_{w}=T_{w B}\left(T^{*} \mathcal{B}\right)$ is the normal bundle at $w$, and the chamber $\mathfrak{C}$ determines a decomposition $N_{w}=N_{w,+} \oplus N_{w,-}$ according to the sign of $A$-weights with respect to the chamber $\mathfrak{C}$. Let $N_{w}^{1 / 2}=N_{w} \cap\left(\left.T^{1 / 2}\right|_{w}\right)$, and similarly define $N_{w,+}^{1 / 2}$ and $N_{w,-}^{1 / 2}$. It follows that the square $\operatorname{root}\left(\frac{\operatorname{det} N_{w,-}}{\operatorname{det} N_{w}^{1 / 2}}\right)^{1 / 2}$ exists in $K_{T}\left(T^{*} \mathcal{B}\right)$.

For any Laurent polynomial $f=\sum_{\mu \in \Lambda} f_{\mu} e^{\mu} \in K_{T}(\mathrm{pt})$, with $f_{\mu} \in \mathbb{Z}\left[q^{1 / 2}\right.$, $\left.q^{-1 / 2}\right], e^{\mu} \in K_{A}(\mathrm{pt})$, define the $A$-degree of $f$ to be the Newton polygon

$$
\operatorname{deg}_{A} f=\text { Convex hull }\left(\left\{\mu \mid f_{\mu} \neq 0\right\}\right) \subset \Lambda \otimes \mathbb{R}
$$

Let $1_{w}$ denote the unit in the $w$-th copy of $K_{T}(p t)$ in $K_{T}\left(\left(T^{*} \mathcal{B}\right)^{A}\right)$. We can now recall the definition of the K-theory stable basis.

Theorem 3.1 [34] For any chamber $\mathfrak{C}$, polarization $T^{1 / 2}$, and alcove $\nabla$, there is a unique map of $K_{T}(\mathrm{pt})$-modules (called the stable envelope):

$$
\operatorname{stab}_{\mathfrak{C}, T^{1 / 2}, \nabla}: K_{T}\left(\left(T^{*} \mathcal{B}\right)^{A}\right) \rightarrow K_{T}\left(T^{*} \mathcal{B}\right)
$$

satisfying the following properties. Denote $\operatorname{stab}_{w}^{\mathfrak{C}, T^{1 / 2}, \nabla}=\operatorname{stab}_{\mathfrak{C}, T^{1 / 2}, \nabla}\left(1_{w}\right)$. Then
(1) (Support) $\operatorname{supp}\left(\operatorname{stab}_{w}^{\mathfrak{C}, T^{1 / 2}, \nabla}\right) \subset \operatorname{FAttr}_{\mathfrak{C}}(w)$.
(2) (Normalization) $\left.\operatorname{stab}_{w}^{\mathfrak{C}, T^{1 / 2}, \nabla}\right|_{w}=\left.(-1)^{\mathrm{rank} N_{w,+}^{1 / 2}}\left(\frac{\operatorname{det} N_{w,-}}{\operatorname{det} N_{w}^{1 / 2}}\right)^{1 / 2} \mathcal{O}_{\operatorname{Attr} \mathcal{C}(w)}\right|_{w}$.
(3) (Degree) $\operatorname{deg}_{A}\left(\left.\operatorname{stab}_{w}^{\mathfrak{C}, T^{1 / 2}, \nabla}\right|_{v}\right) \subset \operatorname{deg}_{A}\left(\left.\operatorname{stab}_{v}^{\mathfrak{C}, T^{1 / 2}, \nabla}\right|_{v}\right)+v \lambda-w \lambda$ for any $v \prec_{\mathfrak{C}}$ $w, \lambda \in \nabla$.

Note that the degree condition depends on the alcove $\nabla$ only, not on a particular $\lambda \in \nabla$. On the other hand, the normalization does not depend on the alcove.

Denote

$$
\operatorname{stab}_{w}^{+}=\operatorname{stab}_{w}^{\mathfrak{C}_{+}, T \mathcal{B}, \nabla_{-}}, \quad \operatorname{stab}_{w}^{-}=\operatorname{stab}_{w}^{\mathfrak{C}_{-}, T^{*} \mathcal{B}, \nabla_{+}} .
$$

We list some basic properties of stable bases [37, Proposition 1] and [44, Lemma 2.2]:
(1) The duality: $\left\langle\operatorname{stab}_{w}^{\mathfrak{c}, T^{1 / 2}, \nabla}, \operatorname{stab}_{v}^{-\mathfrak{C}, T_{\text {opp }}^{1 / 2},-\nabla}\right\rangle=\delta_{v, w} \in K_{T}\left(T^{*} \mathcal{B}\right)$.
(2) $\left.\operatorname{stab}_{w}^{+}\right|_{w}=q^{-\ell(w) / 2} \prod_{\beta>0, w \beta<0}\left(q-e^{w \beta}\right) \prod_{\beta>0, w \beta>0}\left(1-e^{w \beta}\right)$.
(3) $\left.\operatorname{stab}_{w}^{-}\right|_{w}=q^{\ell(w) / 2} \prod_{\beta>0, w \beta<0}\left(1-e^{-w \beta}\right) \prod_{\beta>0, w \beta>0}\left(1-q e^{-w \beta}\right)$.

See Example 3.5 for restrictions of $\operatorname{stab}_{w}^{-}$for all $w \in S_{3}$.
Example 3.2 Let us consider the case $G=\operatorname{SL}(2, \mathbb{C})$. We follow the notations in Example 2.3. The dual Cartan $\mathfrak{h}_{\mathbb{R}}^{*}=\mathbb{R} \alpha$ is one dimensional. The alcoves are $\nabla_{n}:=$ $\left(\frac{n \alpha}{2}, \frac{(n+1) \alpha}{2}\right)$, where $n \in \mathbb{Z}$. The stable basis are given by the following formulae. We refer the readers to [44, Example 1.4] for more details.

For the negative chamber, the cotangent polarization and the fundamental alcove $\nabla_{0}$, we have

$$
\operatorname{stab}_{i d}^{-}=\operatorname{stab}_{i d}^{\mathcal{C}_{d}, T^{*} \mathbb{P}^{1}, \nabla_{0}}=\left[\mathcal{O}_{\mathbb{P}^{1}}\right]+q e^{\alpha}\left[\mathcal{O}_{T_{\infty}^{*} \mathbb{P}^{\mathbb{1}}}\right], \text { and } \quad \operatorname{stab}_{s_{\alpha}}^{-}=\operatorname{stab}_{s_{\alpha}}^{\mathcal{C}_{-}, T^{*} \mathbb{P}^{1}, \nabla_{0}}=-q^{\frac{1}{2}} e^{\alpha}\left[\mathcal{O}_{T_{\infty}^{*} \mathbb{P}^{\mathbb{1}}}\right] \text {. }
$$

For a general alcove $\nabla_{n}$,

$$
\operatorname{stab}_{w}^{\mathfrak{C}_{-}, T^{*} \mathbb{P}^{1}, \nabla_{n}}=e^{-\frac{n}{2} w \alpha} \mathcal{L}_{\frac{n}{2} \alpha} \otimes \operatorname{stab}_{w}^{-},
$$

where $w=i d$ or $s_{\alpha}$.
For the opposite choices,
$\operatorname{stab}_{i d}^{+}=\left[\mathcal{O}_{T_{0}^{* \mathbb{P}}}\right]$, and $\operatorname{stab}_{s_{\alpha}}^{+}=-q^{-\frac{1}{2}} e^{-\alpha}\left[\mathcal{O}_{\mathbb{P}^{1}}\right]+\left(-q^{\frac{1}{2}} e^{-2 \alpha}+\left(q^{-\frac{1}{2}}-q^{\frac{1}{2}}\right) e^{-\alpha}\right)\left[\mathcal{O}_{T_{0}^{*} \mathbb{P}^{1}}\right]$.

### 3.3 The Action of the Affine Hecke Algebra

This section can be thought of as a K-theoretic generalization of the results in Sect.2.3. Recall the affine Hecke algebra $\mathbb{H}$ is a free $\mathbb{Z}\left[q, q^{-1}\right]$ module with basis $\left\{T_{w} e^{\lambda} \mid w \in W, \lambda \in \Lambda\right\}$, such that

- For any $\lambda, \mu \in \Lambda, e^{\lambda} e^{\mu}=e^{\lambda+\mu}$.
- For any simple root $\alpha,\left(T_{s_{\alpha}}+1\right)\left(T_{s_{\alpha}}-q\right)=0$.
- For any $w, y \in W$, such that $\ell(w y)=\ell(w)+\ell(y), T_{w} T_{y}=T_{w y}$.
- For any simple root $\alpha$ and $\lambda \in \Lambda$,

$$
T_{s_{\alpha}} e^{\lambda}-e^{s_{\alpha} \lambda} T_{s_{\alpha}}=(1-q) \frac{e^{s_{\alpha} \lambda}-e^{\lambda}}{1-e^{-\alpha}}
$$

Let $H(W)$ denote the finite Hecke algebra, the subalgebra of $\mathbb{H}$ generated by the $T_{w}^{\prime} s$. Recall the famous theorem of Kazhdan-Lusztig [27] and Ginzburg [21]

$$
\begin{equation*}
\mathbb{H} \simeq K_{G \times \mathbb{C}^{*}}(Z) \tag{5}
\end{equation*}
$$

where $Z=T^{*} \mathcal{B} \times_{\mathcal{N}} T^{*} \mathcal{B}$ is the Steinberg variety, and the right hand side is endowed with the convolution algebra structure [21, Chap. 5]. The isomorphism is constructed as follows. We use the notations in Sect. 2.3. Let $\pi_{i}: T_{Y_{\alpha}}^{*} \rightarrow T^{*} \mathcal{B}$ be the two projections. Let $\mathcal{O}_{\Delta}$ be the structure sheaf of the diagonal $\Delta\left(T^{*} \mathcal{B}\right) \subset T^{*} \mathcal{B} \times T^{*} \mathcal{B}$. Then the isomorphism in (5) sends the simple generator $T_{s_{\alpha}} \in H(W)$ to $-\left[\mathcal{O}_{\Delta}\right]-\pi_{2}^{*} \mathcal{L}_{\alpha} \in$ $K_{G \times \mathbb{C}^{*}}(Z)$ and $e^{\lambda} \in \mathbb{Z}[\Lambda]$ to $\mathcal{O}_{\Delta}(\lambda)$ [38, Proposition 6.1.5].

By convolution, $K_{G \times \mathbb{C}^{*}}(Z)$ acts on $K_{T}\left(T^{*} \mathcal{B}\right)$ [21, Chap. 5]. Since the kernel defin$\operatorname{ing} T_{\alpha}$ is not symmetric, the left and right convolution actions are different. Following [44], we use $T_{\alpha}$ (resp. $T_{\alpha}^{\prime}$ ) to denote the left (resp. right) convolution action of $T_{\alpha}$. From [44, Lemma 3.4], these two operators are adjoint to each other:

$$
\left\langle T_{\alpha}(\mathcal{F}), \mathcal{G}\right\rangle=\left\langle\mathcal{F}, T_{\alpha}^{\prime}(\mathcal{G})\right\rangle, \quad \forall \mathcal{F}, \mathcal{G} \in K_{T}\left(T^{*} \mathcal{B}\right)
$$

One of the main results in [44] is the computation of the affine Hecke algebra action on the stable bases. More precisely, we have

Theorem 3.3 ([44, Proposition 3.3, Theorem 3.5]) Let $\alpha$ be a simple root. Then

$$
\begin{align*}
T_{\alpha}\left(\operatorname{stab}_{w}^{-}\right) & = \begin{cases}(q-1) \operatorname{stab}_{w}^{-}+q^{1 / 2} \operatorname{stab}_{w s_{\alpha}}^{-}, & \text {if } w s_{\alpha}<w ; \\
q^{1 / 2} \operatorname{stab}_{w s_{\alpha}}^{-}, & \text {if } w s_{\alpha}>w .\end{cases}  \tag{6}\\
T_{\alpha}^{\prime}\left(\operatorname{stab}_{w}^{+}\right) & = \begin{cases}(q-1) \operatorname{stab}_{w}^{+}+q^{1 / 2} \operatorname{stab}_{w s_{\alpha}}^{+}, & \text {if } w s_{\alpha}<w ; \\
q^{1 / 2} \operatorname{stab}_{w s_{\alpha}}^{+}, & \text {if } w s_{\alpha}>w .\end{cases} \tag{7}
\end{align*}
$$

In particular,

$$
\operatorname{stab}_{w}^{-}=q^{\frac{\ell\left(w_{0} w\right)}{2}} T_{w_{0} w}^{-1}\left(\operatorname{stab}_{w_{0}}^{-}\right), \quad \text { and } \quad \operatorname{stab}_{w}^{+}=q^{-\frac{\ell(w)}{2}} T_{w^{-1}}^{\prime}\left(\operatorname{stab}_{i d}^{+}\right)
$$

In the proof of this theorem, an elementary but essential method called rigidity was used. Its simplest form says the following. If $p(z) \in \mathbb{C}\left[z^{ \pm 1}\right]$, then

$$
p(z) \text { is bounded as } z^{ \pm 1} \rightarrow \infty \Longleftrightarrow p(z) \text { is a constant. }
$$

More applications of this method can be found in the survey by Okounkov [34]. As a immediate corollary, we obtain recursive formulas for localizations of stable bases [44, Proposition 3.6]. This plays an important role in the identification between the stable basis and motivic Chern classes of Schubert cells (see Sect. 4), which was used to prove some conjectures of Bump, Nakasuji and Naruse in representations of $p$-adic Langlands dual groups in [4] (see Sect. 5).

### 3.4 Root Polynomials and Restriction Formulas

Since localization is such a powerful tool in calculations, it is important to study the localization of the stable basis. The cohomology case was done in [43] (see Sect. 2.4). For the K-theoretic case, we use the root polynomial method. The notations defining root polynomials are very technical, so here we just introduce a simplified version without explanation. Please see [44, Sects. 5, 6] for details. For simplicity, denote $S=\mathbb{Z}\left[q^{1 / 2}, q^{-1 / 2}\right][\Lambda]$ and $Q=\operatorname{Frac}(S)$. Let $Q_{W}^{*}=\operatorname{Hom}(W, Q)$, which is the set of all functions from the set $W$ to the set $Q . Q_{W}^{*}$ is a $Q$-module with basis $f_{w}$ defined by $f_{w}(v)=\delta_{w, v}$, and it is also a commutative ring with product $f_{w} f_{v}=\delta_{w, v} f_{v}$. Indeed $Q_{W}^{*} \cong Q \otimes_{S} K_{T}\left(T^{*} \mathcal{B}\right)$, and one can view stab ${ }_{w}^{ \pm}$inside $Q_{W}^{*}$. For example,

$$
\operatorname{stab}_{w}^{ \pm}=\left.\sum_{v \in W} \operatorname{stab}_{w}^{ \pm}\right|_{v} f_{v},\left.\quad \operatorname{stab}_{w}^{ \pm}\right|_{v} \in S \subset Q
$$

We now introduce the root polynomials. For each $w=s_{i_{1}} \cdots s_{i_{l}}$, denote $\beta_{j}=$ $s_{i_{1}} \cdots s_{i_{j-1}} \alpha_{j}$, and define the root polynomials as

$$
R_{w}=\prod_{j=1}^{l} h_{i_{j}}\left(\beta_{j}\right), \text { where } h_{i}(\beta)=\tau_{s_{i}}-\frac{q-1}{x_{-\beta}} \text { and } x_{\beta}=1-e^{-\beta}
$$

Here $\left\{\tau_{w} \mid w \in W\right\}$ are some formal symbols, generating the finite Hecke algebra $H(W)$, and we assume that they commute with the $x$-variables (so we can forget about the Bernstein relation). See Example 3.5 below. Because of that, when expressing $R_{w}=\sum_{v \leq w} K_{v, w} \tau_{v}, K_{v, w} \in Q$, the coefficients $K_{v, w}$ can be computed easily. To relate with localization formula of the stable basis, we have

Theorem 3.4 ([44, Theorem 6.3, 6.5]) For any $w \in W$, we have

$$
\operatorname{stab}_{w}^{-}=q^{\ell(w) / 2} \sum_{v \geq w} \prod_{\alpha>0, v^{-1} \alpha<0}\left(1-e^{\alpha}\right) \prod_{\alpha>0, v^{-1} \alpha>0}\left(1-q e^{-\alpha}\right) K_{w, v} f_{v} \in K_{T}\left(T^{*} \mathcal{B}\right)
$$

Note that similar formula can be obtained for $\operatorname{stab}_{w}^{+}$by using the operators $T_{v}^{\prime}$, but it is a bit more complicated, and was not included in [44].

Example 3.5 We consider the case of $S L_{3}$, in which case there are two simple roots $\alpha_{1}, \alpha_{2}$ and $W=S_{3}$. Let $w=s_{1} s_{2}$ with $\beta_{1}=\alpha_{1}, \beta_{2}=\alpha_{1}+\alpha_{2}$. Then

$$
\begin{array}{r}
R_{w}=h_{1}\left(\alpha_{1}\right) h_{2}\left(\alpha_{1}+\alpha_{2}\right)=\left(\tau_{s_{1}}-\frac{q-1}{x_{-\alpha_{1}}}\right)\left(\tau_{s_{2}}-\frac{q-1}{x_{-\alpha_{1}-\alpha_{2}}}\right) \\
=\tau_{s_{1} s_{2}}-\frac{q-1}{x_{-\alpha_{1}}} \tau_{s_{2}}-\frac{q-1}{x_{-\alpha_{1}-\alpha_{2}}} \tau_{s_{1}}+\frac{(q-1)^{2}}{x_{-\alpha_{1}} x_{-\alpha_{1}-\alpha_{2}}} .
\end{array}
$$

For instance, we have

$$
\begin{aligned}
\left.\operatorname{stab}_{s_{1}}^{-}\right|_{s_{1} s_{2}} & =q^{\ell\left(s_{1}\right) / 2} \prod_{\alpha>0, s_{2} s_{1} \alpha<0}\left(1-e^{\alpha}\right) \prod_{\alpha>0, s_{2} s_{1} \alpha>0}\left(1-q e^{-\alpha}\right) K_{s_{1}, s_{1} s_{2}} \\
& =q^{1 / 2}\left(1-e^{\alpha_{1}}\right)\left(1-e^{\alpha_{1}+\alpha_{2}}\right)\left(1-q e^{-\alpha_{2}}\right)\left(-\frac{q-1}{1-e^{\alpha_{1}+\alpha_{2}}}\right) \\
& =-q^{1 / 2}(q-1)\left(1-e^{\alpha_{1}}\right)\left(1-q e^{-\alpha_{2}}\right) .
\end{aligned}
$$

Similarly, we can get restriction formula for all stable basis with $w \in S_{3}$ (with $i j i:=$ $s_{i} s_{j} s_{i}, x_{\alpha}=1-e^{-\alpha}, \hat{x}_{\alpha}=1-q e^{-\alpha}, x_{ \pm i \pm j}=x_{ \pm \alpha_{i} \pm \alpha_{j}}$, and $\hat{x}_{ \pm i \pm j}=\hat{x}_{ \pm \alpha_{i} \pm \alpha_{j}}$ ):

$$
\begin{aligned}
q^{-3 / 2} \operatorname{stab}_{121}^{-} & =x_{-1} x_{-2} x_{-1-2} f_{121}, \\
q^{-1} \operatorname{stab}_{12}^{-} & =x_{-1} x_{-1-2} \hat{x}_{2} f_{12}-(q-1) x_{-1} x_{-1-2} f_{121}, \\
q^{-1} \operatorname{stab}_{21}^{-} & =x_{-2} x_{-1-2} \hat{x}_{1} f_{21}-(q-1) x_{-2} x_{-1-2} f_{121}, \\
q^{-1 / 2} \operatorname{stab}_{1}^{-} & =\hat{x}_{2} \hat{x}_{1+2} x_{-1} f_{1}-(q-1) \hat{x}_{2} x_{-1} f_{12}-(q-1) \hat{x}_{1} x_{-1-2} f_{21}+(q-1)^{2} x_{-1-2} f_{121}, \\
q^{-1 / 2} \operatorname{stab}_{2}^{-} & =\hat{x}_{1} \hat{x}_{1+2} x_{-2} f_{2}-(q-1) \hat{x}_{1} x_{-2} f_{21}-(q-1) \hat{x}_{2} x_{-1-2} f_{12}+(q-1)^{2} x_{-1-2} f_{121}, \\
\operatorname{stab}_{e}^{-} & =\hat{x}_{1} \hat{x}_{2} \hat{x}_{1+2} f_{e}-(q-1) \hat{x}_{2} \hat{x}_{1+2} f_{1}-(q-1) \hat{x}_{1} \hat{x}_{1+2} f_{2}+(q-1)^{2} \hat{x}_{2} f_{12}, \\
& +(q-1)^{2} \hat{x}_{1} f_{21}-\left[(q-1)^{3}+q(q-1) x_{-1} x_{-2}\right] f_{121} .
\end{aligned}
$$

## 4 Motivic Chern Classes of the Schubert Cells

In this section, as another application of the relationship between the affine Hecke algebra and the stable bases, we relate the stable bases to motivic Chern classes for the Schubert cells [4, 22].

### 4.1 Definition of Motivic Chern Classes

The motivic Chern classes are K-theoretic generalizations of the Chern-SchwartzMacPherson classes (see Sect.2.5). Let us first recall the definition of the motivic Chern classes, following [4, 14, 22]. Recall $A$ is the maximal torus in $G$. For any quasi-projective complex smooth $A$-variety $X$, let $G_{0}^{A}(v a r / X)$ be the free abelian
group generated by the isomorphism classes of algebraic morphisms $[f: Z \rightarrow X]$ where $Z$ is a quasi-projective $A$-variety and $f$ is an $A$-equivariant morphism modulo the usual additivity relations

$$
[f: Z \rightarrow X]=[f: U \rightarrow X]+[f: Z \backslash U \rightarrow X]
$$

for $U \subset Z$ an open $A$-invariant subvariety. Then there exists a unique natural transformation

$$
M C_{y}: G_{0}^{A}(\operatorname{var} / X) \rightarrow K_{A}(X)[y]
$$

satisfying the following properties:

- It is functorial with respect to $A$-equivariant proper morphisms of smooth, quasiprojective varieties.
- It satisfies the following normalization condition

$$
M C_{y}\left(\left[i d_{X}: X \rightarrow X\right]\right)=\lambda_{y}\left(T^{*} X\right):=\sum y^{i}\left[\wedge^{i} T^{*} X\right] \in K_{A}(X)[y]
$$

Here $y$ is a formal variable. The existence in the non-equivariant case was proved in [14]. The equivariant case is established in [4, 22].

### 4.2 Motivic Chern Classes of the Schubert Cells

Recall $X(w)^{\circ}=B w B / B \subset \mathcal{B}$ and $Y(w)^{\circ}=B^{-} w B / B \subset \mathcal{B}$ are the Schubert cells in $\mathcal{B}$, where $B^{-}$is the opposite Borel group. Recall the Serre duality functor on $X$

$$
\mathcal{D}(\mathcal{F})=R \mathcal{H o m}\left(\mathcal{F}, \omega_{\mathcal{B}}^{\bullet}\right)
$$

where $\mathcal{F} \in K_{A}(\mathcal{B})$ and $\omega_{\mathcal{B}}^{\bullet}=\mathcal{L}_{2 \rho}[\operatorname{dim} \mathcal{B}]$ is the dualizing complex of $\mathcal{B}$. We extend it to $K_{A}(\mathcal{B})\left[y, y^{-1}\right]$ by sending $y$ to $y^{-1}$. Let $i: \mathcal{B} \hookrightarrow T^{*} \mathcal{B}$ be the inclusion. Then the relation between the motivic Chern classes and the stable bases is following, which is a K-theoretic generalization of Theorem 2.7.

Theorem 4.1 ([4]) For any $w \in W$, we have

$$
q^{-\frac{\ell(w)}{2}} \mathcal{D}\left(i^{*} \operatorname{stab}_{w}^{+}\right)=M C_{-q^{-1}}\left(\left[X(w)^{\circ} \hookrightarrow \mathcal{B}\right]\right) \in K_{A}(\mathcal{B})\left[q, q^{-1}\right]
$$

and

$$
q^{\frac{\ell(w)}{2}-\operatorname{dim} \mathcal{B}_{i}^{*}\left(\operatorname{stab}_{w}^{-}\right) \otimes\left[\omega_{X}^{\bullet}\right]=M C_{-q^{-1}}\left(\left[Y(w)^{\circ} \hookrightarrow \mathcal{B}\right]\right) \in K_{A}(\mathcal{B})\left[q, q^{-1}\right] . . . . ~}
$$

Example 4.2 Let $G=\mathrm{SL}(2, \mathbb{C})$. We follow the notations in Example 2.3. By definition,

$$
M C_{y}\left(\left[Y\left(s_{\alpha}\right)^{\circ} \hookrightarrow \mathbb{P}^{1}\right]\right)=\left[\mathcal{O}_{\infty}\right]
$$

and

$$
M C_{y}\left(\left[Y(i d)^{\circ} \hookrightarrow \mathbb{P}^{1}\right]\right)=M C_{y}\left(\left[i d_{\mathcal{B}}: \mathbb{P}^{1} \rightarrow \mathbb{P}^{1}\right]\right)-M C_{y}\left(\left[Y\left(s_{\alpha}\right)^{\circ} \hookrightarrow \mathbb{P}^{1}\right]\right)=1+y\left[T^{*} \mathbb{P}^{1}\right]-\left[\mathcal{O}_{\infty}\right] .
$$

Let us check the second equation in Theorem 4.1 for $w=i d$. It suffices to check the equality after restricting to the fixed points. From Example 3.2, we have

$$
\operatorname{stab}_{i d}^{-}=\left[\mathcal{O}_{\mathbb{P}^{1}}\right]+q e^{\alpha}\left[\mathcal{O}_{T_{\infty}^{*} \mathbb{P}^{1}}\right]
$$

Therefore,

$$
\left.\left(q^{-1} i^{*}\left(\operatorname{stab}_{i d}^{-}\right) \otimes\left[\omega_{X}^{\bullet}\right]\right)\right|_{0}=q^{-1}\left(1-q e^{-\alpha}\right)\left(-e^{\alpha}\right)=1-q^{-1} e^{\alpha}
$$

and

$$
\left.\left(q^{-1} i^{*}\left(\operatorname{stab}_{i d}^{-}\right) \otimes\left[\omega_{X}^{\bullet}\right]\right)\right|_{\infty}=q^{-1}\left(1-q e^{\alpha}+q e^{\alpha}\left(1-e^{-\alpha}\right)\right)\left(-e^{-\alpha}\right)=\left(1-q^{-1}\right) e^{-\alpha}
$$

On the other hand,

$$
\left.M C_{-q^{-1}}\left(\left[Y(i d)^{\circ} \hookrightarrow \mathbb{P}^{1}\right]\right)\right|_{0}=1-q^{-1} e^{\alpha}
$$

and

$$
\left.M C_{-q^{-1}}\left(\left[Y(i d)^{\circ} \hookrightarrow \mathbb{P}^{1}\right]\right)\right|_{\infty}=1-q^{-1} e^{-\alpha}-\left(1-e^{-\alpha}\right)=\left(1-q^{-1}\right) e^{-\alpha}
$$

Thus, the second equation in Theorem 4.1 holds for $w=i d$. The other cases can be checked similarly.

It will be proved in [5] that

$$
(-q)^{-\operatorname{dim} \mathcal{B}} i^{*}\left(g r\left(i_{w!} \mathbb{Q}_{Y(w)^{\circ}}^{H}\right)\right) \otimes\left[\omega_{X}^{\bullet}\right]=M C_{-q^{-1}}\left(\left[Y(w)^{\circ} \hookrightarrow X\right]\right),
$$

where $i_{w}: Y(w)^{\circ} \hookrightarrow \mathcal{B}$ is the inclusion, $\mathbb{Q}_{Y(w)^{\circ}}^{H}$ is the $A$-equivariant shifted mixed Hodge module on $Y(w)^{\circ}$, and $g r\left(i_{w!} \mathbb{Q}_{Y(w)^{\circ}}^{H}\right)$ is the associated $A \times \mathbb{C}^{*}$-equivariant sheaf on the cotangent bundle of $\mathcal{B}$ [46, Sect. 2]. Since $i^{*}$ is an isomorphism, we get

$$
q^{\frac{\ell(w)}{2}} \operatorname{stab}_{w}^{-}=(-1)^{\operatorname{dim} \mathcal{B}} g r\left(i_{w!} \mathbb{Q}_{Y(w)^{\circ}}^{H}\right) .
$$

This is the K-theoretic analogue of Theorem 2.9.

## 5 Unramified Principal Series of $\boldsymbol{p}$-adic Langlands Dual Group

The relation between stable basis and motivic Chern classes is used in [4] to prove some conjectures of Bump, Nakasuji and Naruse [17, 32] about unramified principal series of the $p$-adic Langlands dual group. In the proof, we need to first build a relation between the stable basis (or motivic Chern classes of Schubert cells) and unramified principal series. This is explained in this section.

### 5.1 The Two Bases in Iwahori-Invariants of an Unramified Principal Series Representation

Let $F$ be a finite extension of $\mathbb{Q}_{p} .{ }^{1}$ Let $\mathcal{O}_{F}$ be the ring of integers, with a uniformizer $\varpi \in \mathcal{O}_{F}$, and residue field $\mathbb{F}_{q}$. Let $\check{G}$ be the split reductive group over $F$, which is the Langlands dual group of $G$. Let $\check{A} \subset \check{B}$ be the corresponding dual maximal torus and Borel subgorup. Let $I$ be the Iwahori subgroup.

Let $\tau$ be an unramified character of $\check{A}$, i.e., complex characters of $\check{A}(F)$ which are trivial on $\check{A}\left(\mathcal{O}_{F}\right)$. Recall the principal series representation $I(\tau):=\operatorname{Ind}_{\check{B}(F)}^{\mathscr{G}(F)}(\tau)$ is the induced representation, which consists of locally constant functions $f$ on $\check{G}(F)$ such that $f(b g)=\tau(b) \delta^{1 / 2}(b) f(g), b \in \check{B}(F)$, where $\delta(b)=\prod_{\alpha>0}\left|\alpha^{\vee}(b)\right|_{F}$ is the modulus function on the Borel subgroup.

It is well known that the affine Hecke algebra in Sect. 3.3 can also be realized as

$$
\mathbb{H}=\mathbb{C}_{c}[I \backslash \check{G}(F) / I],
$$

which is the so-called Iwahori-Hecke algebra. The algebra structure on the latter space is induced by convolution. Again by convolution, the Iwahori-Hecke algebra $\mathbb{H}$ acts from the right on the Iwahori-invariants $I(\tau)^{I}$.

We say an unramified character $\tau$ is regular if the stabilizer $W_{\tau}$ in the Weyl group is trivial. In this case, the space $\operatorname{Hom}_{\check{G}(F)}\left(I(\tau), I\left(w^{-1} \tau\right)\right)$ is one-dimensional. It is spanned by the following intertwining operator $\mathcal{A}_{w}^{\tau}$ :

$$
\mathcal{A}_{w}^{\tau}(f)(g):=\int_{\check{N}_{w}} f(w n g) d n,
$$

where $\check{N}_{w}=\check{N}(F) \cap w^{-1} \check{N}^{-}(F) w$ with $\check{N}$ (resp. $\check{N}^{-}$) being the unipotent radical of the Borel subgroup $\check{B}$ (resp. $\check{B}^{-}$).

There are two bases of $I(\tau)^{I}$. The first basis $\left\{\varphi_{w}^{\tau} \mid w \in W\right\}$, which is called the standard basis, is induced by the following decomposition

[^19]$$
\check{G}(F)=\sqcup_{w \in W} \check{B}(F) w I .
$$

That is, $\varphi_{w}^{\tau}$ is characterized by the following conditions:
(1) $\varphi_{w}^{\tau}$ is supported on $\check{B}(F) w I$;
(2) $\varphi_{w}^{\tau}(b w g)=\tau(b) \delta^{1 / 2}(b)$ for any $b \in \check{B}(F), g \in I$.

The other basis $\left\{f_{w}^{\tau} \mid w \in W\right\}$, which is the so-called Casselman's basis [18], is characterized by

$$
\mathcal{A}_{v}^{\tau}\left(f_{w}^{\tau}\right)(1)=\delta_{v, w}
$$

These bases played an important role in the computation of the Macdonald spherical function [18] and the Casselman-Shalika formula for the spherical Whittaker function [19].

### 5.2 The Comparison

Since $\tau$ is an unramified character of $\check{A}$, it is equivalent to a point $\tau \in A$. Therefore, we can evaluate the base ring $K_{A}(p t)$ at $\tau \in A$. We let $\mathbb{C}_{\tau}$ denote this evaluation representation of $K_{T}(p t)=K_{A \times \mathbb{C}^{*}}(p t)$. For any $\mathcal{F} \in K_{A \times \mathbb{C}^{*}}\left(T^{*} \mathcal{B}\right)$, let $\mathcal{F}_{-\rho}$ denote $\mathcal{F} \otimes \mathcal{L}(-\rho)$.

Now we can state the relation between these bases. The following theorem is a shadow of the two geometric realizations of the affine Hecke algebra

$$
K_{G \times \mathbb{C}^{*}}(Z) \simeq \mathbb{H} \simeq \mathbb{C}_{c}[I \backslash \check{G}(F) / I]
$$

We refer the readers to [44, Sects. 8.2 and 8.3] for a more precise statement of the following theorem.

Theorem 5.1 ([44, Theorem 8.4]) There is a unique isomorphism of right $\mathbb{H}$ modules

$$
\Psi: K_{A \times \mathbb{C}^{*}}\left(T^{*} \mathcal{B}\right) \otimes_{K_{A \times \mathbb{C}^{*}}(\mathrm{pt})} \mathbb{C}_{\tau} \rightarrow I(\tau)^{I},
$$

such that the equivariant parameter $q$ is mapped to the cardinality of the residue field $\mathbb{F}_{q}$ and

$$
\begin{aligned}
\frac{q^{\ell(w)}}{\prod_{\beta>0, w \beta>0}\left(1-e^{w \beta}\right) \prod_{\beta>0, w \beta<0}\left(q-e^{w \beta}\right)}\left(\iota_{w *} 1\right)_{-\rho} \mapsto & f_{w}^{\tau}, \\
\left(\operatorname{stab}_{w}^{-}\right)_{-\rho} \mapsto & q^{-\ell(w) / 2} \varphi_{w}^{\tau} .
\end{aligned}
$$

As applications of this theorem, in [44], we showed that the left Weyl group action on the K-theory side corresponds to the suitably normalized intertwiners on the representation theory side [44, Corollary 8.6], and provided some K-theoretic interpretations
of Macdonald's formula [18] [44, Theorem 8.8] and the Casselman-Shalika formula [19] [44, Theorem 8.11].

It is an interesting question to study the transition matrix between the two bases in $I(\tau)^{I}$. Define the following transition matrix coefficients $m_{u, w}$ by

$$
\sum_{w \geq u} \varphi_{w}=\sum_{w \in W} m_{u, w} f_{w} .
$$

For instance, the following special case

$$
m_{i d, w}=\prod_{\alpha>0, w^{-1} \alpha<0} \frac{1-q^{-1} e^{\alpha}(\tau)}{1-e^{\alpha}(\tau)}
$$

is the Gindikin-Karpelevich formula.
The authors of [4] used Theorems 5.1, 4.1 and some functorial properties for the motivic Chern classes to prove the following conjectures of Bump, Nakasuji and Naruse [17, 32].

Theorem 5.2 ([4]) For any $u \leq w \in W$,
(1)

$$
m_{u, w}=\prod_{\alpha>0, u \leq s_{\alpha} w<w} \frac{1-q^{-1} e^{\alpha}(\tau)}{1-e^{\alpha}(\tau)}
$$

if and only if the Schubert variety $Y(u)$ is smooth at the torus fixed point $e_{w}$.
(2) As a rational function of $\tau \in A$, the product

$$
\prod_{\alpha>0, u \leq s_{\alpha} w<w}\left(1-e^{\alpha}\right) m_{u, w}
$$

has no poles on the maximal torus $A$.

## 6 Wall Crossings for the Stable Bases

Note that the definition of the stable bases depend on the alcoves in $\mathfrak{h}_{\mathbb{R}}^{*}$. Stable bases for different alcoves are related by the so-called wall $R$-matrices[37]. In this section, we study the wall R-matrices for the Springer resolution, which will be used for the categorification in the next section. The main reference for these two sections is [45].

### 6.1 Wall Crossing Matrix

By uniqueness of the stable basis, it is immediate to see that for any $\mu \in \Lambda$,

$$
\begin{equation*}
\operatorname{stab}_{y}^{\mathfrak{C}, T^{1 / 2}, \nabla+\mu}=e^{-y \mu} \mathcal{L}_{\mu} \otimes \operatorname{stab}_{y}^{\mathfrak{C}, T^{1 / 2}, \nabla} \tag{8}
\end{equation*}
$$

Because of this property, instead of crossing all the walls on $H_{\alpha^{\vee}, n}$, it is enough to just cross the walls on the 0 hyperplanes $H_{\alpha^{\vee}, 0}$. From now on, let $\nabla_{1}, \nabla_{2}$ be two alcoves sharing a wall on $H_{\alpha^{\vee}, 0}$, and $\left(\lambda_{1}, \alpha^{\vee}\right)>0$ for any $\lambda_{1} \in \nabla_{1}$.

Another useful fact is [37, Theorem 1]:

$$
\operatorname{stab}_{y}^{\mathfrak{C}, T^{1 / 2}, \nabla_{2}}= \begin{cases}\operatorname{stab}_{y}^{\mathfrak{C}, T^{1 / 2}, \nabla_{1}}+f_{y}^{\nabla_{2} \leftarrow \nabla_{1}} \operatorname{stab}_{y s_{\alpha}}^{\mathfrak{C}, T^{1 / 2}, \nabla_{1}}, & \text { if } y s_{\alpha} \prec_{\mathfrak{C}} y ; \\ \operatorname{stab}_{y}^{\mathfrak{C}, T^{1 / 2}, \nabla_{1}}, & \text { if } y s_{\alpha} \succ_{\mathfrak{C} y,},\end{cases}
$$

where $f_{y}^{\nabla_{2} \leftarrow \nabla_{1}} \in K_{T}(\mathrm{pt})$.
Then we have
Theorem 6.1 ([45]) For any $y \in W$, we have

$$
\begin{gathered}
\operatorname{stab}_{y}^{\mathfrak{C}, T \mathcal{B}, \nabla_{2}}= \begin{cases}\operatorname{stab}_{y}^{\mathfrak{C}, T \mathcal{B}, \nabla_{1}}+\left(q^{1 / 2}-q^{-1 / 2}\right) \operatorname{stab}_{y}^{\mathfrak{C}, T \mathcal{B}, \nabla_{1}}, & \text { if } y s_{\alpha} \prec_{\mathfrak{C}} y ; \\
\operatorname{stab}_{y}^{\mathfrak{C}, T \mathcal{B}, \nabla_{1}}, & \text { if } y s_{\alpha} \succ_{\mathfrak{C}} y .\end{cases} \\
\operatorname{stab}_{y}^{-\mathfrak{C}, T^{*} \mathcal{B}, \nabla_{2}}= \begin{cases}\operatorname{stab}_{y}^{-\mathfrak{C}, T^{*} \mathcal{B}, \nabla_{1}}+\left(q^{1 / 2}-q^{-1 / 2}\right) \operatorname{stab}_{y}^{-\mathfrak{C}, T^{*} \mathcal{B}, \nabla_{1}}, & \text { if } y s_{\alpha} \succ_{\mathfrak{C}} y ; \\
\operatorname{stab}_{y}^{-\mathfrak{C}, T^{*} \mathcal{B}, \nabla_{1}}, & \text { if } y s_{\alpha} \prec_{\mathfrak{C}} y .\end{cases}
\end{gathered}
$$

This theorem is firstly proved when $\alpha$ is a simple root, by computing the action of $T_{\alpha}$ and $T_{\alpha}^{\prime}$ on the stable basis $\operatorname{stab}_{y}^{\mathfrak{C}, T \mathcal{B}, \nabla}$ and $\operatorname{stab}_{y}^{-\mathfrak{C}, T^{*} \mathcal{B}, \nabla}$ for any alcove $\nabla$. That is, we use rigidity to obtain a general version of Theorem 3.3. For the non-simple root case, we use the following equality [4] to reduce to the simple root case

$$
\begin{equation*}
w\left(\operatorname{stab}_{y}^{\mathfrak{C}, T \mathcal{B}, \nabla}\right)=\operatorname{stab}_{w y}^{w \mathfrak{C}, T \mathcal{B}, \nabla} \tag{9}
\end{equation*}
$$

where the LHS of the equation above denotes the left action of $w \in W$ on $\operatorname{stab}_{y}^{\mathfrak{C}, T \mathcal{B}, \nabla} \in$ $K_{T}\left(T^{*} \mathcal{B}\right)$.

### 6.2 Wall Crossing and Affine Hecke Algebra Actions

Combining Theorems 3.3 and 6.1, we get the following general formulae, which show that the wall crossing matrices and the affine Hecke algebra action (see Sect. 3) are compatible.

Theorem 6.2 ([45]) For any $x, y \in W$, we have

$$
\begin{align*}
\operatorname{stab}_{y}^{\mathfrak{C}_{-}, T^{*} \mathcal{B}, x \nabla_{+}} & =q_{x}^{-1 / 2} T_{x}\left(\operatorname{stab}_{y x}^{\mathfrak{C}_{-}, T^{*} \mathcal{B}, \nabla_{+}}\right)  \tag{10}\\
\operatorname{stab}_{y}^{\mathfrak{C}_{+}, T \mathcal{B}, x \nabla_{-}} & =q_{x}^{1 / 2}\left(T_{x^{-1}}^{\prime}\right)^{-1}\left(\operatorname{stab}_{y x}^{\mathfrak{C}_{+}, T \mathcal{B}, \nabla_{-}}\right) . \tag{11}
\end{align*}
$$

Therefore, Theorem 3.3, Eqs. (8) and (11) determine all the stable basis stab ${ }_{y}^{\mathfrak{C}_{+}, T \mathcal{B}, \nabla}$ for the dominant Weyl chamber $\mathfrak{C}_{+}$. The stable basis for the other chambers can be computed by Eq. (9). Thus all the stable basis element $\operatorname{stab}_{y}^{\mathfrak{C}, T \mathcal{B}, \nabla}$ can be calculated.

For general symplectic resolutions, Bezrukavnikov and Okounkov [11, 35] conjecture that the representation coming from the derived equivalence is isomorphic to the monodromy representation coming from quantum cohomology. The monodromy matrices of the quantum connection of $T^{*} \mathcal{B}$ is computed in [16, 20]. The above theorem shows that the monodromy matrices coincide with the wall R-matrices for the K stable bases. The relation to derived equivalences is explained in the next section.

## 7 Categorification and Localization in Positive Characteristic

In this final section, we give a categorification of the stable bases, and study its relation with representation of $\mathfrak{g}$ over positive characteristic fields under the localization equivalence of Bezrukavnikov, Mirković and Rumynin [9, 10].

### 7.1 Categorification of the Stable Basis via Affine Braid Group Action

Let $B_{\mathrm{aff}}$ (resp. $B_{\mathrm{aff}}^{\prime}$ ) be the affine braid group (resp. the extended affine braid group) with generators $\widetilde{s}_{\alpha}, \alpha \in I_{\text {aff }}$.

Let $G_{\mathbb{Z}}$ be a split $\mathbb{Z}$-form of the complex algebraic group $G$ and $A_{\mathbb{Z}} \subset B_{\mathbb{Z}} \subset G_{\mathbb{Z}}$ be the maximal torus and a Borel subgroup, respectively. Let $T_{\mathbb{Z}}=A_{\mathbb{Z}} \times_{\mathbb{Z}}\left(\mathbb{G}_{m}\right)_{\mathbb{Z}}$. Bezrukavnikov and Riche constructed an extended affine braid group action on $D_{T_{\mathbb{Z}}}^{\mathrm{b}}\left(T^{*} \mathcal{B}_{\mathbb{Z}}\right):=D^{\mathrm{b}} \operatorname{Coh}^{T_{\mathbb{Z}}}\left(T^{*} \mathcal{B}_{\mathbb{Z}}\right)[12,38]$, denoted by $J_{\widetilde{w}}^{\mathrm{R}}, \widetilde{w} \in B_{\text {aff }}^{\prime}$ (here R denotes the right action). Inspired by Theorem 6.2, we give the following definition.

Definition 7.1 Let $\lambda \in \Lambda_{\mathbb{Q}}$ be regular. We define $\mathfrak{s t a b}_{\lambda}^{\mathbb{Z}}(y) \in D_{T_{\mathbb{Z}}}^{\mathrm{b}}\left(T^{*} \mathcal{B}_{\mathbb{Z}}\right), y \in W$ as follows:

$$
\begin{aligned}
\mathfrak{s t a b} \mathfrak{b}_{\lambda_{0}}^{\mathbb{Z}}(i d) & =\mathcal{L}_{-\rho} \otimes \mathcal{O}_{T_{i d}^{*} \mathcal{B}_{\mathbb{Z}}}, \quad \lambda_{0} \in \nabla_{-}, \\
\mathfrak{s t a b}_{\lambda_{0}}^{\mathbb{Z}}(y) & =J_{\tilde{y}}^{\mathrm{R}} \mathfrak{s t a b}_{\lambda_{0}}^{\mathbb{Z}}(i d), \quad \lambda_{0} \in \nabla_{-}, \\
\mathfrak{s t a b} b_{\lambda}^{\mathbb{Z}}(y) & =\left(J_{\tilde{x}}^{\mathbb{R}}\right)^{-1} \mathfrak{s t a b}_{\lambda_{0}}^{\mathbb{Z}}(y x), \quad y, x \in W, x \lambda_{0}=\lambda, \lambda_{0} \in \nabla_{-}, \\
\mathfrak{s t a b} b_{\lambda}^{\mathbb{Z}}(y) & =e^{-y \mu} J_{\mu}^{\mathrm{R}} \mathfrak{s t a b} \mathfrak{b}_{\lambda_{1}}^{\mathbb{Z}}(y), \quad y \in W, \mu+\lambda_{1}=\lambda, \mu \in Q, \lambda_{1} \in W \nabla_{-} .
\end{aligned}
$$

As an immediate corollary of Theorems 3.3, 6.2, Eq. (8) and [12, Theorem 1.3.1, Proposition 1.4.3, Theorem 1.6.1], we have the following theorem, which gives a categorification of the stable basis.

Theorem 7.2 [45]Applying the derived tensor $-\otimes_{\mathbb{Z}}^{L} \mathbb{C}$ to $\mathfrak{s t a b}_{\lambda}^{\mathbb{Z}}(w) \in D_{T_{\mathbb{Z}}}^{\mathrm{b}}\left(T^{*} \mathcal{B}_{\mathbb{Z}}\right)$, and taking the class in the Grothendieck group, we get $\mathcal{L}_{-\rho} \otimes \operatorname{stab}_{w}^{\mathfrak{C}_{+}, T \mathcal{B}, \nabla} \in K_{T}$ $\left(T^{*} \mathcal{B}\right)$, where $\nabla$ is the alcove containing $\lambda$.

### 7.2 Verma Modules in Positive Characteristic

In this section, we briefly sketch the relation between the K-theory stable bases and Verma modules for Lie algebras over positive characteristic fields. We refer the readers to [45, Sect. 9] for the details.

We consider the level- $p$ configuration of $\rho$-shifted affine hyperplanes, that is, $H_{\alpha^{\vee}, n}^{p}=\left\{\mu \in \Lambda_{\mathbb{Q}} \mid\left\langle\alpha^{\vee}, \mu+\rho\right\rangle=n p\right\}$. Let $A_{0}$ be the fundamental alcove, i.e., it contains $(\epsilon-1) \rho$ for small $\epsilon>0$. Let $W$ act on $\Lambda$ via the level- $p$ dot action, $w: \mu \mapsto w \bullet \mu=w(\mu+\rho)-\rho$.

Let $k$ be an algebraically closed field of characteristic $p$, and $p$ is greater than the Coxeter number. For any $k$-variety $X$, let $X^{(1)}$ be the Frobenius twist. Let $U\left(\mathfrak{g}_{k}\right)$ be the universal enveloping algebra of $\mathfrak{g}_{k}$ with the Frobenius center $\mathcal{O}\left(\mathfrak{g}_{k}^{*(1)}\right)$ and the Harish-Chandra center $\mathcal{O}\left(\mathfrak{h}_{k}^{*} /(W, \bullet)\right)$. Let $\lambda \in \mathfrak{h}_{k}^{*}$ be regular (i.e., does not lie on any hyperplane $H_{\alpha^{\vee}, n}^{p}$ ) and integral (i.e., belongs to the image of the derivative $d$ : $\Lambda \rightarrow \mathfrak{h}_{k}^{*}$ ). Let $U\left(\mathfrak{g}_{k}\right)^{\lambda}$ be the quotient of $U\left(\mathfrak{g}_{k}\right)$ by the central ideal corresponding to $W \bullet \lambda \in \mathfrak{h}_{k}^{*} /(W, \bullet)$. Let $\operatorname{Mod}_{\chi}\left(U\left(\mathfrak{g}_{k}\right)^{\lambda}\right)$ be the category of finitely generated $U\left(\mathfrak{g}_{k}\right)^{\lambda}-$ modules on which the Frobenius center $\mathcal{O}\left(\mathfrak{g}_{k}^{*(1)}\right)$ acts by the generalized character $\chi \in \mathfrak{g}_{k}^{*(1)}$.

Let $\mathcal{D}^{\lambda}$ be the ring of $\mathcal{L}_{\lambda}$-twisted differential operators on $\mathcal{B}$, and $D^{\mathrm{b}}\left(\operatorname{Coh}_{\chi} \mathcal{D}^{\lambda}\right)$ be the full subcategory of coherent $\mathcal{D}^{\lambda}$-modules that are set-theoretically supported on $\mathcal{B}_{\chi}^{(1)}$ as a coherent sheaf on $T^{*} \mathcal{B}_{k}^{(1)}$, where $\mathcal{B}_{\chi}^{(1)}$ is the Springer fiber. Then the global section functor $R \Gamma_{\mathcal{D}^{\lambda}, \chi}: D^{\mathrm{b}}\left(\operatorname{Coh}_{\chi} \mathcal{D}^{\lambda}\right) \rightarrow D^{\mathrm{b}}\left(\operatorname{Mod}_{\chi} U\left(\mathfrak{g}_{k}\right)^{\lambda}\right)$ is an equivalence according to [10, Theorem 3.2], whose inverse is denoted by $\mathfrak{L}^{\lambda_{0}}$. Let $T^{*} \mathcal{B}_{\chi}^{(1) \wedge}$ be the completion of $T^{*} \mathcal{B}^{(1)}$ at $\mathcal{B}_{\chi}^{(1)}$. Then for all integral $\lambda \in \mathfrak{h}^{*}$, the Azumaya algebra $\mathcal{D}^{\lambda}$ splits on $T^{*} \mathcal{B}_{\chi}^{(1) \wedge}[10$, Theorem 5.1.1]. In particular, there is a Morita equivalence

$$
\operatorname{Coh}_{\mathcal{B}_{\chi}^{(1)}}\left(T^{*} \mathcal{B}^{(1)}\right) \cong \operatorname{Coh}_{\chi} \mathcal{D}^{\lambda} .
$$

The equivalence above depends on the choice of a splitting bundle of $\mathcal{D}^{\lambda}$ [10, Remark 5.2.2]. As has been done in [8,12], for $\lambda_{0} \in A_{0}$, we normalize the choice of the splitting bundle by [9, Remark 1.3.5]. In this section, we denote this splitting bundle by $\mathcal{E}^{s}$. That is, $\left.\mathcal{E} n d_{T^{*} \mathcal{B}_{\chi}^{(1) \wedge}}\left(\mathcal{E}^{s}\right) \cong \mathcal{D}^{\lambda_{0}}\right|_{T^{*} \mathcal{B}_{\chi}^{(1) \wedge}}$. This bundle fixes the equivalence $\operatorname{Coh}_{\mathcal{B}_{\chi}^{(1)}}\left(T^{*} \mathcal{B}^{(1)}\right) \cong \operatorname{Coh}_{\chi} \mathcal{D}^{\lambda_{0}}$. Composing this equivalence with $\mathfrak{L}^{\lambda_{0}}$, we have

$$
\begin{equation*}
\gamma_{\chi}^{\lambda_{0}}: D^{\mathrm{b}} \operatorname{Mod}_{\chi}\left(U\left(\mathfrak{g}_{k}\right)^{\lambda_{0}}\right) \cong D^{\mathrm{b}} \operatorname{Coh}_{\mathcal{B}_{\chi}^{(1)}}\left(T^{*} \mathcal{B}^{(1)}\right) \tag{12}
\end{equation*}
$$

which is also referred to as the localization equivalence.
Let $U\left(\mathfrak{g}_{k}\right)_{\chi}^{\lambda_{0}}$ be the completion of $U\left(\mathfrak{g}_{k}\right)^{\lambda_{0}}$ at the central character $\chi \in \mathcal{N}_{k}^{(1)}$. Then, $E n d_{T^{*} \mathcal{B}_{\chi}^{(1) \wedge}}\left(\mathcal{E}^{s}\right) \cong U\left(\mathfrak{g}_{k}\right)_{\chi}^{\lambda_{0}}$. The localization functor $\gamma_{\chi}^{\lambda_{0}}$ can then be written as $\otimes_{U\left(\mathfrak{g}_{k}\right)^{\lambda_{0}}} \mathcal{E}^{s}$. We also have the completed version of the equivalence [12, Remark 2.5.5]

$$
\gamma_{\chi}^{\lambda_{0}}: D^{\mathrm{b}} \operatorname{Mod}\left(U\left(\mathfrak{g}_{k}\right)_{\chi}^{\lambda_{0}}\right) \cong D^{\mathrm{b}} \operatorname{Coh}\left(T^{*} \mathcal{B}_{\chi}^{(1) \wedge}\right)
$$

From now on we consider the case $\chi=0$. In this case, we abbreviate $\gamma_{\chi}^{\lambda_{0}}$ simply as $\gamma^{\lambda_{0}}$, and the twisted Springer fiber $\mathcal{B}_{\chi}^{(1)}$ is the zero section $\mathcal{B}^{(1)} \subseteq T^{*} \mathcal{B}^{(1)}$.

The bundle $\mathcal{E}^{s}$ on $T^{*} \mathcal{B}_{0}^{(1) \wedge}$ has a natural $T_{k}$-equivariant structure [8, § 5.2.4], where $T_{k}=A_{k} \times\left(\mathbb{G}_{m}\right)_{k}$. Taking the ring of endomorphisms, we get a $T_{k}$-action on $U\left(\mathfrak{g}_{k}\right)_{0}^{\lambda_{0}}$ compatible with that on $\mathcal{N}^{(1)}$. In particular, the $\left(\mathbb{G}_{m}\right)_{k} \subseteq T_{k}$-action provides a nonnegative grading on $U\left(\mathfrak{g}_{k}\right)_{0}^{\lambda_{0}}$, referred to as the Koszul grading. The localization equivalences above can be made into equivalences of equivariant categories. Let $\operatorname{Mod}_{0}^{\mathrm{gr}}\left(U\left(\mathfrak{g}_{k}\right)^{\lambda_{0}}, A_{k}\right)$ be the category of finite-dimensional Koszul-graded modules of $U\left(\mathfrak{g}_{k}\right)_{0}^{\lambda_{0}}$, which are endowed with compatible actions of the subgroup $A_{k} \subseteq G_{k}$. The compatibility is in the sense that the action of $A_{k}$ differentiates to the action of the subalgebra $\mathfrak{h}_{k} \subseteq \mathfrak{g}_{k}$. Then, we have [8, Theorem 1.6.7]

$$
\gamma^{\lambda_{0}}: D^{\mathrm{b}} \operatorname{Mod}_{0}^{\mathrm{gr}}\left(U\left(\mathfrak{g}_{k}\right)^{\lambda_{0}}, A_{k}\right) \cong D_{T_{k}}^{\mathrm{b}} \operatorname{Coh}_{\mathcal{B}_{k}^{(1)}}\left(T^{*} \mathcal{B}_{k}^{(1)}\right)
$$

and the completed version

$$
\gamma^{\lambda_{0}}: D^{\mathrm{b}} \operatorname{Mod}^{\mathrm{gr}}\left(U\left(\mathfrak{g}_{k}\right)_{0}^{\lambda_{0}}, A_{k}\right) \cong D_{T_{k}}^{\mathrm{b}} \operatorname{Coh}\left(T^{*} \mathcal{B}_{0}^{(1) \wedge}\right)
$$

By the functor of taking finite vectors, we get [8, Theorem 5.1.1]

$$
\gamma^{\lambda_{0}}: D^{\mathrm{b}} \operatorname{Mod}^{\mathrm{gr}}\left(\mathcal{A}^{\lambda_{0}}, A_{k}\right) \cong D_{T_{k}}^{\mathrm{b}} \operatorname{Coh}\left(T^{*} \mathcal{B}_{k}^{(1)}\right)
$$

 it has the property that $\left(\mathcal{A}^{\lambda_{0}}\right)_{0}^{\wedge} \cong U\left(\mathfrak{g}_{k}\right)_{0}^{\lambda_{0}}$. Using the correspondence given by taking completion and taking finite vectors, we will freely pass between $\operatorname{Mod}^{\text {gr }}\left(U\left(\mathfrak{g}_{k}\right)_{0}^{\lambda_{0}}, A_{k}\right)$ and $\operatorname{Mod}^{\mathrm{gr}}\left(\mathcal{A}^{\lambda_{0}}, A_{k}\right)$; similarly for $D_{T_{k}}^{\mathrm{b}} \operatorname{Coh}\left(T^{*} \mathcal{B}_{0}^{(1) \wedge}\right)$ and $D_{T_{k}}^{\mathrm{b}} \operatorname{Coh}\left(T^{*} \mathcal{B}_{k}^{(1)}\right)$.

For any $\lambda$ in the $W_{\text {aff }}^{\prime}$-orbit of $\lambda_{0}$, we can define the localization functor

$$
\gamma^{\lambda}: D^{\mathrm{b}} \operatorname{Mod}^{\mathrm{gr}}\left(\mathcal{A}^{\lambda_{0}}, A_{k}\right) \cong D_{T_{k}}^{\mathrm{b}} \operatorname{Coh}\left(T^{*} \mathcal{B}_{k}^{(1)}\right)
$$

by precomposing with the affine braid group action functors [45, Sect. 9.3].
For the Lie algebra $\mathfrak{b}_{k}$, recall the Verma module $Z^{\mathfrak{b}}(\lambda):=U(\mathfrak{g}) \otimes_{U(\mathfrak{b})} k_{\lambda}$. Then we have

Theorem 7.3 ([45]) Let $k$ be an algebraically closed field of characteristic $p$, and $p$ is greater than the Coxeter number. Assume $\lambda$ to be regular and integral, then in $D_{T_{k}}^{\mathrm{b}}\left(T^{*} \mathcal{B}_{k}^{(1)}\right)$, we have isomorphisms

$$
e^{\rho} \mathfrak{s t a b}_{-\frac{\lambda+\rho}{\rho}}^{k}(y) \cong \gamma^{\lambda} Z^{\mathfrak{b}}(y \bullet \lambda+2 \rho),
$$

where $\mathfrak{s t a b}_{-\frac{\lambda+\rho}{p}}^{k}(y)=\mathfrak{s t a b}_{-\frac{\lambda+\rho}{p}}^{\mathbb{Z}}(y) \otimes_{\mathbb{Z}}^{L} k$.
To prove this theorem, one needs to use the affine braid group action on $\operatorname{Mod}_{\chi} U\left(\mathfrak{g}_{k}\right)^{\lambda}$ constructed in [9], which iteratively produces all the Verma modules [25], and also its compatibility with the localization equivalences $\gamma^{\lambda}$. Together with the iterative definition of $\mathfrak{s t a b}{ }_{\mathbb{Z}}^{k}(y)$, the theorem follows.

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# Characteristic Classes of Orbit Stratifications, the Axiomatic Approach 

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#### Abstract

Consider a complex algebraic group $G$ acting on a smooth variety $M$ with finitely many orbits, and let $\Omega$ be an orbit. The following three invariants of $\Omega \subset M$ can be characterized axiomatically: (1) the equivariant fundamental class $[\bar{\Omega}, M] \in H_{G}^{*}(M)$, (2) the equivariant Chern-Schwartz-MacPherson class $\mathrm{c}^{\mathrm{sm}}(\Omega, M) \in H_{G}^{*}(M)$, and (3) the equivariant motivic Chern class $\mathrm{mC}(\Omega, M) \in$ $K_{G}(M)[y]$. The axioms for Chern-Schwartz-MacPherson and motivic Chern classes are motivated by the axioms for cohomological and K-theoretic stable envelopes of Okounkov and his coauthors. For $M$ a flag variety and $\Omega$ a Schubert cell-an orbit of the Borel group acting-this implies that CSM and MC classes coincide with the weight functions studied by Rimányi-Tarasov-Varchenko. In this paper we review the general theory and illustrate it with examples.


Keywords Equivariant characteristic classes • Chern-Schwartz-MacPherson class • Motivic Chern class • Schubert classes • Axiomatic characterisation

## 1 Introduction

An effective way of studying the subvariety $X$ of a smooth complex variety $M$ is assigning a characteristic class to $X \subset M$, living in the cohomology or K-theory of the ambient space $M$. When $M$ is a $G$-representation and $X$ is invariant, then the

[^20]characteristic class can be considered in the richer $G$-equivariant cohomology or K-theory of $M$.

In this paper we consider three flavors of characteristic classes: the fundamental class (in cohomology), the Chern-Schwartz-MacPherson (CSM) class (in cohomology) and the motivic Chern class (MC class) in K-theory. These characteristic classes encode fine geometric and enumerative properties of $X \subset M$, and their theory overlaps with the recent advances relating geometry with quantum integrable systems.

We will be concerned with the following situation: a complex algebraic group $G$ acts on a smooth variety $M$ with finitely many orbits, and we want to find the characteristic classes listed above associated to the orbits or their closures. This situation is frequent in Schubert calculus and other branches of enumerative geometry. By the nature how characteristic classes are defined, such calculation assumes resolutions of the singularities of the orbit closures. The main message of the theorems we review in Sects. 2, 3, 4, is that such a resolution calculation can be replaced with an axiomatic approach for all three types of classes mentioned above. For the fundamental class this fact has been known for two decades [16, 29], but for the other two mentioned classes this is a very recent development stemming from the notion of Okounkov's stable envelopes [23, 27].

As an application of the axiomatic approach the characteristic classes of ordinary and matrix Schubert cells can be computed (in type A, for arbitrary partial flag manifolds)-we present the formulas in Sect. 5. For the motivic Chern class one of the axioms concerns convex geometry, namely the containment of one convex polytope in another one. Throughout our exposition we try to illustrate this fascinating connection between complex geometry and the theory of convex polytopes with 2and 3-dimensional pictures.

It is a classical fact that characteristic classes of geometrically relevant varieties often exhibit positivity (or alternating sign) properties. Having computed lots of examples of equivariant motivic Chern classes, we observe these properties for the motivic Chern class as well. In Sect. 6 we collect three different positivity conjectures on MC classes of Schubert cells. One of these conjectures is a K-theoretic counterpart of properties of CSM classes studied by Aluffi-Mihalcea-Schürmann-Su [2, 3]. Another one of our conjectures is a new phenomenon for characteristic classes: not the sign of some coefficients are conjectured, but rather that the coefficient polynomials are log-concave, cf. [21].

We finish the paper with studying the characteristic classes of varieties in the source and in the target of a GIT map. The fundamental class of an invariant subvariety and the fundamental class of its image under the GIT map are essentially the same (more precisely, one lives in a quotient ring, and the other one is a representative of it). A familiar fact illustrating this phenomenon is that the fundamental class of both the matrix Schubert varieties and the ordinary Schubert varieties are the Schur functions. However, for the CSM and the MC classes the situation is different: these classes before and after the GIT map are not the same (or rather, one does not represent the other). Yet, there is direct relation between them, which we prove in Sect. 8.

The goal of most of this paper is to review and illustrate concepts, although in Sect. 4.4 we complement a proof of [19] with an alternative argument. The mathe-
matical novelties are the conjectures in Sect. 6, the careful treatment of the pull-back property of motivic Segre classes in Sect. 8.1, and its consequence, Theorem 8.12 which relates motivic Segre classes of varieties in the source and in the target of a GIT quotient.

## 2 Fundamental Class in Equivariant Cohomology

By equivariant cohomology we will mean Borel's version, developed in [8]. Our general reference is [9].

Suppose a complex algebraic group $G$ acts on the smooth complex algebraic variety $M$. Suppose also that the action has finitely many orbits, and let $\mathcal{O}$ be the set of orbits. For an orbit $\Omega \in \mathcal{O}$ let $G_{\Omega}$ denote the stabilizer subgroup of a point in $\Omega$. Then additively we have

$$
H_{G}^{*}(M ; \mathbb{Q}) \simeq \bigoplus_{\Omega \in \mathcal{O}} H_{G}^{*-2 \operatorname{codim}(\Omega)}(\Omega ; \mathbb{Q}) .
$$

The decomposition holds because each cohomology group $H_{G}^{*}(\Omega ; \mathbb{Q}) \simeq H^{*}\left(B G_{\Omega} ; \mathbb{Q}\right)$ is concentrated in even degrees. Suppose further that the normal bundle $\nu_{\Omega}$ of each orbit $\Omega$ has nonzero equivariant Euler class

$$
\begin{equation*}
e\left(\nu_{\Omega}\right) \neq 0 \in H_{G}^{*}(\Omega ; \mathbb{Q}) . \tag{1}
\end{equation*}
$$

Then a class in $H_{G}^{*}(M ; \mathbb{Q})$ is determined by the restrictions to orbits, cf. [6, Sect. 9]. The topic of this paper is that certain characteristic classes associated with orbits are determined by less information: they are determined by some properties of their restrictions to the orbits.

Example 2.1 For $M=\mathbb{C}, G=\mathbb{C}^{*}$ with the natural action on $\mathbb{C}$. We have the short exact sequence

$$
\begin{array}{ccc}
0 \longrightarrow H_{\mathbb{C}^{*}}^{*-2}(\{0\}, \mathbb{Q}) \underset{t}{\longrightarrow} & H_{\mathbb{C}^{*}}^{*} \\
\| & \| \\
\mathbb{Q}[t] & \mathbb{Q}) \underset{e v_{t=0}}{\longrightarrow} & H_{\mathbb{C}^{*}}^{*}(\mathbb{C} \\
\| & \| \\
\mathbb{Q}[t] & \mathbb{Q}) \longrightarrow 0 .
\end{array}
$$

In this example the restriction to $\{0\}$ is obviously an isomorphism. In general, when $M$ is an equivariantly formal space (see [20]) and $G$ is a torus, then the restriction to the fixed point set is a monomorphism.

The most natural characteristic class associated to an orbit $\Omega$ is the fundamental class of its closure, which we will denote by $[\bar{\Omega}]=[\bar{\Omega}, M] \in H_{G}^{*}(M, \mathbb{Q})$. In an appropriate sense, it is the Poincaré dual of the homology fundamental class of $\bar{\Omega}$. Such classes are studied in Schubert calculus under the name of Schubert classes, in singularity theory under the name of Thom polynomials, and in the theory of quivers
under the name of quiver polynomials. Here is an axiomatic characterization of the equivariant fundamental class, which leads to an effective method to calculate them.

Theorem 2.2 ([16, 29]) Suppose $M$ has finitely many $G$-orbits and (1) holds. Then the equivariant fundamental classes of orbit closures are determined by the conditions:
(i) (support condition) the class $[\bar{\Omega}]$ is supported on $\bar{\Omega}$;
(ii) (normalization condition) $[\bar{\Omega}]_{\left.\right|_{\Omega}}=e\left(\nu_{\Omega}\right)$;
(iii) (degree condition) $\operatorname{deg}([\bar{\Omega}])=2 \operatorname{codim}(\Omega)$.

We obtain an equivalent system of conditions if we replace condition (i) with any of the following two:
(iv) (homogeneous equations, verl) $[\bar{\Omega}]_{\Theta}=0$, if $\Theta \not \subset \bar{\Omega}$,
(v) (homogeneous equations, ver 2 ) $[\bar{\Omega}]_{\mid \Theta}=0$, if $\operatorname{codim}(\Theta) \leq \operatorname{codim}(\Omega), \Theta \neq \Omega$.

We also obtain an equivalent system of conditions if we replace condition (iii) with
(vi) (modified degree condition) $\operatorname{deg}([\Omega])_{\Theta \Theta}<\operatorname{deg}([\Theta])_{\Theta \Theta}$ for $\Theta \neq \Omega$.

The advantage of (iv) or (v) over $(i)$ is that they are local conditions, and hence explicitly computable. The modified degree condition is only added to illustrate the similarity with analogous characterizations of other characteristic classes in Sects.3, 4.

Example 2.3 Let $G=\mathrm{GL}_{2}(\mathbb{C}) \times \mathrm{GL}_{3}(\mathbb{C})$ act on the vector space $\operatorname{Hom}\left(\mathbb{C}^{2}, \mathbb{C}^{3}\right)$ via $(A, B) \cdot X=B X A^{-1}$. The orbits of this action are $\Omega^{i}=\{X: \operatorname{dim}(\operatorname{ker}(X))=i\}$ for $i=0,1,2$. The fundamental class $\left[\overline{\Omega^{1}}\right]$ is a homogeneous degree 2 polynomial in $\mathbb{Z}\left[a_{1}, a_{2}, b_{1}, b_{2}, b_{3}\right]^{S_{2} \times S_{3}}$ (where $\operatorname{deg} a_{i}=\operatorname{deg} b_{i}=1$ ). The constraints put on this polynomial by the conditions of Theorem 2.2 are $\phi_{0}\left(\left[\overline{\Omega^{1}}\right]\right)=0$ (condition (v)) and $\phi_{1}\left(\left[\overline{\Omega^{1}}\right]\right)=\left(b_{2}-a_{2}\right)\left(b_{3}-a_{2}\right)$ (condition (ii)), where $\phi_{0}$ is the "restriction to $\Omega^{0}$ " map, namely $a_{1} \mapsto t_{1}, a_{2} \mapsto t_{2}, b_{1} \mapsto t_{1}, b_{2} \mapsto t_{2}, b_{3} \mapsto b_{3}$, and $\phi_{1}$ is the "restriction to $\Omega^{1 "}$ map, namely $a_{1} \mapsto t_{1}, a_{2} \mapsto a_{2}, b_{1} \mapsto t_{1}, b_{2} \mapsto b_{2}, b_{3} \mapsto b_{3}$. Calculation shows that the only solution to these two constraints is $\left[\overline{\Omega^{1}}\right]=A_{1}^{2}-A_{2}+B_{2}-$ $A_{1} B_{1}$, where $A_{i}$ 's ( $B_{i}$ 's) are the elementary symmetric polynomials of the $a_{i}$ 's $\left(b_{i}\right.$ 's).

More details on this calculation as well as on analogous calculations for 'similar' representations (e.g. Dynkin quivers, $\Lambda^{2} \mathbb{C}^{n}, S^{2} \mathbb{C}^{n}$ ) can be found e.g. in [15-17]. However, these fundamental classes can also be computed by other methods (resolutions, degenerations). The real power of Theorem 2.2 is that it is applicable even in situations where deeper geometric information such as resolutions or degenerations are not known-e.g. contact singularities, matroid representations spaces, see [14, 29].

## 3 Chern-Schwartz-MacPherson Classes of Orbits

Another cohomological characteristic class associated with an invariant subset $X$ of $M$ is the Chern-Schwartz-MacPherson (CSM) class $\mathrm{c}^{\mathrm{sm}}(X, M) \in H_{G}^{*}(M ; \mathbb{Q})$, see e.g. [2, 22, 25, 39] for foundational literature and [18] for the version we consider. The CSM class is an inhomogeneous class, whose lowest degree component equals the fundamental class $[\bar{X}, M$ ].

In general CSM classes are defined for (equivariant) constructible functions on $M$. Here we will only consider CSM classes of locally closed invariant smooth subvarieties of $M$ (corresponding to the indicator functions of such varieties).

According to Aluffi [1] the CSM classes can be calculated along the following lines: one finds $Y$ a partial completion of $X$ which maps properly to $M$ and such that $Y \backslash X$ is a smooth divisor with simple normal crossings $D=\bigcup_{i=1}^{m} D_{i}$. That is, we have a diagram

where $f$ is a proper map. Then

$$
\begin{equation*}
\mathrm{c}^{\mathrm{sm}}(X, M)=f_{*}\left(c\left(T_{Y}\right)-\sum_{i} c\left(T_{D_{i}}\right)+\sum_{i<j} c\left(T_{D_{i} \cap D_{j}}\right)-\cdots\right) . \tag{2}
\end{equation*}
$$

In our situation, this method demanding knowledge about a resolution, can be replaced with an axiomatic characterization.

For an orbit $\Omega \subset M$ let $x_{\Omega} \in \Omega$, and let $G_{\Omega}$ be the stabilizer subgroup of $x_{\Omega}$. Denote $T_{\Omega}=T_{x_{\Omega}} \Omega$, and $\nu_{\Omega}=T_{x_{\Omega}} M-T_{\Omega}$ as $G_{\Omega}$-representations. By the degree of an inhomogeneous cohomology class $a=a_{0}+a_{1}+\cdots+a_{d}$ with $a_{i} \in H_{G}^{i}(\Omega)$ and $a_{d} \neq 0$ we mean $d$.

Theorem 3.1 ([18, 34]) Suppose M has finitely many G-orbits and (1) holds. Then the equivariant CSM classes of the orbits are determined by the conditions:
( $i$ ') (divisibility condition) for any $\Theta$ the restriction $\mathrm{c}^{\mathrm{sm}}(\Omega, M)_{\Theta \Theta}$ is divisible by $c\left(T_{\Theta}\right)$ in $H_{G}^{*}(\Theta ; \mathbb{Q})$;
(ii') (normalization condition) $\mathrm{c}^{\mathrm{sm}}(\Omega, M)_{\mid \Omega}=e\left(\nu_{\Omega}\right) c\left(T_{\Omega}\right)$;
(iii') (smallness condition) $\operatorname{deg}\left(\mathrm{c}^{\mathrm{sm}}(\Omega, M)_{\mid \Theta}\right)<\operatorname{deg}\left(\mathrm{c}^{\mathrm{sm}}(\Theta, M)_{\mid \Theta}\right)$ for $\Theta \neq \Omega$.

Example 3.2 Continuing Example 2.3, let us calculate $c^{\mathrm{sm}}\left(\Omega^{1}\right)$. It is an inhomogeneous polynomial in $\mathbb{Z}\left[a_{1}, a_{2}, b_{1}, b_{2}, b_{3}\right]^{S_{2} \times S_{3}}$. The constraints put on this polynomial by the conditions of Theorem 3.1 are $\phi_{0}\left(\mathrm{c}^{\mathrm{sm}}\left(\Omega^{1}\right)\right)=0, \phi_{1}\left(\mathrm{c}^{\mathrm{sm}}\left(\Omega^{1}\right)\right)=$ $\left(b_{2}-a_{2}\right)\left(b_{3}-a_{2}\right)\left(1+b_{2}-t_{1}\right)\left(1+b_{3}-t_{1}\right)\left(1+t_{1}-a_{2}\right)$, and $\operatorname{deg}\left(\phi_{2}\left(\mathrm{c}^{\mathrm{sm}}\left(\Omega^{1}\right)\right)\right)<$
6.

Here $\phi_{1}, \phi_{2}$ are as in Example 2.3, and $\phi_{3}$ is the "restriction to $\Omega^{2 "}$ map which turns out to be the identity map. Calculation shows that the unique solution to these constraints is (see notations in Example 2.3)

$$
\begin{aligned}
c^{\mathrm{sm}}\left(\Omega^{1}\right)=\left(A_{1}-A_{2}+B_{2}-A_{1} B_{1}\right)+\left(-A_{1}^{3}+2 A_{1}^{2} B_{1}-A_{1} B_{1}^{2}\right. & \left.-A_{1} B_{2}+B_{1} B_{1}-B_{3}\right)+ \\
\cdots & +\left(-3 A_{1} A_{2}^{2}+\cdots+2 B_{2} B_{3}\right) .
\end{aligned}
$$

Details on such calculations, other examples of interpolation calculations of CSM classes, as well as more conceptual ways of presenting the long CSM polynomials are in [18, 28].

## 4 Characteristic Classes in Equivariant K-Theory

### 4.1 The Motivic Chern Class

We consider the topological equivariant K-theory constructed by Segal [36] or Thomason's algebraic K-theory, see [13, Sect. 5]. It is compatible with the K-theory built from locally free sheaves for algebraic varieties. As before we have a decomposition

$$
K_{G}(M) \simeq \bigoplus_{\Omega \in \mathcal{O}} K_{G}(\Omega)
$$

and $K_{G}(\Omega)$ is isomorphic with the representation ring $R\left(G_{\Omega}\right)$.
The analogue of the Euler class of a vector bundle $E$ in K-theory is provided by the $\lambda$-operation

$$
e^{K}(E)=\lambda_{-1}\left(E^{*}\right)=1-E^{*}+\Lambda^{2} E^{*}-\Lambda^{3} E^{*}+\cdots .
$$

It has the property that for any submanifold $\iota: N \hookrightarrow M$ and an element $\beta \in K_{G}(N)$ over $N$ we have

$$
\iota^{*} \iota_{*}(\beta)=e^{K}\left(\nu_{N}\right) \cdot \beta
$$

The total Chern class in K-theory is defined by

$$
c_{G}^{K}(E)=\lambda_{y}\left(E^{*}\right)=1+y E^{*}+y^{2} \Lambda^{2} E^{*}+y^{3} \Lambda^{3} E^{*}+\cdots \in K_{G}(M)[y] .
$$

Note that, both in cohomology and K-theory, the total Chern class can be interpreted as the Euler class of the bundle $E$ which is equivariant with respect to $G \times \mathbb{C}^{*}$, where $\mathbb{C}^{*}$ acts on the base trivially, but acts on the bundle by scalar multiplication. Identifying $K_{G \times \mathbb{C}^{*}}(M)$ with $K_{G}(M)\left[h, h^{-1}\right]$ and setting $y=-h^{-1}$ we obtain

$$
c_{G}^{K}(E)=e_{G \times \mathbb{C}^{*}}^{K}(E)
$$

We will drop the subscript $G$ in the notation, as we do for cohomological Chern classes.

Our next characteristic class defined for a locally closed subvariety $X$ in the smooth ambient space $M$ is the motivic Chern class (MC class) $\mathrm{mC}(X, M) \in K_{G}(M)[y]$, see [10]. In fact this class is defined more generally for maps $X \rightarrow M$ but we will not need this generality here. The $G$-equivariant MC class

$$
\mathrm{mC}(X, M) \in K_{G}(M)[y]
$$

is a straightforward generalization, see details in [4, 19].
A natural method to calculate the motivic Chern class of a locally closed subvariety $X \subset M$ is through the K-theoretic analogue

$$
\begin{equation*}
\mathrm{mC}(X, M)=f_{*}\left(c^{K}\left(T_{Y}\right)-\sum_{i} c^{K}\left(T_{D_{i}}\right)+\sum_{i<j} c^{K}\left(T_{D_{i} \cap D_{j}}\right)-\cdots\right) \tag{3}
\end{equation*}
$$

of the formula (2).

### 4.2 Axiomatic Characterization of MC Classes

In certain situations the resolution method (3) can be replaced with an axiomatic characterization, which we will now explain in several steps.

Theorem 4.1 ([19, Corollary 4.5, Lemma 5.1]) Suppose M has finitely many $G$ orbits. Then the equivariant MC classes of the orbits satisfy the conditions:
( $i$ ") (support condition) the class $\mathrm{mC}(\Omega, M)$ is supported on $\bar{\Omega}$;
(ii") (normalization condition) $\operatorname{mC}(\Omega, M)_{\mid \Omega}=e^{K}\left(\nu_{\Omega}\right) c^{K}\left(T_{\Omega}\right)$;
(iii") (smallness condition) $\mathcal{N}\left(\mathrm{mC}(\Omega, M)_{\mid \Theta}\right) \subseteq \mathcal{N}\left(\mathrm{mC}(\Theta, M)_{\mid \Theta}\right)$ [see explanation below];
(iv") (the divisibility condition) for any $\Theta$ the restriction $\mathrm{mC}(\Omega, M)_{\mid \Theta}$ is divisible by $c^{K}\left(T_{\Theta}\right)$ in $K_{G}(\Theta)[y]$.

In fact, the local condition ( $i v^{\prime \prime}$ ) implies the condition ( $i$ "), and it also implies the obvious condition
$(v ") \mathrm{mC}(\Omega, M)_{\mid \Theta}=0$ if $\Theta \not \subset \bar{\Omega}$.
The K-theoretic version of the condition (1) follows from (iii').
Let us give a precise formulation of the smallness condition (iii"). The classes $\mathrm{mC}(\Omega, M)_{\mid \Theta}$ and $\mathrm{mC}(\Theta, M)_{\mid \Theta}$ belong to $K_{G}(\Theta)[y]=R\left(G_{\Theta}\right)[y]$. Restrict these
classes to the representation ring $R\left(\mathbb{T}_{\Theta}\right)[y]$, where $\mathbb{T}_{\Theta}$ is a maximal torus of $G_{\Theta}$. For a chosen isomorphism $\mathbb{T}_{\Theta}=\left(\mathbb{C}^{*}\right)^{r}$ we have $R\left(\mathbb{T}_{\Theta}\right)[y]=\mathbb{Z}\left[\alpha_{1}^{ \pm 1}, \ldots, \alpha_{r}^{ \pm 1}, y\right]$ where $\alpha_{i}$ are the K-theoretic Chern roots corresponding to the factors of $\left(\mathbb{C}^{*}\right)^{r}$. The Newton polygon of $\beta=\sum_{I} a_{I}(y) z^{I} \in R\left(\mathbb{T}_{\Theta}\right)[y]$ (multiindex notation) is

$$
\mathcal{N}(\beta)=\text { convex hull }\left\{I \in \mathbb{R}^{r}: a_{I}(y) \neq 0\right\}
$$

in particular, for this notion the parameter $y$ is considered a constant, not a variable. The definition of $\mathcal{N}$ depends on the chosen isomorphism between $R\left(\mathbb{T}_{\Theta}\right)[y]$ and $\mathbb{Z}\left[z_{1}^{ \pm 1}, \ldots, z_{r}^{ \pm 1}, y\right]$. Different such choices result in linearly equivalent convex polygons. The smallness condition compares two such Newton polygons, of course the same isomorphism need to be chosen for the two sides.

Example 4.2 If $M=\mathbb{C}, G=\mathbb{C}^{*}$ with the natural action on $\mathbb{C}$ we have a short exact sequence

$$
\begin{gathered}
0 \longrightarrow K_{\mathbb{C}^{*}}(\{0\}) \\
\| \\
\mathbb{Z}\left[\xi, \xi^{-1}\right]
\end{gathered}
$$

In this example the restriction to $\{0\}$ is an isomorphism. We have

$$
\mathrm{mC}\left(\mathbb{C}^{*}, \mathbb{C}\right)=(1+y) \xi^{-1}, \quad \operatorname{mC}(\{0\}, \mathbb{C})=1-\xi^{-1}
$$

and

$$
\mathcal{N}\left(\mathrm{mC}\left(\mathbb{C}^{*}, \mathbb{C}\right)\right)_{\mid 0}=\{-1\}, \quad \mathcal{N}(\mathrm{mC}(\{0\}, \mathbb{C}))_{\mid 0}=[-1,0] .
$$

Example 4.3 Consider the standard action of $G=\mathrm{SL}_{2}(\mathbb{C})$ on $\mathbb{C}^{2}$. Then we have

$$
\mathcal{N}\left(\mathrm{mC}\left(\mathbb{C}^{2} \backslash\{0\}, \mathbb{C}^{2}\right)\right)_{\mid 0}=\mathcal{N}\left(\mathrm{mC}\left(\{0\}, \mathbb{C}^{2}\right)\right)_{\mid 0}=[-1,1]
$$

Examples 4.2, 4.3 show that the inclusion of Newton polygons in (iii") may or may not be strict. However for an interesting class of actions the inclusion is necessarily strict.

Property 4.4 (cf. [19, Definition 4.4]) We say that the action is positive, iffor each orbit $\Omega, x \in \Omega$ there exists a one dimensional torus $\mathbb{C}^{*} \hookrightarrow G_{x}$ which acts on the normal space $\left(\nu_{\Omega}\right)_{x}$ with positive weights.

Note that Property 4.4 implies that 0 is a vertex of the Newton polygon $\mathcal{N}\left(e^{K}\left(\nu_{\Omega}\right)\right)$ because $e^{K}\left(\nu_{\Omega}\right)_{x}=\prod_{i=1}^{\text {codim }(\Omega)}\left(1-\chi_{i}^{-1}\right)$, where $\chi_{i}$ are the weights of $\mathbb{T}_{\Omega}$ acting on $\left(\nu_{\Omega}\right)_{x}$. In particular $e^{K}\left(\nu_{\Omega}\right) \neq 0$. Let " + " denote the Minkowski sum of polygons.

Theorem 4.5 ([19, Theorem 5.3]) Suppose that the action of $G$ on $M$ has Property 4.4. Then the inclusion in (iii") is strict. Moreover, for any pair of different orbits $\Theta, \Omega$ we have
$(v i ") \mathcal{N}\left(\mathrm{mC}(\Omega, M)_{\mid \Theta}\right) \subseteq \mathcal{N}\left(e^{K}\left(\nu_{\Theta}\right)-1\right)+\mathcal{N}\left(c^{K}\left(T_{\Theta}\right)\right) \subsetneq \mathcal{N}\left(\mathrm{mC}(\Theta, M)_{\mid \Theta}\right)$.

A reformulation of condition (vi") is that $\mathcal{N}\left(\mathrm{mC}(\Omega, M)_{\mid \Theta}\right)$ is contained in $\mathcal{N}\left(\mathrm{mC}(\Theta, M)_{\mid \Theta}\right)$ in such a way that the origin is a vertex of $\mathcal{N}\left(\mathrm{mC}(\Theta, M)_{\mid \Theta}\right)$, but it is not contained in $\mathcal{N}\left(\mathrm{mC}(\Omega, M)_{\mid \Theta}\right)$.

Additionally assuming that the stabilizers $G_{\Omega}$ are connected guarantees that an element $\beta \in K_{\mathbb{T}}(M)$ is determined by the restrictions to orbits, and we obtain

Theorem 4.6 ([19, Theorem 5.5]) If the action of $G$ on $M$ has finitely many orbits, it is positive, and the stabilizers are connected, then the conditions (ii'), (iv"), (vi") determine $\mathrm{mC}(\Omega, M)$ uniquely.

Example 4.7 The Borel subgroup $B_{n}$ of $\mathrm{GL}_{n}(\mathbb{C})$ acts on the full flag variety $\mathrm{Fl}(4)$ with finitely many orbits. The Newton polygon containment (ie. smallness condition (iii")) for two of the orbits, $\Omega=\Omega_{(\{3\},\{4\},\{1\},\{2\})}$ and $\Theta=\Omega_{(\{3\},\{4\},\{2\},\{1\})}$, is illustrated below. For more details about the orbits of this action, their combinatorial codes, and their MC classes see Sect. 7. The description of the Borel-equivariant K-theory of flag varieties can be found in [13, Sect. 6] or in [38]. The figure below illustrates (Newton polygons of) the two MC classes, restricted to $K_{B_{n}}(\Theta)[y]=K_{\mathbb{T}_{\Theta}}(\mathrm{pt})$ [ $y$ ]. It turns out that $\mathbb{T}_{\Theta}=\left(\mathbb{C}^{*}\right)^{4}$ and hence the two Newton polygons live in $\mathbb{R}^{4}$. Moreover, both Newton polygons turn out to be contained in a 3-dimensional subspace (sum of coordinates $=0$ ). Hence the picture shows 3-dimensional convex polytopes.


In blue-Newton polygon of $e^{K}\left(\nu_{\Theta}\right)=\lambda_{-1}\left(\nu_{\Theta}^{*}\right)$.
Arrow in red-cotangent weight $T_{\Theta}^{*}$.
In solid violet-Newton polygon of $\mathrm{mC}\left(\Omega \cap S_{\Theta}\right)$, where $S_{\Theta}$ is a slice to $\Theta$. Edges in violet-Newton polygon of $\mathrm{mC}(\Omega)$.

Remark 4.8 Our smallness conditions (iii") and (vi") are motivated by the analogous smallness condition for $K$-theoretic stable envelopes invented by Okounkov, see
[27, (9.1.10)]. There is, however, a difference. In Okounkov's smallness condition the strictness of the inclusion of Newton polygons is guaranteed by the fact that the small Newton polygon can be shifted slightly within the large Newton polygon, see the first picture below. In fact, in the terminology of [27], the inclusion holds for an open set (interior of an alcove) of "fractional shifts". However, such a "wiggle room" for the inclusions does not necessarily exist for MC classes, see the second and third pictures below. The second picture illustrates the smallness condition (iii") (or $\left(v i\right.$ ") ) for the standard representation $\mathrm{GL}_{2}(\mathbb{C})$ on $\mathbb{C}^{2}$ with $\Omega=\mathbb{C}^{2}-\{0\}$, $\Theta=\{0\}$. The third picture illustrates smallness condition (iii") (or (vi")) for the natural $\mathrm{GL}_{2}(\mathbb{C}) \times \mathrm{GL}_{2}(\mathbb{C})$-representation on $\operatorname{Hom}\left(\mathbb{C}^{2}, \mathbb{C}^{2}\right)$ (a.k.a. $A_{2}$ quiver representation with dimension vector $(2,2)$ ) with $\Omega=\{$ rank 1 maps $\}$ and, $\Theta=\{0\}$.


### 4.3 The Key Idea of the Proof of Theorem 4.1

At the heart of the arguments proving the statements of Sect. 4.2 is showing that conditions (iii") and (vi") hold. The crucial step in the proof is the special case when $\Theta$ is a point. To test an inclusion of Newton polygons it is enough to restrict an action to each one dimensional torus $\mathbb{T}_{0}=\mathbb{C}^{*} \hookrightarrow \mathbb{T}$, with $K_{\mathbb{T}_{0}}(p t)=\mathbb{Z}\left[\xi, \xi^{-1}\right]$. The argument reduces to examining the limit

$$
\begin{equation*}
\lim _{\xi \rightarrow \infty} \frac{\mathrm{mC}(\Omega, M)_{\Theta}}{e^{K}\left(\nu_{\Theta}\right)} . \tag{4}
\end{equation*}
$$

We prove that the limit is equal to

$$
\chi_{y}\left(\Omega \cap M_{\Theta}^{-}\right),
$$

where $M_{\Theta}^{-}$is the Białynicki-Birula minus-cell associated to $\Theta$, (see [7] or [12, Sect. 4.1])

$$
\begin{equation*}
M_{\Theta}^{-}=\left\{x \in M \mid \lim _{t \rightarrow \infty} t x \in \Theta\right\}, \tag{5}
\end{equation*}
$$

and $\chi_{y}$ is the Hirzebruch $\chi_{y}$-genus. The proof of this statement in full generality is given in [41] and in the particular cases needed in the proof of Theorem 4.6 in [19, Theorem 5.3].

Since the limit (4) is finite the degree of the denominator is at least the degree of the numerator. This observation leads to the proof of (iii"). Moreover, if the action is positive then the limit is equal to $\chi_{y}(\emptyset)=0$, thus the degree of the denominator must be strictly larger than the degree of the numerator. This observation leads to the proof of ( $v i$ ").

### 4.4 Strict Inclusion for Homogeneous Singularities

In this section we give a rigorous alternative argument proving the containment property (iii") of Newton polygons of mC classes. We believe that it sheds more light on this intriguing connection between algebraic and convex geometry.

Suppose $\mathbb{T}=\mathbb{C}^{*}$ acts on a vector space $M=\mathbb{C}^{n}$ via scalar multiplication. Let $\Omega \subset$ $\mathbb{C}^{n}$ be an invariant subvariety, where $0 \notin \Omega$. Denote by $Z \subset \mathbb{P}^{n}$ the projectivization of $\Omega$. We have a diagram

where $\widetilde{\mathbb{C}^{n}}$ is the blow-up of $\mathbb{C}^{n}$ at 0 and $\iota: \mathbb{P}^{n-1} \rightarrow \widetilde{\mathbb{C}^{n}}$ is the inclusion of the special fiber, $\Omega^{\prime}=q^{-1}(Z) \backslash \iota(Z)$. Let $h=[\mathcal{O}(1)] \in K\left(\mathbb{P}^{n-1}\right)$. Then

$$
\frac{\mathrm{mC}\left(\Omega^{\prime}, \widetilde{\mathbb{C}^{n}}\right)_{\mathbb{P}^{n-1}}}{\lambda_{-1}\left(\nu_{\mathbb{P}^{n-1} / \widetilde{\mathbb{C}}^{n}}^{*}\right)}=\left(\frac{1+y \xi^{-1} h}{1-\xi^{-1} h}-1\right) \mathrm{mC}\left(Z, \mathbb{P}^{n-1}\right)=\frac{(1+y) \xi^{-1} h}{1-\xi^{-1} h} \mathrm{mC}\left(Z, \mathbb{P}^{n-1}\right)
$$

Applying the localization formula for the map $p$ we obtain

$$
\begin{equation*}
\frac{\mathrm{mC}\left(\Omega, \mathbb{C}^{n}\right)}{e^{K}\left(\mathbb{C}^{n}\right)}=p_{*}\left(\frac{(1+y) \xi^{-1} h}{1-\xi^{-1} h} \mathrm{mC}\left(Z, \mathbb{P}^{n-1}\right)\right) \tag{6}
\end{equation*}
$$

Let $u=h^{-1}-1 \in K\left(\mathbb{P}^{n-1}\right)$. Note that $u^{n}=0$. We have the expansion

$$
\begin{aligned}
\frac{(1+y) \xi^{-1} h}{1-\xi^{-1} h}= & \frac{(1+y) \xi^{-1}}{1+u-\xi^{-1}}=\frac{(1+y) \xi^{-1}}{\left(1-\xi^{-1}\right)\left(1+\frac{u}{1-\xi^{-1}}\right)} \\
& =\frac{(1+y) \xi^{-1}}{1-\xi^{-1}}\left(1-\frac{u}{1-\xi^{-1}}+\cdots+(-1)^{n-1} \frac{u^{n-1}}{\left(1-\xi^{-1}\right)^{n-1}}\right)
\end{aligned}
$$

The expression under the push-forward $p_{*}$ in (6) is of the form

$$
\frac{1}{\left(1-\xi^{-1}\right)^{n}} P\left(u, \xi^{-1}, y\right)
$$

The polynomial $P$ is of degree at most $n$ in $\xi^{-1}$, at most $n-1$ in $u$ and it is divisible by $(y+1) \xi^{-1}$. We have ${ }^{1} p_{*}\left(u^{k}\right)=1$ for $k<n$. After the push-forward we obtain

$$
\frac{\operatorname{mC}\left(\Omega, \mathbb{C}^{n}\right)}{e^{K}\left(\mathbb{C}^{n}\right)}=\frac{Q\left(\xi^{-1}, y\right)}{\left(1-\xi^{-1}\right)^{n}}
$$

with the polynomial $Q\left(\xi^{-1}, y\right)=P\left(1, \xi^{-1}, y\right)$ again divisible by $(y+1) \xi^{-1}$ and of degree at most $n$ in $\xi^{-1}$.

Remark 4.9 The argument above shows that $\mathrm{mC}\left(\Omega, \mathbb{C}^{n}\right)$ is divisible by $(y+1)$. In a similar way this divisibility can be proven for any quasihomogenous subvariety. Such variety can be presented as a quotient (by a finite group) of homogeneous one. Application of the Lefschetz-Riemann-Roch formula [11, Theorem 5.1] gives an explicit formula for $\mathrm{mC}\left(\Omega, \mathbb{C}^{n}\right)$.

### 4.5 An Example: Quadratic Cone—limits with Different Choices of 1-Parameter Subgroup

This example is based on the computations made in [24]. Let $\Omega \subset \mathbb{C}^{4}$ be the open set given by the inequality $z_{1} z_{2}-z_{3} z_{4} \neq 0$. This is an open orbit of the natural action of $\mathbb{C}^{*} \times O(4, \mathbb{C})$ on $\mathbb{C}^{4}$.

On the other hand, identifying $\mathbb{C}^{4}$ with $2 \times 2$ matrices $\left(\begin{array}{ll}z_{1} & z_{3} \\ z_{4} & z_{2}\end{array}\right)$ the set $\Omega$ is equal to the set of nondegenerate matrices. Its motivic Chern class was computed in [19, Sect. 8].

Denote the basis characters of the torus $\mathbb{T}=\left(\mathbb{C}^{*}\right)^{3}$ by $\alpha, \beta, \gamma$. Suppose $\mathbb{T}$ acts on $\mathbb{C}^{4}$ with characters

$$
\alpha \beta, \alpha / \beta, \alpha \gamma, \alpha / \gamma
$$

The action preserves the variety $\Omega$. The action has a unique fixed point at 0 . We study $\Omega \cap\left(\mathbb{C}^{4}\right)_{\{0\}}^{-}$, the Białynicki-Birula minus-cell (defined by (5)) intersected with the orbit, depending on the choice of the one parameter subgroup. We apply the formula [24, Formula 3] for $n=4$, which allows us to compute the motivic Chern class of $\Omega$ :

$$
\operatorname{mC}\left(\Omega, \mathbb{C}^{4}\right)=(1+y)^{2}\left(\frac{1}{\alpha^{4}} y^{2}+\left(\frac{\beta}{\alpha^{3}}+\frac{\gamma}{\alpha^{3}}+\frac{1}{\alpha^{3} \beta}+\frac{1}{\alpha^{3} \gamma}-\frac{1}{\alpha^{2}}-\frac{1}{\alpha^{4}}\right) y+\frac{1}{\alpha^{2}}\right)
$$

[^21]Below we present a sample of choices of one parameter subgroups and computed limits of motivic Chern classes:

| $(\alpha, \beta, \gamma)$ | Minus-cell $\left(\mathbb{C}^{4}\right)_{\{0\}^{-}}$ | Inequality | $\Omega \cap\left(\mathbb{C}^{4}\right)_{\{0\}^{-}}$ | limit of $\frac{\operatorname{mC}\left(\Omega \cap\left(\mathbb{C}^{4}\right)_{\{0\}}^{-}\right)_{0}}{e^{K}\left(\nu_{0}\right)}$ |
| :---: | :---: | :---: | :---: | :---: |
| $(\xi, 1,1)$ | $\{0\}$ | $0 \neq 0$ | $\emptyset$ | 0 |
| $\left(\xi^{-1}, 1,1\right)$ | $\mathbb{C}^{4}$ | $z_{1} z_{2}-z_{3} z_{4} \neq 0$ | cone complement <br> $\left[\mathbb{C}^{4}\right]-\left[\mathbb{C}^{*}\right]\left[\mathbb{P}^{1}\right]^{2}-[0]$ | $y^{4}+(y+1)(1-y)^{2}-1$ |
| $\left(\xi^{-1}, \xi^{2}, 1\right)$ | $z_{1}=0$ | $z_{3} z_{4} \neq 0$ | $\mathbb{C} \times\left(\mathbb{C}^{*}\right)^{2}$ | $-y(y+1)^{2}$ |

## 5 Flag Manifolds and Quiver Representation Spaces

After fixing some combinatorial codes we will define two related geometric objects: a flag variety and a quiver representation space, together with the description of the Borel orbits in these spaces. Then we will present formulas for the motivic Chern classes of the orbits.

### 5.1 Combinatorial Codes

Let $n, N$ be non-negative integers, and let $\mu=\left(\mu_{1}, \ldots, \mu_{N}\right) \in \mathbb{N}^{N}$ with $\sum_{i=1}^{N} \mu_{i}=$ $n$. Denote $\mu^{(i)}=\sum_{j=1}^{i} \mu_{j}$.

Let $\mathcal{I}_{\mu}$ denote the collection of $N$-tuples $\left(I_{1}, \ldots, I_{N}\right)$ with $I_{j} \subset\{1, \ldots, n\}, I_{i} \cap$ $I_{j}=\emptyset$ unless $i=j,\left|I_{i}\right|=\mu_{i}$. For $I \in \mathcal{I}_{\mu}$ we will use the following notation: $I^{(j)}=$ $\cup_{i=1}^{j} I_{i}=\left\{i_{1}^{(j)}<\cdots<i_{\mu^{(j)}}^{(j)}\right\}$.

For $I \in \mathcal{I}_{\mu}$ define $\ell(I)=\left|\left\{(a, b) \in\{1, \ldots, n\}^{2}: a>b, a \in I_{j}, b \in I_{k}, j<k\right\}\right|$.

### 5.2 Flag Variety

Let $\mathrm{Fl}_{\mu}$ denote the partial flag variety parameterizing chains of subspaces $V_{\bullet}=$ $\left(0 \subset V_{1} \subset V_{2} \subset \cdots \subset V_{N-1} \subset V_{N}=\mathbb{C}^{n}\right)$ with $\operatorname{dim}\left(V_{i}\right)=\mu^{(i)}$. Consider the natural action of the Borel subgroup $B_{n} \subset \mathrm{GL}_{n}(\mathbb{C})$ on $\mathrm{Fl}_{\mu}$. Let $\mathcal{F}_{i}$ be the tautological bundle over $\mathrm{Fl}_{\mu}$ whose fiber over the point ( $V_{\bullet}$ ) is $V_{i}$. Let the $B_{n}$-equivariant Ktheoretic Chern roots of $\mathcal{F}_{i}$ be $\alpha_{j}^{(i)}$, for $i=1, \ldots, N, j=1, \ldots, \mu^{(i)}$, that is, in $K_{B_{n}}\left(\mathrm{Fl}_{\mu}\right)$ we have $\sum_{j=1}^{\mu^{(i)}} \alpha_{j}^{(i)}=\mathcal{F}_{i}$. Then we have

$$
\begin{equation*}
K_{B_{n}}\left(\mathrm{Fl}_{\mu}\right)[y]=\mathbb{Z}\left[\left(\alpha_{j}^{(i)}\right)^{ \pm 1}, y\right]^{S_{\mu}(1)} \times \cdots \times S_{\mu^{(N-1)}} /(\text { certain ideal) } . \tag{7}
\end{equation*}
$$

Here the symmetric group $S_{\mu^{(k)}}$ permutes the variables $\alpha_{1}^{(k)}, \ldots, \alpha_{\mu^{(k)}}^{(k)}$.
Definition 5.1 For $I \in \mathcal{I}_{\mu}$ define the Schubert cell

$$
\Omega_{I}=\left\{\left(V_{i}\right) \in \mathrm{Fl}_{\mu}: \operatorname{dim}\left(V_{p} \cap \mathbb{C}_{\text {last }}^{q}\right)=\#\left\{i \in I^{(p)}: i>n-q\right\}, \forall p, q\right\} \subset \mathrm{Fl}_{\mu}
$$

where $\mathbb{C}_{\text {last }}^{q}$ is the span of the last $q$ standard basis vectors in $\mathbb{C}^{n}$. Its codimension in $\mathrm{Fl}_{\mu}$ is $\ell(I)$.

### 5.3 Quiver Representation Space

Consider $\operatorname{Rep}_{\mu}=\oplus_{j=1}^{N-1} \operatorname{Hom}\left(\mathbb{C}^{\mu^{(j)}}, \mathbb{C}^{\mu^{(j+1)}}\right)$ with the action of

$$
\mathrm{GL} \times B_{n}=\prod_{i=1}^{N-1} \mathrm{GL}_{\mu^{(i)}}(\mathbb{C}) \times B_{n}
$$

given by

$$
\left(g_{i}, \ldots, g_{N}\right) \cdot\left(a_{j}\right)_{j=1, \ldots, N-1}=\left(g_{j+1} \circ a_{j} \circ g_{j}^{-1}\right)_{j=1, \ldots, N-1}
$$

For $i=1, \ldots, N$ let $\mathcal{F}_{i}$ be the tautological rank $\mu^{(i)}$ bundle over the classifying space of the $i$ 'th component of $\mathrm{GL} \times B_{n}$. Let $\alpha_{j}^{(i)}$ be the K-theoretic Chern roots of $\mathcal{F}_{i}$, for $i=1, \ldots, N, j=1, \ldots, \mu^{(i)}$. Then we have

$$
\begin{equation*}
K_{\mathrm{GL} \times B_{n}}\left(\operatorname{Rep}_{\mu}\right)[y]=\mathbb{Z}\left[\left(\alpha_{j}^{(i)}\right)^{ \pm 1}, y\right]^{S_{\mu(1)} \times \cdots \times S_{\mu}(N-1)} \tag{8}
\end{equation*}
$$

Notice that the K-theory algebra in (7) is a quotient of the K-theory algebra in (8).
Definition 5.2 For $I \in \mathcal{I}_{\mu}$ define the matrix Schubert cell $M \Omega_{I} \subset \operatorname{Rep}_{\mu}$ by

$$
\begin{aligned}
& M \Omega_{I}=\left\{\left(a_{i}\right) \in \operatorname{Rep}_{\mu}: a_{i} \text { is injective } \forall i,\right. \\
& \\
& \left.\quad \operatorname{dim}\left(\left(a_{N-1} \circ \ldots \circ a_{p}\right)\left(\mathbb{C}^{\mu(p)}\right) \cap \mathbb{C}_{\text {last }}^{q}\right)=\#\left\{i \in I^{(p)}: i>n-q\right\}, \forall p, q\right\},
\end{aligned}
$$

where $\mathbb{C}_{\text {last }}^{q}$ is the span of the last $q$ standard basis vectors in $\mathbb{C}^{n}$. Its codimension in $\operatorname{Rep}_{\mu}$ is $\ell(I)$.

Remark 5.3 Of course we have that $\mathrm{Fl}_{\mu}=\operatorname{Rep}_{\mu} / / \mathrm{GL}$ as a GIT quotient space. Namely, let $\operatorname{Rep}_{\mu}^{s s}$ be the subset consisting only injective maps. Then the natural map

$$
\operatorname{Rep}_{\mu}^{s s} \rightarrow \mathrm{Fl}_{\mu}, \quad\left(a_{i}\right) \mapsto\left(\left(a_{N-1} \circ \ldots \circ a_{p}\right)\left(\mathbb{C}^{\mu^{(p)}}\right)\right)_{p}
$$

is a topological quotient by GL, in fact $\operatorname{Rep}_{\mu}^{s s} \rightarrow \mathrm{Fl}_{\mu}$ is a principal bundle, cf. [42]. Under this map, we have $M \Omega_{I}$ maps to $\Omega_{I}$. This correspondence between $\mathrm{Fl}_{\mu}$ and Rep $_{\mu}$ will be used in Sect. 8.

### 5.4 Weight Functions

In this section we define some explicit functions that will be used in naming the equivariant motivic Chern classes of Schubert and matrix Schubert cells. For more details on these functions see e.g. [31-33] or references therein.

For $I \in \mathcal{I}_{\mu}, j=1, \ldots, N-1, a=1, \ldots, \mu^{(j)}, b=1, \ldots, \mu^{(j+1)}$ define

$$
\psi_{I, j, a, b}(\xi)= \begin{cases}1-\xi & \text { if } i_{b}^{(j+1)}<i_{a}^{(j)} \\ (1+y) \xi & \text { if } i_{b}^{(j+1)}=i_{a}^{(j)} \\ 1+y \xi & \text { if } i_{b}^{(j+1)}>i_{a}^{(j)} .\end{cases}
$$

Define the weight function

$$
W_{I}=\operatorname{Sym}_{S_{\mu^{(1)}} \times \cdots \times S_{\mu}(N-1)} U_{I}
$$

where

$$
U_{I}=\prod_{j=1}^{N-1} \prod_{a=1}^{\mu^{(j)}} \prod_{b=1}^{\mu^{(j+1)}} \psi_{I, j, a, b}\left(\alpha_{a}^{(j)} / \alpha_{b}^{(j+1)}\right) \cdot \prod_{j=1}^{N-1} \prod_{1 \leq a<b \leq \mu^{(j)}} \frac{1+y \alpha_{b}^{(j)} / \alpha_{a}^{(j)}}{1-\alpha_{b}^{(j)} / \alpha_{a}^{(j)}}
$$

Here the symmetrizing operator is defined by

$$
\operatorname{Sym}_{S_{\mu^{(1)}} \times \cdots \times S_{\mu^{(N-1)}}} U_{I}=\sum_{\sigma \in S_{\mu^{(1)}} \times \cdots \times S_{\mu^{(N-1)}}} U_{I}\left(\sigma\left(\alpha_{a}^{(j)}\right)\right)
$$

where the $j$ th component of $\sigma$ (an element of $S_{\mu^{(j)}}$ ) permutes the $\alpha^{(j)}$ variables. For

$$
c_{\mu}=\prod_{j=1}^{N-1} \prod_{a=1}^{\mu^{(j)}} \prod_{b=1}^{\mu^{(j)}}\left(1+y \alpha_{b}^{(j)} / \alpha_{a}^{(j)}\right), \quad c_{\mu}^{\prime}=\prod_{j=1}^{N-1} \prod_{a=1}^{\mu^{(j+1)}} \prod_{b=1}^{\mu^{(j)}}\left(1+y \alpha_{b}^{(j)} / \alpha_{a}^{(j+1)}\right)
$$

define the modified weight functions

$$
\widetilde{W}_{I}=W_{I} / c_{\mu}, \quad \hat{W}_{I}=W_{I} / c_{\mu}^{\prime} .
$$

Observe that $\widetilde{W}_{I}, \hat{W}_{I}$ are not Laurent polynomials, but rather ratios of two such.

Weight functions were defined by Tarasov and Varchenko in relation with hypergeometric solutions to qKZ differential equations, see e.g. [37].

### 5.5 Motivic Chern Classes of Schubert Cells Given by Weight Functions

Theorem 5.4 For the motivic Chern classes of matrix Schubert and ordinary Schubert cells we have

$$
\begin{align*}
\mathrm{mC}\left(M \Omega_{I}, \operatorname{Rep}_{\mu}\right)=W_{I} & \in K_{\mathrm{GL} \times B_{n}}\left(\operatorname{Rep}_{\mu}\right)[y]  \tag{9}\\
\mathrm{mC}\left(\Omega_{I}, \mathrm{Fl}_{\mu}\right)=\left[\widetilde{W}_{I}\right] & \in K_{B_{n}}\left(\mathrm{Fl}_{\mu}\right)[y] \tag{10}
\end{align*}
$$

First let us comment on the two statements of the theorem. As we saw in (8) the ring $K_{\mathrm{GL} \times B_{n}}\left(\operatorname{Rep}_{\mu}\right)[y]$ is a Laurent polynomial ring, and claim (9) states which Laurent polynomial is the sought MC class. However, the ring $K_{\mathrm{GL} \times B_{n}}\left(\mathrm{Fl}_{\mu}\right)[y]$ is a quotient ring of that Laurent polynomial ring, cf. (7), hence (10) only names a representative of the sought MC class. Moreover, the function $\widetilde{W}_{I}$ is a rational function, so the equality (10) is meant in the following sense: the restriction to each torus fixed of the two sides of (10) are the same-in particular the fix point restrictions of $\widetilde{W}_{I}$ are Laurent polynomials.
Remark 5.5 If we define the motivic Segre class of $X \subset M$ by

$$
\mathrm{mS}(X, M)=\mathrm{mC}(X, M) / c^{K}(T M)
$$

then we can rephrase Theorem 5.4 in the more symmetric form

$$
\operatorname{mS}\left(M \Omega_{I}, \operatorname{Rep}_{\mu}\right)=\hat{W}_{I}, \quad \operatorname{mS}\left(\Omega_{I}, \mathrm{Fl}_{\mu}\right)=\left[\hat{W}_{I}\right] .
$$

This is due to the fact that $c_{\mathrm{GL} \times B_{n}}^{K}\left(T \operatorname{Rep}_{\mu}\right)=c_{\mu}^{\prime}$, and $\left[c_{B_{n}}^{K}\left(T \mathrm{Fl}_{\mu}\right)\right]=c_{\mu}^{\prime} / c_{\mu}$. That is, $c_{\mu}$ is the K-theoretic total Chern class of the fibers of the GIT quotient map of Remark 5.3.

Now let us review different strategies that can be used to prove Theorem 5.4some of them already present in the literature. One can
(a) prove (9) by resolution of singularities;
(b) prove (9) by the interpolation Theorem 4.6;
(c) prove (10) by resolution of singularities;
(d) prove (10) by the interpolation Theorem 4.6;
(e) prove that (9) implies (10) via Remark 5.5.

Carrying out (b) and/or (d) has the advantage of not having to construct a resolution with normal crossing divisors. Carrying out (a) and/or (c) has the advantage of not having to prove the sophisticated Newton polygon properties of the weight functions.

Historically, first the Newton polygon properties of weight functions were proved (in the context of K-theoretic stable envelopes), see [33, Sect.3.5]. Those proofs are complicated algebraic arguments that were unknown before even by the experts of weight functions. Nevertheless, those arguments provide a complete proof of Theorem 5.4, see details in [19].

It is remarkable that resolutions for $\overline{M \Omega}_{I}, \bar{\Omega}_{I}$ can be constructed from which the motivic Chern classes are calculated easily. More importantly, for well chosen resolutions the formulas obtained for $\mathrm{mC}\left(M \Omega_{I}, \operatorname{Rep}_{\mu}\right), \mathrm{mC}\left(\Omega_{I}, \mathrm{Fl}_{\mu}\right)$ are exactly the defining formulas of weight functions. As a result, such an argument reproves Theorem 5.4, and through that, the Newton polygon properties of weight functions without calculation. In hindsight, such an argument would be a more natural proof of Theorem 5.4. Since formally proving Theorem 5.4 is not needed anymore (as we mentioned above, [33, Sect. 3.5] provides a proof) we will not present a proof based on resolution of singularities in detail. Also, such an argument based on resolution is implicitly present in [30] and an analogous argument is explicitly presented in the elliptic setting in [35]. As a final remark let us mention our appreciation of the creativity of those who defined weight functions without seeing the resolution calculation for $\mathrm{mC}\left(M \Omega_{I}, \operatorname{Rep}_{\mu}\right)$ and $\mathrm{mC}\left(\Omega_{I}, \mathrm{Fl}_{\mu}\right)$ !

The implication (e) is not formally needed to complete the proof of Theorem 5.4. Yet, we find it important enough, that we include the proof of the general such statement (namely motivic Chern classes before vs after GIT quotient) in Sect. 8 .

## 6 MC Classes of Schubert Cells in Full Flag Manifolds: Positivity and Log-Concavity

Characteristic classes of singularities often display positivity properties. Motivic Chern classes are not exceptions: in this section we present three conjectures on the signs (and concavity) of coefficients of MC classes in certain expansions.

### 6.1 Positivity

For $\mu=(1,1, \ldots, 1)$ ( $n$ times) the space $\mathrm{Fl}_{\mu}$ is the full flag variety, let us rename it to $\mathrm{Fl}(n)$. Recall the presentation of its equivariant K-theory algebra from (7). For brevity let us rename the "last" set of variables, the "equivariant variables" to $\tau_{i}:=\alpha_{i}^{(n)}$.

The traditional geometric basis of $K_{B_{n}}(\mathrm{Fl}(n))[y]$ consists of the classes of the structure sheaves of the Schubert varieties (the closures of the Schubert cells). For $w=(w(1), \ldots, w(n)) \in S_{n}$ let $[\mathbf{w}]$ denote the class of the structure sheaf of the closure of $\Omega_{(\{w(1)\},\{w(2)\}, \ldots,\{w(n)\})}$. Also, set $\mathrm{mC}[w]=\mathrm{mC}\left(\Omega_{(\{w(1)\},\{w(2)\}, \ldots,\{w(n)\})}\right)$.

Our Theorem 5.4 can be used to expand the MC classes of orbits in the traditional basis of structure sheaves of Schubert varieties. For $n=3$ we obtain the following expansions.

$$
\begin{align*}
& \mathrm{mC}[1,2,3]=\left(\frac{\tau_{1}^{2}}{\tau_{3}^{2}} y^{3}+\left(\frac{\tau_{1}^{2}}{\tau_{2} \tau_{3}}+\frac{\tau_{1}}{\tau_{3}}+\frac{\tau_{2} \tau_{1}}{\tau_{3}^{2}}\right) y^{2}+\left(\frac{\tau_{1}}{\tau_{2}}+\frac{\tau_{1}}{\tau_{3}}+\frac{\tau_{2}}{\tau_{3}}\right) y+1\right)[\mathbf{1}, \mathbf{2}, \mathbf{3}] \\
& -\left(\left(\frac{\tau_{1}^{2}}{\tau_{2} \tau_{3}}+\frac{\tau_{1}^{2}}{\tau_{3}^{2}}\right) y^{3}+\left(\frac{\tau_{1}^{2}}{\tau_{2} \tau_{3}}+\frac{\tau_{1}}{\tau_{2}}+\frac{2 \tau_{1}}{\tau_{3}}+\frac{\tau_{2} \tau_{1}}{\tau_{3}^{2}}\right) y^{2}+\left(\frac{\tau_{1}}{\tau_{2}}+\frac{\tau_{1}}{\tau_{3}}+\frac{\tau_{2}}{\tau_{3}}+1\right) y+1\right)  \tag{1,3,2}\\
& -\left(\left(\frac{\tau_{1}^{2}}{\tau_{3}^{2}}+\frac{\tau_{2} \tau_{1}}{\tau_{3}^{2}}\right) y^{3}+\left(\frac{\tau_{1}^{2}}{\tau_{2} \tau_{3}}+\frac{2 \tau_{1}}{\tau_{3}}+\frac{\tau_{2} \tau_{1}}{\tau_{3}^{2}}+\frac{\tau_{2}}{\tau_{3}}\right) y^{2}+\left(\frac{\tau_{1}}{\tau_{2}}+\frac{\tau_{1}}{\tau_{3}}+\frac{\tau_{2}}{\tau_{3}}+1\right) y+1\right)  \tag{2,1,3}\\
& +\left(\left(\frac{\tau_{1}^{2}}{\tau_{2} \tau_{3}}+\frac{\tau_{1}^{2}}{\tau_{3}^{2}}+\frac{\tau_{1}}{\tau_{3}}+\frac{\tau_{2} \tau_{1}}{\tau_{3}^{2}}+\frac{\tau_{2}}{\tau_{3}}\right) y^{3}+\left(\frac{\tau_{1}^{2}}{\tau_{2} \tau_{3}}+\frac{\tau_{1}}{\tau_{2}}+\frac{3 \tau_{1}}{\tau_{3}}+\frac{\tau_{2} \tau_{1}}{\tau_{3}^{2}}+\frac{2 \tau_{2}}{\tau_{3}}+1\right) y^{2}\right. \\
& \left.+\left(\frac{\tau_{1}}{\tau_{2}}+\frac{\tau_{1}}{\tau_{3}}+\frac{\tau_{2}}{\tau_{3}}+2\right) y+1\right)  \tag{2,3,1}\\
& +\left(\left(\frac{\tau_{1}^{2}}{\tau_{2} \tau_{3}}+\frac{\tau_{1}^{2}}{\tau_{3}^{2}}+\frac{\tau_{1}}{\tau_{2}}+\frac{\tau_{1}}{\tau_{3}}+\frac{\tau_{2} \tau_{1}}{\tau_{3}^{2}}\right) y^{3}+\left(\frac{\tau_{1}^{2}}{\tau_{2} \tau_{3}}+\frac{2 \tau_{1}}{\tau_{2}}+\frac{3 \tau_{1}}{\tau_{3}}+\frac{\tau_{2} \tau_{1}}{\tau_{3}^{2}}+\frac{\tau_{2}}{\tau_{3}}+1\right) y^{2}\right. \\
& \left.+\left(\frac{\tau_{1}}{\tau_{2}}+\frac{\tau_{1}}{\tau_{3}}+\frac{\tau_{2}}{\tau_{3}}+2\right) y+1\right)[\mathbf{3}, \mathbf{1}, \mathbf{2}]  \tag{3,1,2}\\
& -\left(\left(\frac{\tau_{1}^{2}}{\tau_{2} \tau_{3}}+\frac{\tau_{1}^{2}}{\tau_{3}^{2}}+\frac{\tau_{1}}{\tau_{2}}+\frac{2 \tau_{1}}{\tau_{3}}+\frac{\tau_{2} \tau_{1}}{\tau_{3}^{2}}+\frac{\tau_{2}}{\tau_{3}}+1\right) y^{3}+\left(\frac{\tau_{1}^{2}}{\tau_{2} \tau_{3}}+\frac{2 \tau_{1}}{\tau_{2}}+\frac{3 \tau_{1}}{\tau_{3}}+\frac{\tau_{2} \tau_{1}}{\tau_{3}^{2}}+\frac{2 \tau_{2}}{\tau_{3}}+2\right) y^{2}\right. \\
& \left.+\left(\frac{\tau_{1}}{\tau_{2}}+\frac{\tau_{1}}{\tau_{3}}+\frac{\tau_{2}}{\tau_{3}}+2\right) y+1\right)[\mathbf{3}, \mathbf{2}, \mathbf{1}] . \tag{3,2,1}
\end{align*}
$$

$\mathrm{mC}[1,3,2]=\left(\frac{\tau_{1}^{2}}{\tau_{2} \tau_{3}} y^{2}+\left(\frac{\tau_{1}}{\tau_{2}}+\frac{\tau_{1}}{\tau_{3}}\right) y+1\right)[\mathbf{1}, \mathbf{3}, \mathbf{2}]$

$$
\begin{aligned}
& -\left(\left(\frac{\tau_{1}^{2}}{\tau_{2} \tau_{3}}+\frac{\tau_{1}}{\tau_{3}}+\frac{\tau_{2}}{\tau_{3}}\right) y^{2}+\left(\frac{\tau_{1}}{\tau_{2}}+\frac{\tau_{1}}{\tau_{3}}+\frac{\tau_{2}}{\tau_{3}}+1\right) y+1\right)[\mathbf{2}, \mathbf{3}, \mathbf{1}] \\
& -\left(\left(\frac{\tau_{1}^{2}}{\tau_{2} \tau_{3}}+\frac{\tau_{1}}{\tau_{2}}\right) y^{2}+\left(\frac{\tau_{1}}{\tau_{2}}+\frac{\tau_{1}}{\tau_{3}}+1\right) y+1\right)[\mathbf{3}, \mathbf{1}, \mathbf{2}] \\
& +\left(\left(\frac{\tau_{1}^{2}}{\tau_{2} \tau_{3}}+\frac{\tau_{1}}{\tau_{2}}+\frac{\tau_{1}}{\tau_{3}}+\frac{\tau_{2}}{\tau_{3}}+1\right) y^{2}+\left(\frac{\tau_{1}}{\tau_{2}}+\frac{\tau_{1}}{\tau_{3}}+\frac{\tau_{2}}{\tau_{3}}+2\right) y+1\right)[\mathbf{3}, \mathbf{2}, \mathbf{1}]
\end{aligned}
$$

$\mathrm{mC}[2,1,3]=\left(\frac{\tau_{1} \tau_{2}}{\tau_{3}^{2}} y^{2}+\left(\frac{\tau_{1}}{\tau_{3}}+\frac{\tau_{2}}{\tau_{3}}\right) y+1\right)[\mathbf{2}, \mathbf{1}, \mathbf{3}]$
$-\left(\left(\frac{\tau_{2}}{\tau_{3}}+\frac{\tau_{1} \tau_{2}}{\tau_{3}^{2}}\right) y^{2}+\left(\frac{\tau_{1}}{\tau_{3}}+\frac{\tau_{2}}{\tau_{3}}+1\right) y+1\right)[\mathbf{2}, \mathbf{3}, \mathbf{1}]$
$-\left(\left(\frac{\tau_{1}}{\tau_{2}}+\frac{\tau_{1}}{\tau_{3}}+\frac{\tau_{2} \tau_{1}}{\tau_{3}^{2}}\right) y^{2}+\left(\frac{\tau_{1}}{\tau_{2}}+\frac{\tau_{1}}{\tau_{3}}+\frac{\tau_{2}}{\tau_{3}}+1\right) y+1\right)[\mathbf{3}, \mathbf{1}, \mathbf{2}]$
$+\left(\left(\frac{\tau_{1}}{\tau_{2}}+\frac{\tau_{1}}{\tau_{3}}+\frac{\tau_{2} \tau_{1}}{\tau_{3}^{2}}+\frac{\tau_{2}}{\tau_{3}}+1\right) y^{2}+\left(\frac{\tau_{1}}{\tau_{2}}+\frac{\tau_{1}}{\tau_{3}}+\frac{\tau_{2}}{\tau_{3}}+2\right) y+1\right)[\mathbf{3}, \mathbf{2}, \mathbf{1}]$.
$\mathrm{mC}[2,3,1]=\left(\frac{\tau_{2}}{\tau_{3}} y+1\right)[\mathbf{2}, \mathbf{3}, \mathbf{1}]$

$$
-\left(\left(\frac{\tau_{2}}{\tau_{3}}+1\right) y+1\right)[\mathbf{3}, \mathbf{2}, \mathbf{1}] .
$$

$\mathrm{mC}[3,1,2]=\left(\frac{\tau_{1}}{\tau_{2}} y+1\right)[\mathbf{3}, \mathbf{1}, \mathbf{2}]$
$-\left(\left(\frac{\tau_{1}}{\tau_{2}}+1\right) y+1\right)[\mathbf{3}, \mathbf{2}, \mathbf{1}]$.
$\mathrm{mC}[3,2,1]=[\mathbf{3}, \mathbf{2}, \mathbf{1}]$.
The following conjecture can be verified in the formulas above, and we also verified it for larger flag varieties ( $n \leq 5$ ).

Conjecture 6.1 Let $p, w \in S_{n}$. The coefficient of $[\mathbf{w}]$ in the expansion of $\mathrm{mC}[p]$ is a Laurent polynomial in $\tau_{1}, \ldots, \tau_{n}, y$ whose terms have sign $(-1)^{\ell(p)-\ell(w)}$.

We have been informed by the authors of [4] that they also observed the sign behavior described in Conjecture 6.1.

### 6.2 Log Concavity

To illustrate a new feature of the coefficients of the [w]-expansions, let us make the substitution $\tau_{i}=1(\forall i)$, that is, consider the non-equivariant motivic Chern classes of the Schubert cells. For $n=4$ and for the open cell we obtain

$$
\begin{aligned}
\mathrm{mC}[1,2,3,4] & =(y+1)^{6}[\mathbf{1}, \mathbf{2}, \mathbf{3}, \mathbf{4}] \\
& -(y+1)^{5}(2 y+1)[\mathbf{1}, \mathbf{2}, \mathbf{4}, \mathbf{3}] \\
& -(y+1)^{5}(2 y+1)[\mathbf{1}, \mathbf{3}, \mathbf{2}, \mathbf{4}] \\
& -(y+1)^{5}(2 y+1)[\mathbf{2}, \mathbf{1}, \mathbf{3}, \mathbf{4}] \\
& +(y+1)^{4}\left(5 y^{2}+4 y+1\right)[\mathbf{1}, \mathbf{3}, \mathbf{4}, \mathbf{2}] \\
& +(y+1)^{4}\left(5 y^{2}+4 y+1\right)[\mathbf{1}, \mathbf{4}, \mathbf{2}, \mathbf{3}] \\
& +(y+1)^{4}(2 y+1)^{2}[\mathbf{2}, \mathbf{1}, \mathbf{4}, \mathbf{3}] \\
& +(y+1)^{4}\left(5 y^{2}+4 y+1\right)[\mathbf{2}, \mathbf{3}, \mathbf{1}, \mathbf{4}] \\
& +(y+1)^{4}\left(5 y^{2}+4 y+1\right)[\mathbf{3}, \mathbf{1}, \mathbf{2}, \mathbf{4}] \\
& -(y+1)^{3}\left(8 y^{3}+11 y^{2}+5 y+1\right)[\mathbf{1}, \mathbf{4}, \mathbf{3}, \mathbf{2}] \\
& -(y+1)^{3}\left(14 y^{3}+14 y^{2}+6 y+1\right)[\mathbf{2}, \mathbf{3}, \mathbf{4}, \mathbf{1}] \\
& -(y+1)^{3}\left(13 y^{3}+13 y^{2}+6 y+1\right)[\mathbf{2}, \mathbf{4}, \mathbf{1}, \mathbf{3}] \\
& -(y+1)^{3}(2 y+1)\left(6 y^{2}+4 y+1\right)[\mathbf{3}, \mathbf{1}, \mathbf{4}, \mathbf{2}] \\
& -(y+1)^{3}\left(8 y^{3}+11 y^{2}+5 y+1\right)[\mathbf{3}, \mathbf{2}, \mathbf{1}, \mathbf{4}] \\
& -(y+1)^{3}\left(14 y^{3}+14 y^{2}+6 y+1\right)[\mathbf{4}, \mathbf{1}, \mathbf{2}, \mathbf{3}] \\
& +(y+1)^{2}\left(24 y^{4}+36 y^{3}+21 y^{2}+7 y+1\right)[\mathbf{2}, \mathbf{4}, \mathbf{3}, \mathbf{1}] \\
& +(y+1)^{2}\left(23 y^{4}+35 y^{3}+22 y^{2}+7 y+1\right)[\mathbf{3}, \mathbf{2}, \mathbf{4}, \mathbf{1}] \\
& +(y+1)^{2}(3 y+1)\left(10 y^{3}+10 y^{2}+4 y+1\right)[\mathbf{3}, \mathbf{4}, \mathbf{1}, \mathbf{2}] \\
& +(y+1)^{2}\left(23 y^{4}+35 y^{3}+22 y^{2}+7 y+1\right)[\mathbf{4}, \mathbf{1}, \mathbf{3}, \mathbf{2}] \\
& +(y+1)^{2}\left(24 y^{4}+36 y^{3}+21 y^{2}+7 y+1\right)[\mathbf{4}, \mathbf{2}, \mathbf{1}, \mathbf{3}] \\
& -(y+1)\left(44 y^{5}+85 y^{4}+66 y^{3}+29 y^{2}+8 y+1\right)[\mathbf{3}, \mathbf{4}, \mathbf{2}, \mathbf{1}] \\
& -(y+1)\left(49 y^{5}+91 y^{4}+69 y^{3}+30 y^{2}+8 y+1\right)[\mathbf{4}, \mathbf{2}, \mathbf{3}, \mathbf{1}] \\
& -(y+1)\left(44 y^{5}+85 y^{4}+66 y^{3}+29 y^{2}+8 y+1\right)[\mathbf{4}, \mathbf{3}, \mathbf{1}, \mathbf{2}] \\
& +\left(64 y^{6}+163 y^{5}+169 y^{4}+98 y^{3}+37 y^{2}+9 y+1\right)[\mathbf{4}, \mathbf{3}, \mathbf{2}, \mathbf{1}]
\end{aligned}
$$

In the expression above the permutations are ordered primary by length, secondary by lexicographical order.

The coefficient of $[\mathbf{5}, \mathbf{4}, \mathbf{3}, \mathbf{2}, \mathbf{1}]$ in $\mathrm{mC}[1,2,3,4,5]$ is

$$
1+14 y+92 y^{2}+377 y^{3}+1120 y^{4}+2630 y^{5}+4972 y^{6}+7148 y^{7}+7024 y^{8}+4063 y^{9}+1024 y^{10}
$$

Definition 6.2 The polynomial $\sum_{k=0}^{d} a_{k} y^{k}$ is said to be log-concave if $a_{k}^{2} \geq a_{k-1} a_{k+1}$ for $0<k<d$. It is strictly log-concave if the equality is strict.

See Huh's survey article [21] for the role of log-concavity in geometry and combinatorics.

Conjecture 6.3 The coefficients in the [w]-expansion of the non-equivariant $\mathrm{mC}[p]$ classes are strictly log-concave.

We checked Conjecture 6.3 for $n \leq 6$.

### 6.3 Positivity in New Variables

There is another kind of positivity which does not follow from the conjectured properties above. Namely, let us substitute $\frac{\tau_{i}}{\tau_{i+1}}=s_{i}+1$ and $y=-1-\delta$. For example

$$
\begin{aligned}
\mathrm{mC}[4,3,1,2] & =\left(1+y \frac{\tau_{1}}{\tau_{2}}\right)[\mathbf{4}, \mathbf{3}, \mathbf{1}, \mathbf{2}]-y[\mathbf{4}, \mathbf{3}, \mathbf{2}, \mathbf{1}] \\
& =-\left(s_{1}+\delta\left(1+s_{1}\right)\right)[\mathbf{4}, \mathbf{3}, \mathbf{1}, \mathbf{2}]+(1+\delta)[\mathbf{4}, \mathbf{3}, \mathbf{2}, \mathbf{1}]
\end{aligned}
$$

The coefficient of $[\mathbf{4}, \mathbf{3}, \mathbf{2}, \mathbf{1}]$ in $\mathrm{mC}[1,4,3,2]$ in the new variables equals

$$
\begin{aligned}
& \left(1+s_{1}\right)^{3}\left(1+s_{2}\right)^{2}\left(1+s_{3}\right)+ \\
& \delta\left(9+17 s_{1}+12 s_{1}^{2}+3 s_{1}^{3}+14 s_{2}+28 s_{1} s_{2}+22 s_{1}^{2} s_{2}+6 s_{1}^{3} s_{2}+6 s_{2}^{2}+12 s_{1} s_{2}^{2}+\right. \\
& 10 s_{1}^{2} s_{2}^{2}+3 s_{1}^{3} s_{2}^{2}+8 s_{3}+15 s_{1} s_{3}+11 s_{1}^{2} s_{3}+3 s_{1}^{3} s_{3}+13 s_{2} s_{3}+26 s_{1} s_{2} s_{3}+21 s_{1}^{2} s_{2} s_{3}+ \\
& \left.6 s_{1}^{3} s_{2} s_{3}+6 s_{2}^{2} s_{3}+12 s_{1} s_{2}^{2} s_{3}+10 s_{1}^{2} s_{2}^{2} s_{3}+3 s_{1}^{3} s_{2}^{2} s_{3}\right)+ \\
& \delta^{2}\left(21+28 s_{1}+15 s_{1}^{2}+3 s_{1}^{3}+26 s_{2}+40 s_{1} s_{2}+26 s_{1}^{2} s_{2}+6 s_{1}^{3} s_{2}+9 s_{2}^{2}+15 s_{1} s_{2}^{2}+\right. \\
& 11 s_{1}^{2} s_{2}^{2}+3 s_{1}^{3} s_{2}^{2}+16 s_{3}+22 s_{1} s_{3}+13 s_{1}^{2} s_{3}+3 s_{1}^{3} s_{3}+22 s_{2} s_{3}+35 s_{1} s_{2} s_{3}+24 s_{1}^{2} s_{2} s_{3}+ \\
& \left.6 s_{1}^{3} s_{2} s_{3}+9 s_{2}^{2} s_{3}+15 s_{1} s_{2}^{2} s_{3}+11 s_{1}^{2} s_{2}^{2} s_{3}+3 s_{1}^{3} s_{2}^{2} s_{3}\right)+ \\
& \delta_{3}^{3}\left(14+14 s_{1}+6 s_{1}^{2}+s_{1}^{3}+14 s_{2}+18 s_{1} s_{2}+10 s_{1}^{2} s_{2}+2 s_{1}^{3} s_{2}+4 s_{2}^{2}+6 s_{1} s_{2}^{2}+\right. \\
& 4 s_{1}^{2} s_{2}^{2}+s_{1}^{3} s_{2}^{2}+9 s_{3}+10 s_{1} s_{3}+5 s_{1}^{2} s_{3}+s_{1}^{3} s_{3}+11 s_{2} s_{3}+15 s_{1} s_{2} s_{3}+9 s_{1}^{2} s_{2} s_{3}+ \\
& \left.2 s_{1}^{3} s_{2} s_{3}+4 s_{2}^{2} s_{3}+6 s_{1} s_{2}^{2} s_{3}+4 s_{1}^{2} s_{2}^{2} s_{3}+s_{1}^{3} s_{2}^{2} s_{3}\right) .
\end{aligned}
$$

Conjecture 6.4 The coefficients of the $s_{i}, \delta$-monomials in the [w]-expansion of the $\mathrm{mC}[p]$ classes have sign $(-1)^{\ell(w)}$. Note that here the sign depends only on the length the permutation $w$, not on the length of $p$, as it was in Conjecture 6.1.

Remark 6.5 The positivity properties of Sects. 6.1 and 6.3 imply positivity properties of the restrictions of MC classes to fixed points. These "local" positivity properties are another instances of the positivity in the $\delta$-variables discussed in [40, Sect. 15]. It is implied by the Conjecture 6.4 and [5].

## 7 Local Picture

To illustrate the pretty convex geometric objects encoded by motivic Chern classes in this section we present some pictures. Consider $\mathrm{Fl}(3)$ and its six Schubert cells parameterized by permutations. According to Theorem 5.4 the six fixed point restrictions of these classes satisfy some strict containment properties. Below we present the Newton polygons (and some related weights) of the six fixed point restrictions (rows) of the six Schubert cells (columns). The Newton polygons live in the $\tau_{1}+\tau_{2}+\tau_{3}=0$ plane of $\mathbb{R}^{3}$ (with coordinates $\tau_{i}$ ), hence instead of 3D pictures we only draw the mentioned plane.


In blue-Newton polygon of $e^{K}\left(\nu_{\Theta}\right)=\lambda_{-1}\left(\nu_{\Theta}^{*}\right)$.
In red-cotangent weights of $T_{\Theta}^{*}$.
In violet-Newton polygon of $\mathrm{mC}\left(\Omega \cap S_{\Theta}\right)$, where $S_{\Theta}$ is a slice to the orbit.
Edges in violet-Newton polygon of $\mathrm{mC}(\Omega)$.

## 8 Transversality and GIT Quotients

Both CSM and MC classes have their Segre version (cf. Remark 5.5): for example for the K-theory case we have

$$
\mathrm{mS}(f: Z \rightarrow M):=\mathrm{mC}(f) / \lambda_{y}\left(T^{*} M\right)
$$

### 8.1 Transversality and Motivic Segre Class

Motivic Segre classes behave well for transversal pull-back, for a fine notion of transversality. To formulate this notion let us recall the definition of the motivic Chern class given in [19]. Consider the case of $\mathrm{mC}(f)$ where $f: U \rightarrow M$ is an inclusion of smooth varieties, but the inclusion is not necessarily closed (or equivalently, proper).
Definition 8.1 A proper normal crossing extension (PNC) of $f: U \rightarrow M$ is a morphism $\bar{f}: Y \rightarrow M$ and an inclusion $j: U \hookrightarrow Y$ such that
(1) $f=\bar{f} \circ j$,
(2) $Y$ is smooth
(3) $\bar{f}$ is proper
(4) The exceptional divisor $D:=Y \backslash j(U)=\bigcup_{i=1}^{s} D_{i}$ is a simple normal crossing divisor (i.e. the $D_{i}$ 's are smooth hypersurfaces in transversal position).

If $f$ is $G$-equivariant for some group $G$ acting on $U$ and $M$, then we require all maps to be $G$-equivariant in the definition.

Remark 8.2 Notice that if $f$ is injective then $\bar{f}: Y \rightarrow M$ is a resolution of the closure of $f(U)$.

We can use the existence of proper normal crossing extensions to define the ( $G$ equivariant) motivic Chern class:

Definition 8.3 Suppose that $f: U \rightarrow M$ is a map of smooth $G$-varieties and let $\bar{f}: Y \rightarrow M$ be a proper normal crossing extension of $f$. For $I \subset \underline{s}=\{1,2, \ldots, s\}$ let $D_{I}=\bigcap_{i \in I} D_{i}, f_{I}=\bar{f} \mid D_{I}$, in particular $f_{\emptyset}=\bar{f}$. Then

$$
\begin{equation*}
\mathrm{mC}(f):=\sum_{I \subset \underline{s}}(-1)^{|I|} f_{I *} \lambda_{y}\left(T^{*} D_{I}\right) \tag{11}
\end{equation*}
$$

It is explained in [19] that this is a good definition: independent of the PNC chosen.
Definition 8.4 Let $N$ be a smooth variety. Then $g: N \rightarrow M$ is motivically transver$\operatorname{sal}^{2}$ to $f: U \rightarrow M$ if there is a PNC $\bar{f}: Y \rightarrow M$ for $f$ such that $g$ is transversal to all the $D_{I}$ 's.

Theorem 8.5 If $g: N \rightarrow M$ is motivically transversal to $f: U \rightarrow M$ then

$$
\mathrm{mS}(\tilde{f})=g^{*} \mathrm{mS}(f)
$$

where $\tilde{f}: \tilde{U} \rightarrow N$ is the map in the pull-back diagram

[^22]

The proof is a straightforward consequence of the fact that a motivically transversal pull-back of a PNC is PNC. We can extend the notion to non-smooth varieties $Z$ requiring the existence of a stratification of $Z$ such that $g$ is motivically transversal to the restrictions of $f$ to these strata.

Being motivically transversal is a very restrictive condition. However it holds in some special situations:

Proposition 8.6 Let $U, M$ be smooth and a connected group $G$ is acting on $M$. Assume that $f: U \rightarrow M$ is transversal to all $G$-orbits. Then $f$ is motivically transversal to all $G$-invariant subvarieties of $M$.

The proof is simple and can be found in the proof of Lemma 5.1 in [19].

### 8.2 G-equivariant Motivic Segre Class as a Universal Motivic Segre Class for Degeneracy Loci

Now we prove a statement (Corollary 8.8) expressing the fact that the $G$-equivariant motivic Segre class is a universal formula for motivic Chern classes of degeneracy loci. The analogous statement for the Segre version of the CSM class was proved in [26]. In K-theory the proof is simpler because reference to classifying spaces and maps can be avoided. Suppose that $\pi_{P}: P \rightarrow M$ is a principal $G$-bundle over the smooth $M$ and $A$ is a smooth $G$-variety. Then we can define a map

$$
\begin{equation*}
a: K_{G}(A) \rightarrow K\left(P \times_{G} A\right) \tag{12}
\end{equation*}
$$

by association: For any $G$-vector bundle $E$ over $A$ the associated bundle $P \times_{G} E$ is a vector bundle over $P \times_{G} A$.

In the rest of the paper we will use the notation $\mathrm{mS}(A, B)$ for $\mathrm{mS}(f: A \rightarrow B)$ when the map $f$ is clear from the context. The diagrams

will be useful when reading the proof of the next Proposition.

Proposition 8.7 Let $Y \subset A$ be $G$-invariant. Then

$$
\operatorname{mS}\left(P \times_{G} Y, P \times_{G} A\right)=a\left(\operatorname{mS}_{G}(Y, A)\right)
$$

Proof To calculate the left hand side we need to calculate $\lambda_{y}\left(T^{*} P \times_{G} A\right)$ first. For the tangent space we have

$$
T\left(P \times_{G} A\right)=\pi_{A}^{*}(T M) \oplus\left(P \times_{G} T A\right)
$$

where $\pi_{A}: P \times{ }_{G} A \rightarrow M$ is the projection. Consequently

$$
\lambda_{y}\left(T^{*}\left(P \times_{G} A\right)\right)=\pi_{A}^{*} \lambda_{y}\left(T^{*} M\right) \cdot a\left(\lambda_{y}^{G}\left(T^{*} A\right)\right)
$$

Assume first that $Y$ is a closed submanifold of $A$. To ease the notation we will use the $\lambda(M)=\lambda_{y}\left(T^{*} M\right)$ abbreviation. Then

$$
\mathrm{mS}\left(P \times_{G} Y, P \times_{G} A\right)=\frac{i_{*}^{P} \lambda\left(P \times_{G} Y\right)}{\lambda\left(P \times_{G} A\right)}=\frac{i_{*}\left(\pi_{Y}^{*} \lambda(M) a\left(\lambda^{G}(Y)\right)\right.}{\pi_{A}^{*} \lambda(M) a\left(\lambda^{G}(A)\right)}
$$

where $\pi_{Y}: P \times_{G} Y \rightarrow M$ is the projection and $i: Y \rightarrow A, i_{P}: P \times_{G} Y \rightarrow P \times_{G}$ $A$ are the inclusions. Then noticing that $\pi_{Y}=\pi_{A} \circ i_{P}$ and applying the adjunction formula we arrive at the right hand side.

For general $Y$ we can use Definition 8.3 to reduce the calculation to the smooth case.

Corollary 8.8 Suppose that $\sigma: M \rightarrow P \times_{G} A$ is a section motivically transversal to $P \times_{G} Y$. Then

$$
\begin{equation*}
\mathrm{mS}(Y(\sigma), M)=\sigma^{*} a\left(\mathrm{mS}_{G}(Y, A)\right) \tag{13}
\end{equation*}
$$

where $Y(\sigma)=\sigma^{-1}\left(P \times_{G} Y\right)$ is the $Y$-locus of the section $\sigma$.
If $A$ is a vector space then $\sigma^{*}$ can be identified with the identity map $K\left(P \times_{G} A\right) \simeq$ $K(P / G)=K(M)$.

### 8.3 GIT Quotients

With the applications in mind we use the following simple version of GIT quotient: Let $V$ be a $G$-vector space for a connected algebraic group $G$ and assume that $P \subset V$ is an open $G$-invariant subset such that $\pi: P \rightarrow P / G$ is a principal $G$-bundle over the smooth $M:=P / G$. (We want $P$ to be a right $G$-space so we define $p g:=g^{-1} p$.)

To state the main theorem we first introduce the K-theoretic Kirwan map $\kappa$ : $K_{G}(V) \rightarrow K(M)$ as the composition of the pull back $K\left(P \times_{G} V\right) \rightarrow K(M)$ via the zero section and the association map $K_{G}(V) \rightarrow K\left(P \times_{G} V\right)$ given in Equation (12) (cf. the RHS of (13)). Notice again that this definition is simpler than the cohomological analogue.

Theorem 8.9 Let $V$ be a $G$-vector space and assume that $P \subset V$ is an open $G$ invariant subset such that $\pi: P \rightarrow P / G$ is a principal $G$-bundle over the smooth $M:=P / G$. Let $Y \subset M$. Then

$$
\mathrm{mS}(Y, M)=\kappa\left(\mathrm{mS}_{G}\left(\pi^{-1}(Y), V\right)\right)
$$

To prove this theorem we apply Corollary 8.8 to the universal section: The inclusion $j: P \rightarrow V$ is obviously $G$-equivariant therefore induces a section $\sigma_{j}: M \rightarrow$ $P \times_{G} V$. In other words $\sigma_{j}(m)=[p, p]$ for any $p$ with $\pi(p)=m$.
Example 8.10 For $V=\operatorname{Hom}\left(\mathbb{C}^{k}, \mathbb{C}^{n}\right)$ with the $G=\mathrm{GL}_{k}(\mathbb{C})$-action we can choose $P:=\Sigma^{0}\left(\mathbb{C}^{k}, \mathbb{C}^{n}\right)$, the set of injective maps, so $P / G=\operatorname{Gr}_{k}\left(\mathbb{C}^{n}\right)$. Then $P \times_{G} V$ is the vector bundle $\operatorname{Hom}\left(\gamma^{k}, \mathbb{C}^{n}\right)$ for the tautological subbundle $\gamma^{k}$ and $\sigma_{j}$ is the section expressing the fact that $\gamma^{k}$ is a subbundle of the trivial bundle of rank $n$. More generally see Remark 5.3.

The reason we call $\sigma_{j}$ universal is the following simple observation: For any $Y \subset M$ the $\pi^{-1}(Y)$-locus of $\sigma_{j}$ is $Y$. Therefore Theorem 8.9 is a consequence of the following:

Proposition 8.11 The universal section $\sigma_{j}$ is motivically transversal to $P \times{ }_{G}$ $\pi^{-1}(Y)$ for any $Y \subset M$.

In fact we can replace $\pi^{-1}(Y)$ with any $G$-invariant (algebraic) subset $Z \subset V$.
Proof The statement is local so it is enough to study the restriction of $\sigma_{j}$ to an open subset of $M$ over which $P \times_{G} V$ is trivial. A local trivialization can be obtained by a $\varphi: W \rightarrow P$ local section of $P$ which is transversal to the fibers i.e. to the $G$-orbits. In this local trivialization the section $\left.\sigma_{j}\right|_{W}$ is the graph of the map $j \varphi: W \rightarrow V$, therefore $\sigma_{j}$ is motivically transversal to $P \times_{G} Z$ if $\varphi$ is motivically transversal to $Z \subset V$, which is implied by Proposition 8.6.

With the same argument we can obtain an equivariant version of Theorem 8.9:
Theorem 8.12 Let $V$ be a $G \times H$-vector space and assume that $P \subset V$ is an open $G \times H$-invariant subset such that $\pi: P \rightarrow P / G$ is a principal $G$-bundle over the smooth $M:=P / G$. Let $Y \subset M$ be $H$-invariant. Then

$$
\mathrm{mS}_{H}(Y, M)=\kappa_{H}\left(\mathrm{mS}_{G \times H}\left(\pi^{-1}(Y), V\right)\right),
$$

where $\kappa_{H}$ is the equivariant Kirwan map: the composition of the pull back $K_{H}\left(P \times_{G}\right.$ $V) \rightarrow K_{H}(M)$ via the zero section and the association map $K_{G \times H}(V) \rightarrow K_{H}\left(P \times_{G}\right.$ $V)$.

As a consequence we proved argument (e) in Sect. 5.4 (cf. Remark 5.3): the calculation of motivic Chern classes of matrix Schubert cells leads directly to the calculation of motivic Chern classes of ordinary Schubert cells. The formulas for motivic Chern classes are modified by the respective total Chern classes of the ambient spaces, while the formulas for the motivic Segre classes are identical (cf. Remark 5.5).

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# A Survey of Recent Developments on Hessenberg Varieties 

Hiraku Abe and Tatsuya Horiguchi


#### Abstract

This article surveys recent developments on Hessenberg varieties, emphasizing some of the rich connections of their cohomology and combinatorics. In particular, we will see how hyperplane arrangements, representations of symmetric groups, and Stanley's chromatic symmetric functions are related to the cohomology rings of Hessenberg varieties. We also include several other topics on Hessenberg varieties to cover recent developments.


Keywords Hessenberg varieties • Flag varieties • Cohomology • Hyperplane arrangements - Representations of symmetric groups • Chromatic symmetric functions.

## 1 Introduction

Hessenberg varieties are subvarieties of the full flag variety which was introduced by F. De Mari, C. Procesi, and M. A. Shayman [21, 22] around 1990. They provide a relatively new research subject, and similarly to Schubert varieties it has been found that geometry, combinatorics, and representation theory interact nicely on Hessenberg varieties. Let $X$ be a complex $n \times n$ matrix considered as a linear map $X: \mathbb{C}^{n} \rightarrow \mathbb{C}^{n}$ and $h:\{1,2, \ldots, n\} \rightarrow\{1,2, \ldots, n\}$ a Hessenberg function, i.e. a non-decreasing function satisfying $h(j) \geq j$ for $1 \leq j \leq n$. The Hessenberg variety (in type $A_{n-1}$ ) associated to $X$ and $h$ is defined as follows:

[^23]$$
\operatorname{Hess}(X, h):=\left\{V_{\bullet} \in F l\left(\mathbb{C}^{n}\right) \mid X V_{j} \subseteq V_{h(j)} \text { for all } 1 \leq j \leq n\right\}
$$
where $F l\left(\mathbb{C}^{n}\right)$ is the flag variety of $\mathbb{C}^{n}$ consisting of sequences $V_{\bullet}=\left(V_{1} \subset V_{2} \subset\right.$ $\cdots \subset V_{n}=\mathbb{C}^{n}$ ) of linear subspaces of $\mathbb{C}^{n}$ with $\operatorname{dim} V_{i}=i$ for $1 \leq i \leq n$. Particular examples include the full flag variety itself, Springer fibres, the Peterson variety, and the permutohedral variety.

Over the past 20 years, plentiful developments of Hessenberg varieties has been made. For example, it has been discovered that hyperplane arrangements and representations of symmetric groups appear when we deal with the cohomology rings of Hessenberg varieties. Also, these representations are determined by Stanley's chromatic symmetric functions of certain graphs, and is related with the StanleyStembridge conjecture in graph theory.

This article is a survey of recent developments on Hessenberg varieties, and it is intended to stimulate future research. We keep our explanations concise and include concrete examples so as to make the important ideas accessible, especially for young mathematicians (e.g. graduate students, postdoctoral fellows, and so on). We also include several other topics on Hessenberg varieties to cover recent developments. For simplicity, we explain most of the results in type $A$, but we make comments for results which hold in arbitrary Lie type.

## 2 Background and Notations

In this section, we recall some background, and establish some notations for the rest of the document.

### 2.1 Definitions and Basic Properties

Let $n$ be a positive integer, and we use the notation $[n]:=\{1,2, \ldots, n\}$ throughout this document. A function $h:[n] \rightarrow[n]$ is a Hessenberg function if it satisfies the following two conditions:
(i) $h(1) \leq h(2) \leq \cdots \leq h(n)$,
(ii) $h(j) \geq j$ for all $j \in[n]$.

Note that $h(n)=n$ by definition. We frequently write a Hessenberg function by listing its values in a sequence, i.e. $h=(h(1), h(2), \ldots, h(n))$. We may identify a Hessenberg function $h$ with a configuration of (shaded) boxes on a square grid of size $n \times n$ which consists of boxes in the $i$-th row and the $j$-th column satisfying $i \leq h(j)$ for $i, j \in[n]$, as we illustrate in the following example.

Example 2.1 Let $n=5$. The Hessenberg function $h=(3,3,4,5,5)$ corresponds to the configuration of the shaded boxes drawn in Fig. 1.

Fig. 1 The configuration corresponding to
$h=(3,3,4,5,5)$


In particular, this identification implies that the set of Hessenberg functions and the set of Dyck paths are in one-to-one correspondence. That is, the number of Hessenberg functions is the Catalan number:

$$
\#\{h:[n] \rightarrow[n]: \text { Hessenberg functions }\}=\frac{1}{n+1}\binom{2 n}{n} .
$$

There is a natural partial order $\subset$ on Hessenberg functions defined as follows. For any two Hessenberg functions $h$ and $h^{\prime}$, we define $h \subset h^{\prime}$ by

$$
h \subset h^{\prime} \Longleftrightarrow h(j) \leq h^{\prime}(j) \quad \forall j \in[n] .
$$

We use the symbol $\subset$ for this order since it corresponds to the inclusion of the configurations of boxes under the above visualization of Hessenberg functions.

The (full) flag variety $F l\left(\mathbb{C}^{n}\right)$ of $\mathbb{C}^{n}$ is the collection of nested linear subspaces $V_{\bullet}=\left(V_{1} \subset V_{2} \subset \cdots \subset V_{n}=\mathbb{C}^{n}\right)$ with $\operatorname{dim} V_{i}=i$ for $i \in[n]$. For an $n \times n$ matrix $X$ considered as a linear map $X: \mathbb{C}^{n} \rightarrow \mathbb{C}^{n}$ and a Hessenberg function $h:[n] \rightarrow[n]$, the Hessenberg variety ${ }^{1}$ (in type $A_{n-1}$ ) associated with $X$ and $h$ is defined as

$$
\begin{equation*}
\operatorname{Hess}(X, h)=\left\{V_{\bullet} \in F l\left(\mathbb{C}^{n}\right) \mid X V_{j} \subset V_{h(j)} \text { for all } j \in[n]\right\} \tag{2.1}
\end{equation*}
$$

References [21, 22]. If $X$ is the zero matrix or $h=(n, n, \ldots, n)$, then it is clear that $\operatorname{Hess}(X, h)=F l\left(\mathbb{C}^{n}\right)$ is the flag variety itself from the definition (2.1). If $X$ is nilpotent and $h=i d=(1,2, \ldots, n)$, then $\operatorname{Hess}(X, h)$ is called a Springer fiber which plays an important role in the geometric representation theory of the symmetric group $S_{n}[64,65]$.

Remark 2.2 The definition (2.1) can be rephrased in terms of the adjoint representation of $\mathrm{GL}(n, \mathbb{C})$. See [21] for the definition in arbitrary Lie type. Also, Goresky-Kottwitz-MacPherson [33, Sect. 2] considered more general Hessenberg varieties which are defined for arbitrary representations of reductive algebraic groups (cf. Chen-Vilonen-Xue [18, Sect. 2]).

[^24]As a general picture of Hessenberg varieties, we remark the following two properties. Suppose that one considers Hessenberg varieties for a fixed matrix $X$. Then, Hessenberg varieties preserves the inclusions:

$$
\begin{equation*}
h \subset h^{\prime} \Rightarrow \operatorname{Hess}(X, h) \subset \operatorname{Hess}\left(X, h^{\prime}\right) \tag{2.2}
\end{equation*}
$$

Also, if $g \in \operatorname{GL}(n, \mathbb{C})$, then we have an isomorphism

$$
\operatorname{Hess}(X, h) \cong \operatorname{Hess}\left(g X g^{-1}, h\right)
$$

by sending $V_{\bullet}$ to $g V_{\bullet}$. This implies that we may assume that $X$ is in a Jordan canonical form.
J. Tymoczko lays the foundation for the study of Hessenberg varieties as follows:

Theorem 2.3 ([72]) Every Hessenberg variety $\operatorname{Hess}(X, h)$ is paved by affines. ${ }^{2}$ In particular, the integral cohomology group of $\operatorname{Hess}(X, h)$ is torsion-free, and the odd-degree cohomology groups vanish.

Remark 2.4 This generalizes the work of Spaltenstein [63] for the Springer fibers and De Mari-Procesi-Shayman [21] for the regular semisimple Hessenberg varieties (which we will define below). For generalizations to arbitrary Lie type, see Precup [52] (cf. Tymoczko [73], De Mari-Procesi-Shayman [21], and De Concini-LusztigProcesi [20]).

By Theorem 2.3 we may denote the Poincaré polynomial of $\operatorname{Hess}(X, h)$ by

$$
\begin{equation*}
\operatorname{Poin}(\operatorname{Hess}(X, h), q):=\sum_{i=0}^{m} \operatorname{dim} H^{2 i}(\operatorname{Hess}(X, h) ; \mathbb{Q}) q^{i} \tag{2.3}
\end{equation*}
$$

where $m:=\operatorname{dim}_{\mathbb{C}} \operatorname{Hess}(X, h)$ and the variable $q$ stands for the grading with $\operatorname{deg}(q)=2$.

In this survey, we will focus on so-called regular nilpotent Hessenberg varieties and regular semisimple Hessenberg varieties which are particularly well-studied. As we will see in Sects. 3 and 4, their cohomology rings has an interesting relation between each other, and these are related with other research areas such as hyperplane arrangements, representation theory, and graph theory.

### 2.2 Regular Nilpotent Hessenberg Varieties

Let $N$ a regular nilpotent matrix of size $n \times n$, i.e. a nilpotent matrix with a single Jordan block. In Jordan canonical form, it is given by

[^25]Fig. 2 Decomposition of $h$ into $h^{(1)}$ and $h^{(2)}$

$h$

$$
N=\left(\begin{array}{rrrrr}
0 & 1 & & & \\
& 0 & 1 & & \\
& & \ddots & \ddots & \\
& & & 0 & 1 \\
& & & & 0
\end{array}\right)
$$

For a Hessenberg function $h:[n] \rightarrow[n], \operatorname{Hess}(N, h)$ is called a regular nilpotent Hessenberg variety. We have the following two main examples for this class of Hessenberg varieties. For $h=(n, n, \ldots, n)$, $\operatorname{Hess}(N, h)$ is the flag variety $F l\left(\mathbb{C}^{n}\right)$ itself. For $h=(2,3,4, \ldots, n, n)$, i.e. $h(j)=j+1$ for $1 \leq j<n$, $\operatorname{Hess}(N, h)$ is called the Peterson variety which is related to the quantum cohomology of partial flag varieties (cf. [14, 49, 56]). Here, $h=(n, n, \ldots, n)$ is the maximum Hessenberg function, and we may think that $h=(2,3,4, \ldots, n, n)$ gives the minimum one in the following sense.

If we have $h(j)=j$ for some $j<n$, then the Hessenberg function $h$ can be decomposed into two Hessenberg functions $h^{(1)}$ and $h^{(2)}$ of smaller sizes defined as follows:

$$
\begin{aligned}
h^{(1)} & :=(h(1), h(2), \ldots, h(j)), \\
h^{(2)} & :=(h(j+1)-j, h(j+2)-j, \ldots, h(n)-j) .
\end{aligned}
$$

Then, $\operatorname{Hess}(N, h)$ is decomposed as the product of regular nilpotent Hessenberg varieties associated with $h^{(1)}:[j] \rightarrow[j]$ and $h^{(2)}:[n-j] \rightarrow[n-j]$ of smaller sizes [24, Theorem 4.5]. In this sense, the condition $h(j) \geq j+1$ for any $j<n$ is essential, and for such Hessenberg functions, $h=(2,3,4, \ldots, n, n)$ is the minimum one. Hence, from (2.2), we may regard regular nilpotent Hessenberg varieties Hess ( $N, h$ ) as a (discrete) family of subvarieties of $F l\left(\mathbb{C}^{n}\right)$ connecting the Peterson variety and the flag variety itself (Fig. 2).

In the following theorem, we summarize some basic properties of $\operatorname{Hess}(N, h)$. For this purpose, let us prepare some notations. Given a Hessenberg function $h$, we denote by $S_{n}^{h}$ the subset of the $n$-th symmetric group $S_{n}$ defined as

$$
\begin{equation*}
S_{n}^{h}:=\left\{w \in S_{n} \mid w^{-1}(w(j)-1) \leq h(j) \text { for all } j \in[n]\right\} \tag{2.4}
\end{equation*}
$$

where we take by convention $w^{-1}(w(j)-1)=0$ whenever $w(j)-1=0$. The condition for $w \in S_{n}^{h}$ in (2.4) is exactly the condition that the permutation flag $^{3}$ associated with $w$ belongs to $\operatorname{Hess}(N, h)$, and this condition is equivalent to the condition that the Schubert cell $X_{w}^{\circ}$ intersects with $\operatorname{Hess}(N, h)$ [73]. That is, we have

$$
w \in S_{n}^{h} \Longleftrightarrow X_{w}^{\circ} \cap \operatorname{Hess}(N, h) \neq \emptyset
$$

Here, the dimension of the Schubert cell $X_{w}^{\circ}$ is equal to $\ell(w)$, which is the length of $w$. It is known from [73] that the dimension of the intersection $X_{w}^{\circ} \cap \operatorname{Hess}(N, h)$ is equal to

$$
\begin{equation*}
\ell_{h}(w):=\#\{(j, i) \mid 1 \leq j<i \leq n, w(j)>w(i), i \leq h(j)\} \tag{2.5}
\end{equation*}
$$

and that the intersections $X_{w}^{\circ} \cap \operatorname{Hess}(N, h)$ form a paving by affines of $\operatorname{Hess}(N, h)$. Combining the works of D. Anderson, E. Insko, B. Kostant, E. Sommers, J. Tymoczko, and A. Yong, we now give some basic properties of $\operatorname{Hess}(N, h)$.

Theorem 2.5 ([12, 47, 49, 62, 73]) Let $\operatorname{Hess}(N, h)$ be a regular nilpotent Hessenberg variety. Then the following hold.
(1) $\operatorname{Hess}(N, h)$ is irreducible, and it is singular in general.
(2) The (complex) dimension of $\operatorname{Hess}(N, h)$ is equal to $\sum_{j=1}^{n}(h(j)-j)$.
(3) The Poincaré polynomial (2.3) for $X=N$ has the following two types of expressions:

$$
\begin{align*}
\operatorname{Poin}(\operatorname{Hess}(N, h), q) & =\sum_{w \in S_{n}^{h}} q^{\ell_{h}(w)}  \tag{2.6}\\
& =\prod_{j=1}^{n}\left(1+q+q^{2}+\cdots+q^{h(j)-j}\right) \tag{2.7}
\end{align*}
$$

For generalizations of Theorem 2.5 to arbitrary Lie type, see $[10,52,53,57,62$, 73].

Note that the (complex) dimension of $\operatorname{Hess}(N, h)$ given in Theorem 2.5 (2) is equal to the number of boxes in which lie strictly below the diagonal under the identification of a Hessenberg function $h$ and a configuration of boxes.

Example 2.6 Let $h=(3,3,4,5,5)$. Then we have that $\operatorname{dim}_{\mathbb{C}} \operatorname{Hess}(N, h)=5$, which is the number of boxes which lie strictly below the diagonal (Fig. 3).

Example 2.7 Let $h=(2,3,3)$. Then $S_{3}^{h}=\{123,213,132,321\}$ where we use the standard one-line notation $w=w(1) w(2) \cdots w(n)$ for permutations in $S_{n}$ throughout this document. In particular, we have geometrically that

[^26]Fig. 3 Boxes which lie strictly below the diagonal


$$
\operatorname{Hess}(N, h) \cap X_{w}^{\circ}=\emptyset \Longleftrightarrow w=312,231 .
$$

Since $\ell_{h}(123)=0, \ell_{h}(213)=\ell_{h}(132)=1, \ell_{h}(321)=2$ for each permutation in $S_{3}^{h}$, we have

$$
\begin{equation*}
\operatorname{Poin}(\operatorname{Hess}(N, h), q)=1+2 q+q^{2}=(1+q)^{2} . \tag{2.8}
\end{equation*}
$$

### 2.3 Regular Semisimple Hessenberg Varieties

Let $S$ be a regular semisimple matrix of size $n \times n$. In Jordan canonical form, it is given by

$$
S=\left(\begin{array}{llll}
c_{1} & & & \\
& c_{2} & & \\
& & \ddots & \\
& & & c_{n}
\end{array}\right)
$$

where $c_{1}, c_{2} \ldots, c_{n}$ are mutually distinct complex numbers. For a Hessenberg function $h:[n] \rightarrow[n]$, $\operatorname{Hess}(S, h)$ is called a regular semisimple Hessenberg variety. It is known that the topology of $\operatorname{Hess}(S, h)$ is independent of the choices of the (distinct) eigenvalues, ${ }^{4}$ and hence one may think that the topology of $\operatorname{Hess}(S, h)$ only depends on $h$. Based on this fact, we have the following two main examples for this class of Hessenberg varieties. For $h=(n, n, \ldots, n), \operatorname{Hess}(S, h)$ is the flag variety $F l\left(\mathbb{C}^{n}\right)$ itself. For $h=(2,3,4, \ldots, n, n), \operatorname{Hess}(S, h)$ is called the permutohedral variety which is the toric variety associated with the fan consisting of the collection of Weyl chambers of the root system of type $A_{n-1}$ [21, Theorem 11]. Similarly to the case for $\operatorname{Hess}(N, h)$ in Sect. 2.2, the Hessenberg function $h=(2,3,4, \ldots, n, n)$ gives the minimum for this class of Hessenberg varieties as well in the following sense.

If we have $h(j)=j$ for some $j<n$, it is known that $\operatorname{Hess}(S, h)$ is not connected but equidimensional. In fact, all the connected components are isomorphic, and each component is decomposed as the product of regular semisimple Hessenberg varieties

[^27]of smaller sizes as in Sect. 2.2. See [70] for detail description of the connected components. Hence the condition $h(j) \geq j+1$ for any $j<n$ is essential, and we may regard regular semisimple Hessenberg varieties $\operatorname{Hess}(S, h)$ as a (discrete) family of subvarieties of the flag variety connecting the permutohedral variety and the flag variety itself. De Mari-Procesi-Shayman proved the following properties of $\operatorname{Hess}(S, h)$ (for arbitrary Lie type).

Theorem 2.8 ([21]) Let $\operatorname{Hess}(S, h)$ be a regular semisimple Hessenberg variety. Then the following hold.
(1) $\operatorname{Hess}(S, h)$ is smooth, and it is connected if and only if $h(j) \geq j+1$ for all $j<n$.
(2) The (complex) dimension of $\operatorname{Hess}(S, h)$ is equal to $\sum_{j=1}^{n}(h(j)-j)$.
(3) The Poincaré polynomial (2.3) for $X=S$ has the following expression:

$$
\operatorname{Poin}(\operatorname{Hess}(S, h), q)=\sum_{w \in S_{n}} q^{\ell_{h}(w)}
$$

where $\ell_{h}(w)$ is defined in (2.5).
From Theorem 2.5 (2) and Theorem 2.8 (2), we see that

$$
\operatorname{dim}_{\mathbb{C}} \operatorname{Hess}(N, h)=\operatorname{dim}_{\mathbb{C}} \operatorname{Hess}(S, h)=\sum_{j=1}^{n}(h(j)-j)
$$

Unlike the situation for $\operatorname{Hess}(N, h)$, any regular semisimple Hessenberg variety $\operatorname{Hess}(S, h)$ intersects with all the Schubert cells $X_{w}^{\circ}$. It is known that the dimension of the intersection $X_{w}^{\circ} \cap \operatorname{Hess}(S, h)$ is equal to $\ell_{h}(w)$ given in (2.5), and that all of the intersections $X_{w}^{\circ} \cap \operatorname{Hess}(S, h)$ form a paving by affines of $\operatorname{Hess}(S, h)$.

Example 2.9 Let $n=3$ and $h=(2,3,3)$. Since $\ell_{h}(123)=0, \ell_{h}(213)=\ell_{h}(132)=$ $\ell_{h}(231)=\ell_{h}(312)=1, \ell_{h}(321)=2$ for each permutation in $S_{3}$, the Poincaré polynomial of $\operatorname{Hess}(S, h)$ is given by

$$
\operatorname{Poin}(\operatorname{Hess}(S, h), q)=1+4 q+q^{2}
$$

## 3 Cohomology

In this section, we explain the structures of the cohomology rings of Hessenberg varieties, focusing on regular nilpotent Hessenberg varieties $\operatorname{Hess}(N, h)$ in Sect. 3.1 and regular semisimple Hessenberg varieties $\operatorname{Hess}(S, h)$ in Sect.3.2. We will also see that these two cohomology rings have an interesting relation in Sect.3.3.

### 3.1 Cohomology Rings of Regular Nilpotent Hessenberg Varieties

The cohomology ring of a regular nilpotent Hessenberg variety $\operatorname{Hess}(N, h)$ has been studied from various viewpoints (e.g. [6, 10, 15, 23, 29, 36, 41, 43, 45, 50]).

In this section we explain an explicit presentation of the cohomology ring of a regular nilpotent Hessenberg variety given by [6] due to M. Harada, M. Masuda, and the authors. We also discuss a relation between this presentation and Schubert polynomials along [42].

We first recall an explicit presentation of the cohomology ring of the flag variety $F l\left(\mathbb{C}^{n}\right)$. Let $E_{i}$ be the $i$-th tautological vector bundle over $F l\left(\mathbb{C}^{n}\right)$; namely, $E_{i}$ is the subbundle of the trivial vector bundle $F l\left(\mathbb{C}^{n}\right) \times \mathbb{C}^{n}$ over $F l\left(\mathbb{C}^{n}\right)$ whose fiber over a point $V_{\bullet}=\left(V_{1} \subset \cdots \subset V_{n}=\mathbb{C}^{n}\right) \in F l\left(\mathbb{C}^{n}\right)$ is exactly $V_{i}$. We denote the negative of the first Chern class of the quotient line bundle $E_{i} / E_{i-1}$ by $\bar{x}_{i}$, i.e.

$$
\begin{equation*}
\bar{x}_{i}:=-c_{1}\left(E_{i} / E_{i-1}\right) \in H^{2}\left(F l\left(\mathbb{C}^{n}\right) ; \mathbb{Q}\right) . \tag{3.1}
\end{equation*}
$$

We will also use the same symbol $\bar{x}_{i}$ for its restrictions to the cohomology of Hessenberg varieties by abuse of notation. It is known that there is a ring isomorphism

$$
\begin{equation*}
H^{*}\left(F l\left(\mathbb{C}^{n}\right) ; \mathbb{Q}\right) \cong \mathbb{Q}\left[x_{1}, \ldots, x_{n}\right] /\left(e_{1}, \ldots, e_{n}\right) \tag{3.2}
\end{equation*}
$$

which sends $x_{i}$ to $\bar{x}_{i}$ where $e_{i}$ is the elementary symmetric polynomial of degree $i$ in the variables $x_{1}, \ldots, x_{n}$ (cf. [30, p161, Proposition 3]).

In order to describe the cohomology ring of a regular nilpotent Hessenberg variety $\operatorname{Hess}(N, h)$, we define polynomials $f_{i, j}$ for $1 \leq j \leq i \leq n$ as follows:

$$
\begin{equation*}
f_{i, j}:=\sum_{k=1}^{j}\left(\prod_{\ell=j+1}^{i}\left(x_{k}-x_{\ell}\right)\right) x_{k} . \tag{3.3}
\end{equation*}
$$

Here, we take by convention $\prod_{\ell=j+1}^{i}\left(x_{k}-x_{\ell}\right)=1$ whenever $i=j$. Note that this definition does not involve $n$.

Theorem 3.1 ([6]) The restriction map

$$
H^{*}\left(F l\left(\mathbb{C}^{n}\right) ; \mathbb{Q}\right) \rightarrow H^{*}(\operatorname{Hess}(N, h) ; \mathbb{Q})
$$

is surjective, and there is a ring isomorphism

$$
\begin{equation*}
H^{*}(\operatorname{Hess}(N, h) ; \mathbb{Q}) \cong \mathbb{Q}\left[x_{1}, \ldots, x_{n}\right] /\left(f_{h(1), 1}, f_{h(2), 2}, \ldots, f_{h(n), n}\right) \tag{3.4}
\end{equation*}
$$

which sends $x_{i}$ to $\bar{x}_{i}=-\left.c_{1}\left(E_{i} / E_{i-1}\right)\right|_{\operatorname{Hess}(N, h)}$.

| $f_{1,1}$ |  |  |
| :--- | :--- | :--- |
| $f_{2,1}$ | $f_{2,2}$ |  |
| $f_{3,1}$ | $f_{3,2}$ | $f_{3,3}$ |$\quad$| $f_{1,1}=x_{1}, \quad f_{2,2}=x_{1}+x_{2}, \quad f_{3,3}=x_{1}+x_{2}+x_{3}$, |
| :--- |
| $f_{2,1}=\left(x_{1}-x_{2}\right) x_{1}, \quad f_{3,2}=\left(x_{1}-x_{3}\right) x_{1}+\left(x_{2}-x_{3}\right) x_{2}$, |
| $f_{3,1}=\left(x_{1}-x_{2}\right)\left(x_{1}-x_{3}\right) x_{1}$. |

Fig. 4 The polynomials $f_{i, j}$ for $1 \leq i \leq j \leq 3$

Remark 3.2 The presentation (3.4) does not hold for the integral coefficients. See [7, Remark 3].

From the presentation (3.4), we can see that the cohomology ring $H^{*}(\operatorname{Hess}(N, h)$; $\mathbb{Q})$ is a complete intersection since the number of generators of the polynomial ring $\mathbb{Q}\left[x_{1}, \ldots, x_{n}\right]$ is equal to the number of generators of the ideal $\left(f_{h(1), 1}, \ldots, f_{h(n), n}\right)$. This implies the following corollary.

Corollary 3.3 ([6]) $H^{*}(\operatorname{Hess}(N, h) ; \mathbb{Q})$ is a Poincaré duality algebra.
Note that a regular nilpotent Hessenberg variety is singular in general, but its cohomology is a Poincaré duality algebra. For arbitrary Lie type, the restriction map is surjective and Corollary 3.3 holds [10] (Fig. 4).

Example 3.4 Let $n=3$. We first assign polynomials $f_{i, j}$ to each box below the diagonal line.

If we take $h=(3,3,3)$, then we obtain from Theorem 3.1 an explicit presentation of the flag variety $\operatorname{Hess}(N, h)=F l\left(\mathbb{C}^{3}\right)$ :

$$
\begin{equation*}
H^{*}(\operatorname{Hess}(N, h) ; \mathbb{Q}) \cong \mathbb{Q}\left[x_{1}, x_{2}, x_{3}\right] /\left(f_{3,1}, f_{3,2}, f_{3,3}\right) \tag{3.5}
\end{equation*}
$$

where the polynomials $f_{3,1}, f_{3,2}, f_{3,3}$ in this presentation are obtained by taking the bottom one in each column (See Fig. 5). It is straightforward to verify that the ideals in (3.2) (with $n=3$ ) and (3.5) are the same. If we take $h^{\prime}=(2,3,3)$, then we also obtain an explicit presentation of the Peterson variety $\operatorname{Hess}\left(N, h^{\prime}\right)$ in $F l\left(\mathbb{C}^{3}\right)$ :

$$
\begin{equation*}
H^{*}\left(\operatorname{Hess}\left(N, h^{\prime}\right) ; \mathbb{Q}\right) \cong \mathbb{Q}\left[x_{1}, x_{2}, x_{3}\right] /\left(f_{2,1}, f_{3,2}, f_{3,3}\right) \tag{3.6}
\end{equation*}
$$

One may notice that the polynomial $f_{3,1}$ in (3.5) is replaced by $f_{2,1}$ in (3.6) when we changed the Hessenberg function from $h$ to $h^{\prime}$. More specifically, the polynomial $f_{2,1}$ does not vanish in $H^{*}(\operatorname{Hess}(N, h))$, but it does vanish in $H^{*}\left(\operatorname{Hess}\left(N, h^{\prime}\right)\right)$. This polynomial $f_{2,1}$ has the following expression as a linear combination of Schubert polynomials $\mathfrak{S}_{w}\left(w \in S_{n}\right)$ :

$$
\begin{equation*}
f_{2,1}=x_{1}^{2}-x_{1} x_{2}=\mathfrak{S}_{312}-\mathfrak{S}_{231} \tag{3.7}
\end{equation*}
$$

Fig. 5 The bottom $f_{i, j}$ 's for
$h=(3,3,3)$ and $h^{\prime}=(2,3,3)$


Fig. 6 The pictures of $h$ and $h^{\prime}$

since $\mathfrak{S}_{312}=x_{1}^{2}$ and $\mathfrak{S}_{231}=x_{1} x_{2}$. As seen in Example 2.7, we also have

$$
\operatorname{Hess}\left(N, h^{\prime}\right) \cap X_{w}^{\circ}=\emptyset \Longleftrightarrow w=312,231 .
$$

These permutations $w=312$, 231 are exactly the ones that appeared in (3.7).
In general, we have a similar interpretation of the presentation (3.4) in Theorem 3.1 by considering a smaller Hessenberg variety $\operatorname{Hess}(N, h) \supset \operatorname{Hess}\left(N, h^{\prime}\right)$ of codimension 1, as suggested by the above example (Fig. 6).

The unique difference in the generators of the ideal appearing in the presentation (3.4) for $H^{*}(\operatorname{Hess}(N, h) ; \mathbb{Q})$ and $H^{*}\left(\operatorname{Hess}\left(N, h^{\prime}\right) ; \mathbb{Q}\right)$ is the polynomial $f_{i-1, j}$. The second author showed that this polynomial can be written as an alternating sum of Schubert polynomials $\mathfrak{S}_{w}$ where the set of permutations $w$ appearing in this sum coincides with the set of minimal length permutations $w$ in $S_{n}$ satisfying

$$
\operatorname{Hess}(N, h) \cap X_{w}^{\circ} \neq \emptyset \text { and } \operatorname{Hess}\left(N, h^{\prime}\right) \cap X_{w}^{\circ}=\emptyset
$$

See [42] for details.

### 3.2 The Cohomology Rings of Regular Semisimple Hessenberg Varieties

Let $\operatorname{Hess}(S, h)$ be a regular semisimple Hessenberg variety. One of the most interesting feature of $\operatorname{Hess}(S, h)$ is the $S_{n}$-representation on its cohomology $H^{*}(\operatorname{Hess}(S, h)$; $\mathbb{C}$ ) constructed by J. Tymoczko [74]. She first constructed an $S_{n}$-representation on a torus equivariant cohomology $H_{T}^{*}(\operatorname{Hess}(S, h) ; \mathbb{C})$ via a combinatorial description
called the GKM-presentation, and she showed that this representation descends to the ordinary cohomology $H^{*}(\operatorname{Hess}(S, h) ; \mathbb{C})$. An alternative geometric construction via a monodromy action of the fundamental group of the space of regular semisimple matrices is explained in Brosnan-Chow [16]. Following the construction [74], N. Teff started to analyze this $S_{n}$-representation in [70, 71], and Shareshian-Wachs [59, 60] announced a beautiful conjecture on this representation using chromatic quasisymmetric functions. In this manuscript, we explain this representation along the construction due to Tymoczko.

Let $T \subset \mathrm{GL}(n, \mathbb{C})$ be the maximal torus consisting of diagonal elements of $\operatorname{GL}(n, \mathbb{C})$. The flag variety $\operatorname{Flag}\left(\mathbb{C}^{n}\right)$ has a natural action of $\operatorname{GL}(n, \mathbb{C})$, and hence the torus $T$ acts on $\operatorname{Flag}\left(\mathbb{C}^{n}\right)$ via its restriction. This $T$-action preserves the regular semisimple Hessenberg variety $\operatorname{Hess}(S, h)$ since all the elements of $T$ commute with the diagonal matrix $S$. It is known that $\operatorname{Hess}(S, h)$ contains all the $T$-fixed points of the flag variety $\operatorname{Flag}\left(\mathbb{C}^{n}\right)\left[21\right.$, Proposition 3] so that $\operatorname{Hess}(S, h)^{T}=\operatorname{Flag}\left(\mathbb{C}^{n}\right)^{T} \cong S_{n}$. Here, the last bijection corresponds $w \in S_{n}$ and the permutation flag associated with $w$.

Let $H_{T}^{*}(\operatorname{Hess}(S, h) ; \mathbb{C})$ be the $T$-equivariant cohomology of $\operatorname{Hess}(S, h)$. Recalling that $\operatorname{Hess}(S, h)$ has no odd-degree cohomology from Theorem 2.3, we can apply localization techniques for $T$-equivariant cohomology which we refer [32, 74] for details. As a conclusion, we obtain the so-called GKM presentation of $H_{T}^{*}(\operatorname{Hess}(S, h) ; \mathbb{C})$ as a subring of a direct sum of polynomial rings $\bigoplus_{w \in S_{n}}$ $\mathbb{C}\left[t_{1}, \ldots, t_{n}\right]$.

Proposition 3.5 ([74, Proposition 4.7]) The equivariant cohomology $H_{T}^{*}(\operatorname{Hess}(S, h) ; \mathbb{C})$ is isomorphic (as rings) to

$$
\left\{\begin{array}{l|l}
\alpha \in \bigoplus_{w \in S_{n}} \mathbb{C}\left[t_{1}, \ldots, t_{n}\right] & \begin{array}{c}
\alpha(w)-\alpha\left(w^{\prime}\right) \text { is divisible by } t_{w(i)}-t_{w(j)} \\
\text { if } w^{\prime}=w(j \text { i)for some } j<i \text { with } i \leq h(j)
\end{array} \tag{3.8}
\end{array}\right\}
$$

where $\alpha(w)$ is the $w$-component of $\alpha$ and $(j i) \in S_{n}$ is the transposition of $j$ and $i$.
We can visualize elements of the subring appearing in (3.8) in terms of the socalled GKM graph whose vertex set is $S_{n}$ and there is an edge between vertices $w, w^{\prime} \in S_{n}$ if there exists $1 \leq j<i \leq n$ with $i \leq h(j)$ satisfying $w^{\prime}=w(j i)$. Additionally, we equip such an edge with the data of the polynomial $\pm\left(t_{w(i)}-t_{w(j)}\right)$ (up to sign) arising in (3.8). This labeled graph is called the GKM graph of Hess(S,h), and we denote it by $\Gamma(h)$. In this language, the condition in (3.8) says that the collection of polynomials $(\alpha(w))_{w \in S_{n}}$ satisfies the following: if $w$ and $w^{\prime}$ are connected in $\Gamma(h)$ by an edge labeled by $t_{w(i)}-t_{w(j)}$, then the difference of the polynomials assigned for $w$ and $w^{\prime}$ must be divisible by the label.

Example 3.6 Let $n=3$. For $h=(3,3,3)$ and $h^{\prime}=(2,3,3)$, the corresponding GKM graphs are depicted in Fig. 7, where we use the one-line notation for each vertex.


Fig. 7 The GKM graphs $\Gamma(h)$ and $\Gamma\left(h^{\prime}\right)$


Fig. 8 Some elements of $H_{T}^{*}(\operatorname{Hess}(S, h))$


Fig. 9 Some elements of $H_{T}^{*}\left(\operatorname{Hess}\left(S, h^{\prime}\right)\right)$

For example, one can verify that the tuples of polynomials in Fig. 8 are elements of $H_{T}^{*}(\operatorname{Hess}(S, h) ; \mathbb{C})=H_{T}^{*}\left(F l\left(\mathbb{C}^{3}\right) ; \mathbb{C}\right)\left(\right.$ and hence of $H_{T}^{*}\left(\operatorname{Hess}\left(S, h^{\prime}\right) ; \mathbb{C}\right)$ ).

We use the symbols $\bar{x}_{1}, \bar{x}_{2}, \bar{x}_{3}$ by abuse of notation (cf. (3.1)) because of the reason which we will explain soon later (See Example 3.7 below). Also, the elements $\bar{y}_{1}, \bar{y}_{2}, \bar{y}_{3}$ given by the tuples of polynomials in Fig. 9 are elements of $H_{T}^{*}\left(\operatorname{Hess}\left(S, h^{\prime}\right) ; \mathbb{C}\right)$ but not of $H_{T}^{*}(\operatorname{Hess}(S, h) ; \mathbb{C})=H_{T}^{*}\left(F l\left(\mathbb{C}^{3}\right) ; \mathbb{C}\right)$.

In what follows, we identify $H_{T}^{*}(\operatorname{Hess}(S, h) ; \mathbb{C})$ and the presentation (3.8), and we do not distinguish them. For each $i=1, \ldots, n$, it is clear that a collection $\left(t_{i}\right)_{w \in S_{n}}$ lies in (3.8). For simplicity, we also write this element as $t_{i}$ by abuse of notation. Then the theory of $T$-equivariant cohomology also shows that there is a ring isomorphism

$$
\begin{equation*}
H^{*}(\operatorname{Hess}(S, h) ; \mathbb{C}) \cong H_{T}^{*}(\operatorname{Hess}(S, h) ; \mathbb{C}) /\left(t_{1}, \ldots, t_{n}\right) \tag{3.9}
\end{equation*}
$$

This means that one can study the ordinary cohomology ring $H^{*}(\operatorname{Hess}(S, h) ; \mathbb{C})$ from the equivariant cohomology ring $H_{T}^{*}(\operatorname{Hess}(S, h) ; \mathbb{C})$.

We now describe the $S_{n}$-action on $H_{T}^{*}(\operatorname{Hess}(S, h) ; \mathbb{C})$ constructed by Tymoczko [74]. For $v \in S_{n}$ and $\alpha=(\alpha(w)) \in \bigoplus_{w \in S_{n}} \mathbb{C}\left[t_{1}, \ldots, t_{n}\right]$, we define an element $v \cdot \alpha$ by the formula

$$
\begin{equation*}
(v \cdot \alpha)(w):=v \cdot \alpha\left(v^{-1} w\right) \text { for all } w \in S_{n} \tag{3.10}
\end{equation*}
$$

where $v \cdot f\left(t_{1}, \ldots, t_{n}\right)=f\left(t_{v(1)}, \ldots, t_{v(n)}\right)$ for $f\left(t_{1}, \ldots, t_{n}\right) \in \mathbb{C}\left[t_{1}, \ldots, t_{n}\right]$ in the right-hand side. This $S_{n}$-action preserves the subset (3.8), and hence it defines an $S_{n}$-action on the equivariant cohomology $H_{T}^{*}(\operatorname{Hess}(S, h) ; \mathbb{C})$ by Proposition 3.5. Since we have $v \cdot t_{i}=t_{v(i)}$ for the classes $t_{i}=\left(t_{i}\right)_{w \in S_{n}}$ defined above, the $S_{n}$ action on $H_{T}^{*}(\operatorname{Hess}(S, h) ; \mathbb{C})$ induces an $S_{n}$-action on the ordinary cohomology $H^{*}(\operatorname{Hess}(S, h) ; \mathbb{C})$ via (3.9). By construction this $S_{n}$-representation preserves the cup product of $H^{*}(\operatorname{Hess}(S, h) ; \mathbb{C})$.

Example 3.7 Let $n=3$ and $h=(3,3,3)$. For $\bar{x}_{1}, \bar{x}_{2}, \bar{x}_{3} \in H_{T}^{2}(\operatorname{Hess}(S, h) ; \mathbb{C})=$ $H_{T}^{2}\left(F l\left(\mathbb{C}^{3}\right) ; \mathbb{C}\right)$ given in Fig. 8, one can easily verify from the definition (3.10) that they are invariant under the $S_{3}$-action;

$$
w \cdot \bar{x}_{i}=\bar{x}_{i} \quad(1 \leq i \leq 3)
$$

for any $w \in S_{3}$. Under the isomorphism (3.9), it follows that these equivariant cohomology classes corresponds to $\bar{x}_{i}=-c_{1}\left(E_{i} / E_{i-1}\right) \in H^{2}\left(F l\left(\mathbb{C}^{3}\right) ; \mathbb{Q}\right)$ introduced in (3.1) which gives a justification for our notation. From (3.2) (or (3.5)), this means that the $S_{3}$-representation on $H^{*}\left(F l\left(\mathbb{C}^{3}\right) ; \mathbb{C}\right)$ is trivial, and the same claim holds for the case $H^{*}\left(F l\left(\mathbb{C}^{n}\right) ; \mathbb{C}\right)$ in general [74, Proposition 4.4].

Example 3.8 Let $n=3$ and $h=(2,3,3)$. For $\bar{y}_{1}, \bar{y}_{2}, \bar{y}_{3} \in H_{T}^{2}(\operatorname{Hess}(S, h) ; \mathbb{C})$ given in Fig. 9, one can also verify that these classes are naturally permuted by the $S_{3}$-action;

$$
w \cdot \bar{y}_{i}=\bar{y}_{w(i)} \quad(1 \leq i \leq 3)
$$

for any $w \in S_{3}$ from the definition (3.10).
Compared to the situation for $\operatorname{Hess}(N, h)$ (e.g. Theorem 3.1), the restriction $\operatorname{map} H^{*}\left(\operatorname{Flag}\left(\mathbb{C}^{n}\right) ; \mathbb{C}\right) \rightarrow H^{*}(\operatorname{Hess}(S, h) ; \mathbb{C})$ is not surjective in general. Hence, to describe the ring $H^{*}(\operatorname{Hess}(S, h) ; \mathbb{C})$ in terms of ring generators and relations among them, we need to find some cohomology classes of $\operatorname{Hess}(S, h)$ which do not come from $H^{*}\left(\operatorname{Flag}\left(\mathbb{C}^{n}\right) ; \mathbb{C}\right)$ by restriction. However, since we have a surjection $H_{T}^{*}(\operatorname{Hess}(S, h) ; \mathbb{C}) \rightarrow H^{*}(\operatorname{Hess}(S, h) ; \mathbb{C})$ by (3.9), the graphical presentation of the equivariant cohomology ring $H_{T}^{*}(\operatorname{Hess}(S, h) ; \mathbb{C})$ can be used to seek those classes as we saw in Example 3.6. For the case $h=(h(1), n, \ldots, n)$, M. Masuda and the authors showed that we can explicitly describe the integral cohomology ring in terms of ring generators and their relations by this approach. See [8] for details. For the case $h=(2,3,4, \ldots, n, n)$, the cohomology ring $H^{*}(\operatorname{Hess}(S, h) ; \mathbb{C})$ is well-understood since $\operatorname{Hess}(S, h)$ is a non-singular projective toric variety in this case [1, 48, 55]. However, the ring structure of $H^{*}(\operatorname{Hess}(S, h) ; \mathbb{C})$ for general $h$ is not well-understood at this moment.

Example 3.9 Let $n=3$ and $h=(2,3,3)$. Recall from Fig. 9 that we have three classes $\bar{y}_{1}, \bar{y}_{2}, \bar{y}_{3}$ in the equivariant cohomology $H_{T}^{2}(\operatorname{Hess}(S, h) ; \mathbb{C})$. We also denote by the same symbol $\bar{y}_{i} \in H^{2}(\operatorname{Hess}(S, h) ; \mathbb{C})$ the image of $\bar{y}_{i}$ under the isomorphism (3.9) by abuse of notation. Then the presentation of $H^{*}(\operatorname{Hess}(S, h) ; \mathbb{C})$ due to [8] is given by

$$
\begin{equation*}
H^{*}(\operatorname{Hess}(S, h) ; \mathbb{C}) \cong \mathbb{C}\left[x_{1}, x_{2}, x_{3}, y_{1}, y_{2}, y_{3}\right] / J, \tag{3.11}
\end{equation*}
$$

where $x_{i}$ and $y_{i}$ correspond to $\bar{x}_{i}=-\left.c_{1}\left(E_{i} / E_{i-1}\right)\right|_{\text {Hess }(S, h)}$ and $\bar{y}_{i}$ respectively, and the ideal $J$ is generated by

$$
\begin{aligned}
& y_{k} y_{k^{\prime}} \text { for } 1 \leq k \neq k^{\prime} \leq 3, \\
& x_{1} y_{k} \text { for } 1 \leq k \leq 3, \\
& x_{3} y_{k}+x_{2} x_{3} \text { for } 1 \leq k \leq 3, \\
& y_{1}+y_{2}+y_{3}-\left(x_{1}-x_{2}\right), \\
& x_{1}+x_{2}+x_{3}, x_{1} x_{2}+x_{1} x_{3}+x_{2} x_{3}, x_{1} x_{2} x_{3} .
\end{aligned}
$$

Also, the $S_{3}$-action on $H^{*}(\operatorname{Hess}(S, h) ; \mathbb{C})$ is given by $w \cdot x_{i}=x_{i}$ and $w \cdot y_{i}=y_{w(i)}$ for $i=1,2,3$ and $w \in S_{3}$, and the irreducible decomposition as $S_{3}$-representation is given by

$$
\begin{equation*}
H^{*}(\operatorname{Hess}(S, h) ; \mathbb{C}) \cong S^{(3)} \oplus\left(\left(S^{(3)}\right)^{\oplus 2} \oplus S^{(2,1)}\right) q \oplus S^{(3)} q^{2}, \tag{3.12}
\end{equation*}
$$

where $S^{\lambda}$ for a partition $\lambda$ of 3 is the irreducible representation of $S_{3}$ corresponding to $\lambda$ (cf. [30]), and $q$ is a formal symbol standing for the cohomology grading with $\operatorname{deg}(q)=2$. Note that $S^{(3)}$ is the trivial representation. In particular, this recovers the Poincaré polynomial of Hess ( $S, h$ ) in Example 2.9.

### 3.3 Regular Nilpotent Versus Regular Semisimple

In the last two sections, we have described the cohomology of regular nilpotent Hessenberg varieties $\operatorname{Hess}(N, h)$ and regular semisimple Hessenberg varieties $\operatorname{Hess}(S, h)$, and we saw that the latter cohomology $H^{*}(\operatorname{Hess}(S, h) ; \mathbb{C})$ admits the $S_{n}$-representation constructed by Tymoczko. Since this representation preserves the cup product, the invariant subgroup $H^{*}(\operatorname{Hess}(S, h) ; \mathbb{C})^{S_{n}}$ in fact forms a subring of $H^{*}(\operatorname{Hess}(S, h) ; \mathbb{C})$. M. Harada, M. Masuda, and the authors [6] showed that the $S_{n}$-representation provides a connection between the topology of $\operatorname{Hess}(N, h)$ and $\operatorname{Hess}(S, h)$ as follows.

Theorem 3.10 ([6]) Let $\operatorname{Hess}(N, h)$ and $\operatorname{Hess}(S, h)$ be a regular nilpotent Hessenberg variety and a regular semisimple Hessenberg variety, respectively. Then, there is a ring isomorphism

$$
H^{*}(\operatorname{Hess}(N, h) ; \mathbb{C}) \cong H^{*}(\operatorname{Hess}(S, h) ; \mathbb{C})^{S_{n}}
$$

which sends $\bar{x}_{i}=-\left.c_{1}\left(E_{i} / E_{i-1}\right)\right|_{\text {Hess }(N, h)}$ to $\bar{x}_{i}=-\left.c_{1}\left(E_{i} / E_{i-1}\right)\right|_{\text {Hess }(S, h)}$.

For the case $h=(n, n, \ldots, n)$, we have $\operatorname{Hess}(N, h)=\operatorname{Hess}(S, h)=F l\left(\mathbb{C}^{n}\right)$, and the isomorphism in Theorem 3.10 is obvious since the $S_{n}$-representation on the cohomology is trivial in this case (See Example 3.7). For the case $h=(2,3,4, \ldots, n, n)$, explicit presentations for the rings $H^{*}(\operatorname{Hess}(N, h) ; \mathbb{C})$ and $H^{*}(\operatorname{Hess}(S, h) ; \mathbb{C})^{S_{n}}$ were given in [29, 36] and [48] respectively, and those presentations are in fact identical although it was not mentioned. In this case, it means that the cohomology ring of the Peterson variety is isomorphic to the $S_{n}$-invariant subring of the cohomology of the permutohedral variety. See the work of P. Brosnan and T. Chow [16])for the geometry behind this phenomenon. We also note that Theorem 3.10 holds for arbitrary Lie type [10].

Example 3.11 Let $n=3$ and $h=(2,3,3)$. Then from (2.8) and (3.12), it is clear that we have the following equalities for dimension;

$$
\operatorname{dim} H^{2 k}(\operatorname{Hess}(N, h) ; \mathbb{C})=\operatorname{dim} H^{2 k}(\operatorname{Hess}(S, h) ; \mathbb{C})^{S_{3}}
$$

for all $k$. However, Theorem 3.10 states more; these are isomorphic as (graded) rings. For this case $h=(2,3,3)$, one can directly construct this isomorphism by using the ring presentations (3.6) and (3.11) (see [8, Remark 4.8.]).

## 4 Combinatorics

In this section, we explain combinatorial objects which are related to Hessenberg varieties. More specifically, we will see how hyperplane arrangements arise to describe the structure of the cohomology rings $H^{*}(\operatorname{Hess}(N, h) ; \mathbb{R})$ of regular nilpotent Hessenberg varieties and how Stanley's chromatic symmetric function of graphs determines the $S_{n}$-representation on the cohomology rings $H^{*}(\operatorname{Hess}(S, h) ; \mathbb{C})$ of regular semisimple Hessenberg varieties.

### 4.1 Hyperplane Arrangements

In this section we explain a connection between Hessenberg varieties and hyperplane arrangements established in [10]. Originally, E. Sommers and J. Tymoczko pointed out that Hessenberg varieties are related to hyperplane arrangements, and they conjectured that the Ponicaré polynomial of a regular nilpotent Hessenberg variety $\operatorname{Hess}(N, h)$ can be described in terms of certain hyperplane arrangement [62]. This conjecture was verified for some Lie types by Sommers-Tymoczko, G. Röhrle, A. Schauenburg [57, 62]). After this, an explicit presentation of the cohomology ring of $\operatorname{Hess}(N, h)$ was provided by M. Harada, M. Masuda, and the authors [6].

Fig. 10 The Weyl arrangement for $n=3$


Fig. 11 The ideal arrangement associated with $h=(2,3,3)$


Motivated by all these works, T. Abe, M. Masuda, S. Murai, T. Sato, and the second author proved that the cohomology ring of a regular nilpotent Hessenberg variety is isomorphic to a ring coming from those hyperplane arrangements [10]. By taking Hilbert series of both sides of this isomorphism, one sees that the conjecture of Sommers and Tymoczko is true. For general reference about hyperplane arrangements, see [51] (Figs. 10 and 11).

Let $V$ be a real vector space of finite dimension. A (central) hyperplane arrangement $\mathcal{A}$ in $V$ is a finite set of linear hyperplanes in $V$. As we see in the following example, Hessenberg functions naturally determine hyperplane arrangements.

Example 4.1 Let $V=\left\{\left(x_{1}, \ldots, x_{n}\right) \in \mathbb{R}^{n} \mid x_{1}+\cdots+x_{n}=0\right\}$, and consider hyperplanes in $V$ given by $H_{i, j}=\left\{x_{j}-x_{i}=0\right\}$ for $1 \leq j<i \leq n$. For each Hessenberg function $h:[n] \rightarrow[n]$, the set

$$
\mathcal{A}_{h}:=\left\{H_{i, j} \mid 1 \leq j<i \leq h(j)\right\}
$$

is called an ideal arrangement associated with $h$. In particular, if $h=(n, n, \ldots, n)$, then $\mathcal{A}_{h}$ is called the Weyl arrangement (of type $A_{n-1}$ ).

Let $\mathcal{R}=\operatorname{Sym}\left(V^{*}\right)$ be the symmetric algebra of $V^{*}$, where $V^{*}$ is the dual space of $V$. We regard $\mathcal{R}$ as an algebra of polynomial functions on $V$. A map $\theta: \mathcal{R} \rightarrow \mathcal{R}$ is an $\mathbb{R}$-derivation if it satisfies
(1) $\theta$ is $\mathbb{R}$-linear,
(2) $\theta(f \cdot g)=\theta(f) \cdot g+f \cdot \theta(g)$ for all $f, g \in \mathcal{R}$.

We denote the module of $\mathbb{R}$-derivations $\theta: \mathcal{R} \rightarrow \mathcal{R}$ by $\operatorname{Der} \mathcal{R}$. If one chooses a linear coordinate system $x_{1}, \ldots, x_{m}$ on $V$, i.e., $x_{1}, \ldots, x_{m}$ is a basis for $V^{*}$, then the $\mathcal{R}$-module $\operatorname{Der} \mathcal{R}$ can be expressed as $\bigoplus_{i=1}^{m} \mathcal{R} \frac{\partial}{\partial x_{i}}$ where $\frac{\partial}{\partial x_{i}}$ denotes the partial derivative with respect to $x_{i}$.

Fig. 12 A basis of the logarithmic derivation module of the hyperplane arrangement $\{x=0\}$ in $\mathbb{R}^{2}$

the vector field $x \frac{\partial}{\partial x}$

the vector field $\frac{\partial}{\partial y}$

Let $\mathcal{A}$ be a hyperplane arrangement in $V$. For each $H \in \mathcal{A}$, let $\alpha_{H} \in V^{*}$ be the defining linear form of $H$ so that $H=\operatorname{ker}\left(\alpha_{H}\right)$. The logarithmic derivation module $D(\mathcal{A})$ of $\mathcal{A}$ is an $\mathcal{R}$-module defined by

$$
D(\mathcal{A}):=\left\{\theta \in \operatorname{Der} \mathcal{R} \mid \theta\left(\alpha_{H}\right) \in \mathcal{R} \alpha_{H} \text { for all } H \in \mathcal{A}\right\}
$$

Geometrically, this consists of polynomial vector fields on $V$ tangent to $\mathcal{A}$.
Example 4.2 Let $V=\mathbb{R}^{2}$ and $H=\{x=0\}$ a hyperplane. We consider a hyperplane arrangement $\mathcal{A}=\{H\}$. An element $\theta \in \operatorname{Der} \mathcal{R}$ can be written as

$$
\theta=f \frac{\partial}{\partial x}+g \frac{\partial}{\partial y}
$$

for some $f, g \in \mathcal{R}=\mathbb{R}[x, y]$. Since we can take $\alpha_{H}=x$, we have $\theta\left(\alpha_{H}\right)=\theta(x)=$ $f$. Hence, $\theta$ belongs to the logarithmic derivation module $D(\mathcal{A})$ if and only if $f$ is divisible by $x$. From this, it is straightforward to see that $D(\mathcal{A})$ is free module over $\mathbb{R}[x, y]$ with a basis $x \frac{\partial}{\partial x}$ and $\frac{\partial}{\partial y}$. Hence, $D(\mathcal{A})$ consists of polynomial vector fields on $\mathbb{R}^{2}$ tangent to $\mathcal{A}$, as desired (Fig. 12).

In Example 4.2, we saw that the logarithmic derivation module $D(\mathcal{A})$ is a free $\mathcal{R}$ module. However, the logarithmic derivation module $D(\mathcal{A})$ of a hyperplane arrangement $\mathcal{A}$ is not free in general (cf. [51, Example 4.34]). A hyperplane arrangement $\mathcal{A}$ is a free arrangement if its logarithmic derivation module $D(\mathcal{A})$ is a free module over $\mathcal{R}$. Note that if $\mathcal{A}$ is a free arrangement in $V$, then $D(\mathcal{A})$ has a basis consisting of $m$ homogeneous elements where $m:=\operatorname{dim} V$ (cf. [51, Proposition 4.18]). Here, a nonzero element $\theta \in \operatorname{Der} \mathcal{R}=\mathcal{R} \otimes V$ is homogeneous if $\theta=\sum_{k=1}^{\ell} f_{k} \otimes v_{k}$ ( $f_{k} \in \mathcal{R}, v_{k} \in V$ ) and all non-zero $f_{k}$ 's are homogeneous polynomials of the same degree.

We now explain the connection with the cohomology rings of regular nilpotent Hessenberg varieties. Let $h$ be a Hessenberg function and $\mathcal{A}_{h}$ the ideal arrangement in $V=\left\{\left(x_{1}, \ldots, x_{n}\right) \in \mathbb{R}^{n} \mid x_{1}+\cdots+x_{n}=0\right\}$ given in Example 4.1. First, it is known to be free (for arbitrary Lie type) by T. Abe, M. Barakat, M. Cuntz, T. Hoge, and H. Terao.

Theorem 4.3 ([9]) Any ideal arrangement $\mathcal{A}_{h}$ is free.

We next define an ideal of $\mathcal{R}$ from the logarithmic derivation module $D\left(\mathcal{A}_{h}\right)$ as follows. Let $Q$ be an $S_{n}$-invariant non-degenerate quadratic form on $V$, which is unique up to a non-zero scalar multiple. We may take $Q=x_{1}^{2}+\cdots+x_{n}^{2} \in \operatorname{Sym}^{2}\left(V^{*}\right)^{S_{n}}$. We define an ideal $\mathfrak{a}(h)$ of $\mathcal{R}$ by

$$
\mathfrak{a}(h):=\left\{\theta(Q) \in \mathcal{R} \mid \theta \in D\left(\mathcal{A}_{h}\right)\right\} .
$$

T. Abe, M. Masuda, S. Murai, T. Sato, and the second author proved the following.

Theorem 4.4 ([10]) There is a ring isomorphism

$$
H^{*}(\operatorname{Hess}(N, h) ; \mathbb{R}) \cong \mathcal{R} / \mathfrak{a}(h)
$$

which sends $x_{i}$ to $\bar{x}_{i}=-\left.c_{1}\left(E_{i} / E_{i-1}\right)\right|_{\text {Hess }(N, h)}$ of (3.1).
In arbitrary Lie type, Theorem 4.4 is proved in [10]. It is known that the Poincare polynomial of $\operatorname{Hess}(N, h)$ has a summation formula such as (2.6) [52]. On the other hand, the Hilbert series of the quotient $\operatorname{ring} \mathcal{R} / \mathfrak{a}(h)$ has a product formula such as (2.7). Theorem 4.4 gives an affirmative answer to a conjecture of Sommers and Tymoczko in [62] which states that these formulas are equal.

Remark 4.5 The quotient ring in the right-hand side of Theorem 4.4 is an example of Solomon-Terao algebras studied in [28] and [11], where [28] considered more general hypersurface singularities. The Solomon-Terao algebra $S T(\mathcal{A}, \eta)$ is defined by a hyperplane arrangement $\mathcal{A}$ and a homogeneous polynomial $\eta$. Motivated by the work of L. Solomon and H. Terao in [61], it is proved in [11] that $S T(\mathcal{A}, \eta)$ for a generic $\eta$ is a complete intersection if and only if the hyperplane arrangement $\mathcal{A}$ is free. In [28] the same equivalence is proved more generally for hypersurface singularities which are holonomic in the sense of K. Saito [58, (3.8)]. This generalization was obtained independently and published slightly earlier.

Theorem 4.4 also tells us that if we can find an explicit $\mathcal{R}$-basis of $D(\mathcal{A})$, then we obtain an explicit presentation of the cohomology ring of $\operatorname{Hess}(N, h)$. In fact, we can recover the presentation (3.4) as follows. First, we recall a well-known criterion for bases of logarithmic derivation modules.

Theorem 4.6 (Saito's criterion, [58], see also [51]) Let $\mathcal{A}$ be a hyperplane arrangement in $V$ and let $\theta_{1}, \ldots, \theta_{m} \in D(\mathcal{A})$ be homogeneous derivations. Then $\theta_{1}, \ldots, \theta_{m}$ form an $\mathcal{R}$-basis of $D(\mathcal{A})$ if and only if $\theta_{1}, \ldots, \theta_{m}$ are $\mathcal{R}$-independent and $\sum_{i=1}^{m} \operatorname{deg} \theta_{i}=|\mathcal{A}|$.

To recover the presentation (3.4) by Saito's criterion, we define derivations $\psi_{i, j}$ for $1 \leq j \leq i \leq n$ on the ambient vector space $V=\left\{\left(x_{1}, \ldots, x_{n}\right) \in \mathbb{R}^{n} \mid x_{1}+\cdots+\right.$ $\left.x_{n}=0\right\}$ of ideal arrangements as follows:

$$
\psi_{i, j}:=\sum_{k=1}^{j}\left(\prod_{\ell=j+1}^{i}\left(x_{k}-x_{\ell}\right)\right)\left(\frac{\partial}{\partial x_{k}}-\frac{1}{n} \bar{\partial}\right) \in \operatorname{Der} \mathcal{R}
$$

where $\bar{\partial}:=\frac{\partial}{\partial x_{1}}+\cdots+\frac{\partial}{\partial x_{n}}$ and we take by convention $\prod_{\ell=j+1}^{i}\left(x_{k}-x_{\ell}\right)=1$ whenever $i=j$. Note that $\frac{\partial}{\partial x_{k}}$ is not an element of $\operatorname{Der\mathcal {R}}$ but $\left(\frac{\partial}{\partial x_{k}}-\frac{1}{n} \bar{\partial}\right)$ is, since $\mathcal{R}=$ $\mathbb{R}\left[x_{1}, \ldots, x_{n}\right] /\left(x_{1}+\cdots+x_{n}\right)$. Using Theorem 4.6 , one can verify that $\left\{\psi_{h(j), j} \mid 1 \leq\right.$ $j \leq n-1\}$ form an $\mathcal{R}$-basis of $D\left(\mathcal{A}_{h}\right)$ [10, Proposition 10.3]. Since $\psi_{i, j}(Q)=$ $\psi_{i, j}\left(x_{1}^{2}+\cdots+x_{n}^{2}\right)=2 f_{i, j}$ in $\mathcal{R}$, we recover the presentation (3.4) from Theorem 4.4.

Abe, Barakat, Cuntz, Hoge, and Terao in [9] gave a theoretical method to construct a basis of $D\left(\mathcal{A}_{h}\right)$ over $\mathcal{R}$ for arbitrary Lie type by classification-free proof. For a construction of an explicit basis for each Lie type, Barakat, Cuntz, and Hoge provided ones for types $E$ and $F$ by computer when the work of [9] was in progress. Also, Terao and Abe worked for types $A$ and $B$, respectively. In [10], an explicit basis was constructed for types $A, B, C, G$. Motivated by this, Enokizono, Nagaoka, Tsuchiya, and the second author in [26] introduced and studied uniform bases for the logarithmic derivation modules of the ideal arrangements. In particular, from Theorem 4.4 which is valid for all Lie types, we obtained explicit presentations of the cohomology rings of regular nilpotent Hessenberg varieties in all Lie types.

### 4.2 Chromatic Symmetric Functions in Graph Theory

In Sect. 3.2, we explained the representation of $S_{n}$ on the cohomology ring $H^{*}(\operatorname{Hess}(S, h) ; \mathbb{C})$ of regular semisimple Hessenberg varieties constructed by J. Tymoczko. Recall that there is a natural correspondence between representations of $S_{n}$ and symmetric functions of degree $n$ (cf. [30, Sect. 7]). In this section, we describe the symmetric function corresponding to the $S_{n}$-representation on $H^{*}(\operatorname{Hess}(S, h) ; \mathbb{C})$ which was conjectured by J. Shareshian and M. Wachs [59, Conjecture 1.2], [60, Conjecture 1.4] in terms of chromatic symmetric functions of a graph determined by $h$. This conjecture is proved by P. Brosnan and T. Chow in [16] and soon after by M. Guay-Paquet [35]. We refer [30] for notations and elementary knowledge of symmetric functions.

Let $h:[n] \rightarrow[n]$ be a Hessenberg function and $G_{h}$ a graph on the vertex set $[n]$ defined as follows; there is an edge between the vertices $i, j \in[n]$ if $j<i \leq h(j)$. Namely, if we visualize $h$ as a configuration of boxes as in Example 2.1, we have an edge between $i$ and $j$ if and only if we have a box in $(i, j)$-th position which is strictly below the diagonal (cf. Fig.3).

Example 4.7 If we take $h=(n, n, \ldots, n)$, then the graph $G_{h}$ is the complete graph ${ }^{5}$ on the vertex set $[n]$. If we take $h=(2,3,4, \ldots, n, n)$, then the graph $G_{h}$ is the path graph on the vertex set $[n]$ (Fig. 13).

[^28]Fig. 13 The complete graph on the vertex set [3]


For a graph $G=(V, E)$ on the vertex set $V=[n]$, Shareshian-Wachs [59, 60] introduced the chromatic quasisymmetric function $X_{G}(\mathbf{x}, q)$ of $G$ as

$$
\begin{equation*}
X_{G}(\mathbf{x}, q)=\sum_{\kappa}\left(\prod_{i=1}^{n} x_{\kappa(i)}\right) q^{\operatorname{asc}(\kappa)} \tag{4.1}
\end{equation*}
$$

where the summation runs over all proper colorings ${ }^{6} \kappa:[n] \rightarrow \mathbb{N}=\{1,2,3, \ldots\}$ of $G$ and $\operatorname{asc}(\kappa):=|\{(j, i) \in E \mid j<i, \kappa(j)<\kappa(i)\}|$ is the number of ascents of $\kappa$. Here, the variable $q$ stands for the grading, and this is a graded version of Stanley's chromatic symmetric function $X_{G}(\mathbf{x})$ of $G$ [67]. In general, $X_{G}(\mathbf{x}, q)$ is quasisymmetric in $\mathbf{x}$-variables but may not be symmetric. However, for our graph $G_{h}$, it is known that $X_{G_{h}}(\mathbf{x}, q)$ is in fact symmetric [60, Theorem 4.5]. In [60, Theroem 6.3], the Schur basis expansion of $X_{G_{h}}(\mathbf{x}, q)$ is determined in terms of combinatorics, where the non-graded version was originally obtained by V. Gasharov [31].

Example 4.8 Let $n=3$ and $h=(2,3,3)$. Then, one can verify that

$$
\begin{equation*}
X_{G_{h}}(\mathbf{x}, q)=s_{(1,1,1)}(\mathbf{x})+\left(2 s_{(1,1,1)}(\mathbf{x})+s_{(2,1)}(\mathbf{x})\right) q+s_{(1,1,1)}(\mathbf{x}) q^{2} . \tag{4.2}
\end{equation*}
$$

Here, $s_{\lambda}(\mathbf{x})$ is the Schur function corresponding to a partition $\lambda$ of 3 .
The following theorem determines the $S_{n}$-representation on $H^{*}(\operatorname{Hess}(S, h) ; \mathbb{C})$ in terms of the combinatorics of the graph $G_{h}$. This beautiful theorem was conjectured by Shareshian-Wachs [59, Conjecture 1.2], [60, Conjecture 1.4], and it was proved by Brosnan-Chow [16] and soon after by Guay-Paquet [35]. We denote by ch the Frobenius character under which symmetric functions of degree $n$ corresponds to representations of $S_{n}$.

Theorem 4.9 ( $[16,35])$ Let $\operatorname{Hess}(S, h)$ be a regular semisimple Hessenberg variety. Then,

$$
\omega X_{G_{h}}(\boldsymbol{x}, q)=\sum_{k=0}^{m} \operatorname{ch} H^{2 k}(\operatorname{Hess}(S, h) ; \mathbb{C}) q^{k}
$$

where $m=\operatorname{dim}_{\mathbb{C}} \operatorname{Hess}(S, h)$ and $\omega$ is the involution on the ring of symmetric functions in $\boldsymbol{x}$-variables sending each Schur function $s_{\lambda}(\boldsymbol{x})$ to $s_{\tilde{\lambda}}(\boldsymbol{x})$ associated with the transpose $\tilde{\lambda}$.

[^29]For the case $h=(n, n, \ldots, n)$, recall from Sect. 3.2 that the $S_{n}$-representation on $H^{*}(\operatorname{Hess}(S, h) ; \mathbb{C})=H^{*}\left(F l\left(\mathbb{C}^{n}\right) ; \mathbb{C}\right)$ is trivial, and it is straightforward to see that $X_{G_{h}}(\mathbf{x}, q)=e_{n}(\mathbf{x}) \sum_{\kappa} q^{\operatorname{asc}(\kappa)}$ in this case where the latter sum is equal to the Poincaré polynomial of $F l\left(\mathbb{C}^{n}\right)$. Since the complete symmetric function $h_{n}(\mathbf{x})=\omega\left(e_{n}(\mathbf{x})\right)$ corresponds to the 1-dimensional trivial representation, the equality in Theorem 4.9 holds in this case. For the case $h=(2,3,4, \ldots, n, n)$, both of $X_{G_{h}}(\mathbf{x}, q)$ and the $S_{n}$-representation on $H^{*}(\operatorname{Hess}(S, h) ; \mathbb{C})$ were well-studied in [55, 66, 69], and the equality was known. See [60, Sect. 1] for details. The proof of Theorem 4.9 given by Brosnan-Chow is geometric in the sense that they used the theory of monodromy actions, whereas Guay-Paquet provided a combinatorial proof using the theory of Hopf algebras.

Example 4.10 Let $n=3$ and $h=(2,3,3)$. Then, from (3.12) and (4.2), we see that

$$
\begin{aligned}
\omega X_{G_{h}}(\mathbf{x}, q) & =s_{(3)}(\mathbf{x})+\left(2 s_{(3)}(\mathbf{x})+s_{(2,1)}(\mathbf{x})\right) q+s_{(3)}(\mathbf{x}) q^{2} \\
& =\operatorname{ch} H^{0}(\operatorname{Hess}(S, h) ; \mathbb{C})+\operatorname{ch} H^{2}(\operatorname{Hess}(S, h) ; \mathbb{C}) q+\operatorname{ch} H^{4}(\operatorname{Hess}(S, h) ; \mathbb{C}) q^{2}
\end{aligned}
$$

Theorem 4.9 is related with the Stanley-Stembridge conjecture in graph theory as we explain in what follows. Recall that the chromatic quasisymmetric function $X_{G}(\mathbf{x}, q)$ is a graded version of Stanley's chromatic symmetric function $X_{G}(\mathbf{x})$ [67] which we obtain by forgetting the grading parameter $q$ in (4.1). Given a poset $P$, we can construct its incomparability graph which has as its vertices the elements of $P$ and we have an edge between two vertices if they are not comparable in $P$. The Stanley-Stembridge conjecture is stated as follows.

Conjecture 4.11 (Stanley-Stembridge conjecture [68, Conjecture 5.5], [67, Conjecture 5.1]) Let $G$ be the incomparability graph of a $(3+1)$-free poset. Then $X_{G}(\mathbf{x})$ is $e$-positive.

Here, a poset is called $(r+s)$-free if the poset does not contain an induced subposet isomorphic to the direct sum of an $r$ element chain and an $s$ element chain, and a symmetric function of degree $n$ is called $e$-positive if it is a non-negative sum of elementary symmetric functions $e_{\lambda}(\mathbf{x})$ for partition $\lambda$ of $n$. See [60], for more information on Stanley-Stembridge conjecture.

Example 4.12 Take a poset on [3]=\{1,2,3\} for which $1<3$ is the unique comparable pair. It is obviously ( $3+1$ )-free. In this case, the incomparability graph $G$ is the path graph given in Fig. 14. Hence, from (4.2), we see that

$$
\begin{aligned}
X_{G}(\mathbf{x}) & =s_{(1,1,1)}(\mathbf{x})+\left(2 s_{(1,1,1)}(\mathbf{x})+s_{(2,1)}(\mathbf{x})\right)+s_{(1,1,1)}(\mathbf{x}) \\
& =3 e_{(3)}(\mathbf{x})+e_{(2,1)}(\mathbf{x})
\end{aligned}
$$

as desired since we have $e_{(3)}(\mathbf{x})=s_{(1,1,1)}(\mathbf{x})$ and $e_{(2,1)}(\mathbf{x})=s_{(2,1)}(\mathbf{x})+s_{(1,1,1)}(\mathbf{x})$.

Fig. 14 The path graph on the vertex set [3]


Guay-Paquet [34] showed that, to solve the Stanley-Stembridge conjecture, it is enough to solve it for posets which are $(3+1)$-free and $(2+2)$-free, and this is precisely the case that the corresponding incomparability graph can be identified with $G_{h}$ for some Hessenberg function $h:[n] \rightarrow[n]$, where $n$ is the number of the elements of the poset, as we encountered in Example 4.12 (See [59, Sect.4] for details). Thus, the Stanley-Stembridge conjecture is now reduced to the following conjecture by Theorem 4.9.

Conjecture 4.13 ([59, Conjecture 5.4], [60, Conjecture 10.4]) The $S_{n}$-representation on $H^{*}(\operatorname{Hess}(S, h) ; \mathbb{C})$ constructed by Tymoczko is a permutation representation, i.e. a direct sum of induced representations of the trivial representation from $S_{\lambda}$ to $S_{n}$ where $S_{\lambda}=S_{\lambda_{1}} \times \cdots \times S_{\lambda_{\ell}}$ for $\lambda=\left(\lambda_{1}, \ldots, \lambda_{\ell}\right)$.

Motivated by the connection to Stanely-Stembridge conjecture, M. Harada and M. Precup verified Conjecture 4.11 for so-called abelian regular semisimple Hessenberg varieties, and they also derived a set of linear relations satisfied by the multiplicities of certain permutation representations. See [39, 40] for details.

## 5 More Topics

### 5.1 Semisimple Hessenberg Varieties

E. Insko and M. Precup studied Hessenberg varieties associated with semisimple matrices which may not be regular semisimple [44]. They determined the irreducible components of semisimple Hessenberg varieties for $h=(2,3,4, \ldots, n, n)$ in arbitrary Lie type. They also proved that irreducible components are smooth and gave an explicit description of their intersections.

### 5.2 Hessenberg Varieties for the Minimal Nilpotent Orbit

P. Crooks and the first author studied Hessenberg varieties associated with minimal nilpotent matrices (i.e. nilpotent matrices with Jordan blocks of size 1 and a single Jordan block of size 2). Explicit descriptions of their Poincaré polynomials and irreducible components are described, and a certain presentation of their cohomology rings are also provided in terms of Schubert classes. See [3] for details.

### 5.3 Regular Hessenberg Varieties

An $n \times n$ matrix $R$ is called regular if the Jordan blocks of $R$ have distinct eigenvalues. For a regular matrix $R, \operatorname{Hess}(R, h)$ is called a regular Hessenberg variety. This class of Hessenberg varieties contains regular nilpotent Hessenberg varieties and regular semisimple Hessenberg varieties, and $\operatorname{Hess}(R, h)$ plays an important role in the work of P. Brosnan and T. Chow [16] proving Shareshian-Wachs conjecture (Theorem 4.9). They are singular varieties in general, however M. Precup proved that the Betti numbers of $\operatorname{Hess}(R, h)$ are palindromic [53]. N. Fujita, H. Zeng, and the first author proved that higher cohomology groups of their structure sheaves vanish and that they degenerate to the regular nilpotent Hessenberg variety $\operatorname{Hess}(N, h)$ if $h(i) \geq i+1$ for all $1 \leq i \leq n$ [5], which was motivated by the works of D. Anderson and J. Tymoczo [12] and L. DeDieu, F. Galetto, M. Harada, and the first author [2].

### 5.4 Poincaré Dual of Hessenberg Varieties

In [5], the Poincaré dual of a regular Hessenberg variety $\operatorname{Hess}(R, h)$ in $H^{*}\left(F l\left(\mathbb{C}^{n}\right)\right)$ was computed in terms of positive roots associated the Hessenberg function $h$, and E. Insko, J. Tymoczko, and A. Woo gave a combinatorial formula for this class using Schubert polynomials [46]. Also, the cohomology class [Hess $(R, h)] \in H^{*}\left(F l\left(\mathbb{C}^{n}\right)\right)$ does not depend on a choice of a regular matrix $R$ if $h(i) \geq i+1$ for all $1 \leq i \leq n$ (See for details [5, 46]).

### 5.5 Additive Bases of the Cohomology Rings of Regular Nilpotent Hessenberg Varieties

M. Enokizono, T. Nagaoka, A. Tsuchiya, and the second author constructed in [27] an additive basis of the cohomology ring of a regular nilpotent Hessenberg variety $\operatorname{Hess}(N, h)$. This basis is obtained by extending the Poincaré duals $\left[\operatorname{Hess}\left(N, h^{\prime}\right)\right] \in$ $H^{*}(\operatorname{Hess}(N, h))$ of smaller regular nilpotent Hessenberg varieties $\operatorname{Hess}\left(N, h^{\prime}\right)$ with $h^{\prime} \subset h$. In particular, all of the classes $\left[\operatorname{Hess}\left(N, h^{\prime}\right)\right] \in H^{*}(\operatorname{Hess}(N, h))$ with $h^{\prime} \subset h$, are linearly independent. On the other hand, M. Harada, S. Murai, M. Precup, J. Tymoczko, and the second author derive in [38] a filtration on the cohomology ring $H^{*}(\operatorname{Hess}(N, h))$ of regular nilpotent Hessenberg varieties, from which they obtain a monomial basis for $H^{*}(\operatorname{Hess}(N, h))$. This basis is different from the one obtained in [27]. From the filtration they additionally obtain an inductive formula for the Poincaré polynomials of $\operatorname{Hess}(N, h)$; moreover, the monomial basis has an interpretation in terms of Schubert calculus.

### 5.6 The Volume Polynomials of Hessenberg Varieties

Recall from Corollary 3.3 that the cohomology ring $H^{*}(\operatorname{Hess}(N, h) ; \mathbb{Q})$ of a regular nilpotent Hessenberg variety is a Poincaré duality algebra. In [10], T. Abe, M. Masuda, S. Murai, T. Sato, and the second author computed the volume polynomial of this ring $H^{*}(\operatorname{Hess}(N, h) ; \mathbb{Q})$, and it precisely gives the volume of a certain embedding of any regular Hessenberg variety associated with $h$ into a projective space [2, 5]. M. Harada, M. Masuda, S. Park, and the second author provided a combinatorial formula for this polynomial in terms of the volumes of certain faces of the Gelfand-Zetlin polytope [37].

### 5.7 Hessenberg Varieties of Parabolic Type

J. Tymoczko and M. Precup showed that the Betti numbers of parabolic Hessenberg varieties decompose into a combination of those of Springer fibers and Schubert varieties associated to the parabolic [54]. As a corollary, they deduced that the Betti numbers of some parabolic Hessenberg varieties in Lie type A are equal to those of a specific union of Schubert varieties.

### 5.8 Twins for Regular Semisimple Hessenberg Varieties

Given a Hessenberg function $h:[n] \rightarrow[n]$, A. Ayzenberg and V. Buchstaber introduced a smooth $T$-manifold $X_{h}$, where $T$ is the maximal torus of $\operatorname{GL}(n, \mathbb{C})$ given in Sect.3.2. This manifold is similar to the regular semisimple Hessenberg variety $\operatorname{Hess}(S, h)$ in some sense. For example, they have the same Betti numbers and their $T$-equivariant cohomology rings are isomorphic as rings. See [13] for details.

### 5.9 Integrable Systems and Hessenberg Varieties

For a Hessenberg function $h:[n] \rightarrow[n]$, we denote by $\mathcal{X}(h)$ the family of Hessenberg varieties associated with $h$. Then, $\mathcal{X}(h)$ in fact have a structure of a vector bundle over the flag variety $F l\left(\mathbb{C}^{n}\right)$. For the case $h=(2,3,4, \ldots, n, n)$, the family $\mathcal{X}(h)$ contains the Peterson variety and the permutohedral variety. In this special case, P . Crooks and the first author showed that $\mathcal{X}(h)$ admits a Poisson structure with an open dense symplectic leaf, and that there is a completely integrable system on $\mathcal{X}(h)$ which contains the Toda lattice as a sub-system [4].

### 5.10 The Poset of Hessenberg Varieties

E. Drellich studied the poset of the Hessenberg varieties $\operatorname{Hess}(X, h)$ in $F l\left(\mathbb{C}^{n}\right)$ for a given $n \times n$ matrix $X$. She proved that if $X$ is not a scalar multiple of the identity matrix, then the Hessenberg functions determine distinct Hessenberg varieties. See [25] for details.

### 5.11 Springer Correspondence for Symmetric Spaces

As mentioned in Remark 2.2, Hessenberg varieties can be defined in a more general setting [33], and those Hessenberg varieties also appear in the works of T. H. Chen, K. Vilonen, and T. Xue [17-19] on Springer correspondence for symmetric spaces.

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# Stability of Bott-Samelson Classes in Algebraic Cobordism 

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#### Abstract

In this paper, we construct stable Bott-Samelson classes in the projective limit of the algebraic cobordism rings of full flag varieties, upon an initial choice of a reduced word in a given dimension. Each stable Bott-Samelson class is represented by a bounded formal power series modulo symmetric functions in positive degree. We make some explicit computations for those power series in the case of infinitesimal cohomology. We also obtain a formula of the restriction of Bott-Samelson classes to smaller flag varieties.


Keywords Schubert calculus • Cobordism • Flag variety • Bott-Samelson resolution

## 1 Introduction

Let k be an algebraically closed field of characteristic 0 . Let $F l_{n}$ be the flag variety of complete flags in $\mathrm{k}^{n}$. It can be identified with the homogeneous space $\mathrm{GL}_{n}(\mathrm{k}) / B$ where $B$ is the Borel subgroup of upper triangular matrices. For each permutation $w \in S_{n}$, the corresponding Schubert variety $X_{w}^{(n)} \subset F l_{n}$ is defined as $\overline{B_{-} w B}$, the closure of the orbits of $w B$ by the action of the opposite Borel subgroup $B_{-}$. If

[^30]$\iota_{n}: F l_{n} \rightarrow F l_{n+1}$ is the natural embedding, the cohomology fundamental classes of these Schubert varieties have the property that $\iota_{n}^{*}\left[X_{w}^{(n+1)}\right]=\left[X_{w}^{(n)}\right]$, i.e., the Schubert classes are stable under the pullback maps. The exact analogue of this property also holds in $K$-theory, in which one defines the Schubert classes as the $K$-theory classes of the structure sheaves of Schubert varieties.

In this paper, we attempt to generalize the above notion of stability to BottSamelson classes in algebraic cobordism. The algebraic cobordism, denoted by $\Omega^{*}$, was introduced by Levine-Morel in [16] and represents the universal object among oriented cohomology theories, a family of functors which includes both the Chow ring $\mathrm{CH}^{*}$ and $K^{0}\left[\beta, \beta^{-1}\right]$, a graded version of the Grothendieck ring of vector bundles. In recent years a lot of energy has been spent to lift results of Schubert calculus to $\Omega^{*}$, in the same way in which Bressler-Evens did in $[1,2]$ for topological cobordism. The first works in this direction were those of Calmés-Petrov-Zanoulline [3] and Hornbostel-Kiritchenko [6] who investigated the algebraic cobordism of flag manifolds. Later, the interest shifted to Grassmann and flag bundles (cf. [4, 7-13]). One of the main difficulties of Schubert calculus in algebraic cobordism is caused by the fact that the fundamental classes of Schubert varieties are not well-defined in general oriented cohomology theories. A candidate for the replacement of Schubert classes is the family of the push-forward classes of Bott-Samelson resolutions of Schubert varieties.

Since a Bott-Samelson variety is defined upon a choice of a reduced word, our stability of Bott-Samelson classes depends on a particular choice of a sequence of reduced words. The followings are the main results in this paper: (1) For a given BottSamelson variety $Y_{n}$ in $F l_{n}$, we construct a sequence of Bott-Samelson varieties $Y_{m}$ over $F l_{m}$ for $m \geq n$ such that their push-forward classes in algebraic cobordism are stable under pullbacks, namely, the identity $\iota_{m}^{*}\left[Y_{m+1} \rightarrow F l_{m+1}\right]=\left[Y_{m} \rightarrow F l_{m}\right]$ holds in $\Omega^{*}\left(F l_{m}\right)$ for all $m \geq n$; (2) For a given Bott-Samelson variety $Y_{n}$ over $F l_{n}$, we find an explicit formula for the pullback $\iota_{n-1}^{*}\left[Y_{n} \rightarrow F l_{n}\right]$ of its push-forward class in $\Omega^{*}\left(F l_{n-1}\right)$.

The pullback maps $\iota_{n}^{*}: \Omega^{*}\left(F l_{n+1}\right) \rightarrow \Omega^{*}\left(F l_{n}\right)$ give rise to a projective system of graded rings. Based on the ring presentation of $\Omega^{*}\left(F l_{n}\right)$ obtained by HornbostelKiritchenko [6], we observe that their graded projective limit, denoted by $\mathcal{R}$, is isomorphic to the graded ring of bounded formal power series in an infinite sequence of variables $x=\left(x_{i}\right)_{i \in \mathbb{Z}}{ }^{0}$ with coefficients in the Lazard ring $\mathbb{L}$ modulo the ideal of symmetric functions of positive degrees in $x$. Our stable sequence of Bott-Samelson classes determine a class in this limit, which we call a stable Bott-Samelson class. On each $\Omega^{*}\left(F l_{n}\right)$ the divided difference operators commute with the pullback maps and therefore lift to the limit $\mathcal{R}$. This gives a method of computing the power series representing stable Bott-Samelson classes, which we apply to the case of a chosen infinitesimal cohomology theory. In particular, we obtain a formula for the power series representing stable Bott-Samelson classes associated to dominant permutations.

In $[11,12]$, the first and second authors obtained determinant formula of the cobordism push-forward classes of so-called Damon-Kempf-Laksov resolutions, generalizing the classical Damon-Kempf-Laksov determinant formula of Schubert classes.

In [10], more explicit formula of Damon-Kempf-Laksov classes were obtained for infinitesimal cohomology. While these resolutions only exist for Schubert varieties associated to vexillary permutations (like for instance Grassmannian elements), their push-forward classes are stable and so is their determinantal formula. On the other hand, Naruse-Nakagawa [17-20] achieved, by considering a different resolution, a stable generalization of the Hall-Littlewood type formula for Schur polynomials in the context of topological cobordism. The differences among these stable expressions, including the ones obtained in this paper, should reflect the geometric nature of the different resolutions, each of which gives a different class in cobordism.

The paper is organized as follows. In Sect. 2, we recall basic facts about the algebraic cobordism ring of flag varieties and, in particular, we identify their projective limit. In Sect.3, we review the definition of Bott-Samelson resolutions and show the stability of their push-forward classes in cobordism based on the choice of a sequence of reduced words. We then focus on infinitesimal cohomology theory and compute, using divided difference operators, the power series representing the limits of the classes associated to dominant permutations. In Sect. 4, we prove a formula for the product of any Bott-Samelson class with the class $\left[F l_{n-1} \rightarrow F l_{n}\right]$, generalizing the restriction formula given in Sect. 3 .

## 2 Preliminary

Let k be an algebraic closed field of characteristic 0 .

### 2.1 Basics on Algebraic Cobordisms

For the reader's convenience, we will briefly recall some basic facts about algebraic cobordism and infinitesimal theories. More details on the construction and the properties of $\Omega^{*}$ can be found in [16], while a more comprehensive description of $I_{n}^{*}$ is given in [10].

Both $\Omega^{*}$ and $I_{n}^{*}$ are examples of oriented cohomology theories, a family of contravariant functors $A^{*}: \mathbf{S m}_{\mathrm{k}} \rightarrow \mathcal{R}^{*}$ from the category of smooth schemes to graded rings, which are furthermore endowed with push-forward maps for projective morphisms. Such functors are required to satisfy, together with some expected functorial compatibilities, the projective bundle formula and the extended homotopy property. These imply that, for every vector bundle $E \rightarrow X$, one is able to describe the evaluation of $A^{*}$ on the associated projective bundle $\mathbb{P}(E) \rightarrow X$ as well as on every $E$-torsor $V \rightarrow X$. The Chow ring $C H^{*}$ is probably the most well-known example of oriented cohomology theory and it should be kept in mind as a first approximation to the general concept.

As a direct consequence of the projective bundle formula one has that every oriented cohomology theory admits a theory of Chern classes, which can be defined
using Grothendieck's method. These satisfy most of the expected properties, like for instance the Whitney sum formula, however it is no longer true that the first Chern class behaves linearly with respect to tensor product: this is a key difference with $C H^{*}$. For a pair of line bundles $L$ and $M$ defined over the same base, classically one has

$$
\begin{equation*}
c_{1}^{C H}(L \otimes M)=c_{1}^{C H}(L)+c_{1}^{C H}(M) \text { and } c_{1}^{C H}\left(L^{\vee}\right)=-c_{1}^{C H}(L) \tag{1}
\end{equation*}
$$

but these equalities in general fail for $c_{1}^{A}$. Instead, in order to describe $c_{1}^{A}(L \otimes M)$, it becomes necessary to introduce a formal group law, a power series in two variables defined over the coefficient ring $F_{A} \in A^{*}(\operatorname{Spec} k)[[u, v]]$ satisfying some requirements. Similarly, expressing $c_{1}^{A}\left(L^{\vee}\right)$ in terms of $c_{1}^{A}(L)$ requires one to consider the formal inverse $\chi_{A} \in A^{*}(\operatorname{Spec} k)[[u]]$. The analogues of (1) then become

$$
\begin{equation*}
c_{1}^{A}(L \otimes M)=F_{A}\left(c_{1}^{A}(L), c_{1}^{A}(M)\right) \text { and } c_{1}^{A}\left(L^{\vee}\right)=\chi_{A}\left(c_{1}^{A}(L)\right) \tag{2}
\end{equation*}
$$

It is a classical result of Lazard [15] that every formal group law ( $R, F_{R}$ ) can be obtained from the universal one ( $\mathbb{L}, F_{\mathbb{L}}$ ), which is defined over a ring later named after him. He also proved that, as a graded ring, $\mathbb{L}=\bigoplus_{m \leq 0} \mathbb{L}^{m}$ is isomorphic to a polynomial ring in countably many variables $y_{i}$, each appearing in degree $-i$ for $i \geq 1$. In the case of a field of characteristic 0 , Levine and Morel were able to prove that the coefficient ring of algebraic cobordism is isomorphic to $\mathbb{L}$ and that its formal group law $F_{\Omega}$ coincides with the universal one, which from now on we will simply denote $F$. The universality of $\Omega^{*}$ does not restrict itself only to its coefficient ring, in fact, Levine and Morel were able to prove the following theorem.

Theorem 2.1 ([16, Theorem 1.2.6]) $\Omega^{*}$ is universal among oriented cohomology theories on $\mathbf{S m}_{\mathrm{k}}$. That is, for any other oriented cohomology theory $A^{*}$ there exists a unique morphism

$$
\vartheta_{A}: \Omega^{*} \rightarrow A^{*}
$$

of oriented cohomology theories.
It essentially follows formally from this result that for any given formal group law ( $R, F_{R}$ ) the functor $\Omega^{*} \otimes_{\mathbb{L}} R$ is universal among the oriented cohomology theories with $R$ as coefficient ring and $F_{R}$ as associated law. This procedure can be used to produce functors, like the infinitesimal theories $I_{n}^{*}$, whose formal group laws are far simpler than the universal one and as a consequence more suitable for explicit computations. More precisely the projection $\mathbb{L} \rightarrow \mathbb{Z}\left[y_{n}\right] /\left(y_{n}^{2}\right)$, which maps $y_{i}$ to 0 unless $i=n$, gives rise to the following formal group law $F_{I_{n}}$ on $\mathbb{Z}\left[y_{n}\right] /\left(y_{n}^{2}\right)$ :

$$
\begin{equation*}
F_{I_{n}}(u, v)=u+v+y_{n} \cdot \frac{1}{d_{n}} \sum_{j=1}^{n}\binom{n+1}{j} u^{j} v^{n+1-j} \tag{3}
\end{equation*}
$$

Here one has $d_{n}=p$, if $n+1$ is a power of a prime $p$, and $d_{n}=1$ otherwise. In our computations we will only consider the case $n=2$, for which (3) becomes

$$
u \boxplus v:=F_{I_{2}}(u, v)=u+v+y_{2}\left(u^{2} v+u v^{2}\right)=(u+v)\left(1+y_{2} u v\right)
$$

with the formal inverse being $\boxminus u:=\chi_{I_{2}}(u)=-u$. For the remainder of the paper we will write $\gamma$ instead of $y_{2}$.

Let us finish this overview by discussing fundamental classes, another aspect in which a general oriented cohomology theory differs from $\mathrm{CH}^{*}$. While in $\mathrm{CH}^{*}$ it is possible to associate such a class to every equi-dimensional scheme, for a general oriented cohomology theory $A^{*}$ one has to restrict to schemes whose structure morphism is a local complete intersection. In particular, since not all Schubert varieties satisfy this requirement, it becomes necessary to find an alternative definition for Schubert classes. One possible option is to choose a family of resolutions of singularities and replace the fundamental classes of Schubert varieties with the pushforwards of the associated resolutions.

### 2.2 Algebraic Cobordism of Flag Varieties and Their Limit

For any integers $a, b$ such that $a \leq b$, let $[a, b]:=\{a, a+1, \ldots, b\}$. Let $\mathrm{k}^{\mathbb{Z}>0}$ be the infinite dimensional vector space generated by a formal basis $\left(e_{i}\right)_{i \in \mathbb{Z}_{>0}}$. For each $m \in \mathbb{Z}_{>0}$, let $E_{m}$ be the subspace of $\mathrm{k}^{\mathbb{Z}_{>0}}$ generated by $e_{1}, \ldots, e_{m}$. We set $E_{0}=0$. We often identify $E_{m}$ with $\mathrm{k}^{m}$ the space of column vectors.

For each $n \in \mathbb{Z}_{>0}$, the flag variety $F l_{n}$ consists of flags $U_{\bullet}=\left(U_{i}\right)_{i \in[1, n-1]}$ of subspaces in $E_{n}$ where $U_{i} \subset U_{i+1}$ and $\operatorname{dim} U_{i}=i$ for each $i \in[1, n-1]$. Note that this implies $U_{n}=E_{n}$. For a fixed $n$, let $\mathcal{U}_{i}^{(n)}, i \in[0, n]$ denote the tautological vector bundles of $F l_{n}$ and $\mathcal{E}_{i}$ the trivial bundles of fiber $E_{i}$. In particular, $\mathcal{U}_{0}^{(n)}=0$ and $\mathcal{U}_{n}^{(n)}=\mathcal{E}_{n}$.

Let $\mathrm{GL}_{n}(\mathrm{k})=\mathrm{GL}\left(E_{n}\right)$ be the general linear group. We consider the maximal torus $T_{n} \subset \mathrm{GL}_{n}(\mathrm{k})$ given by the matrices having $\left(e_{i}\right)_{i \in[1, n]}$ as a basis of eigenvectors and the Borel subgroup $B_{n} \subset \mathrm{GL}_{n}(\mathrm{k})$ given by the upper triangular matrices stabilizing the flag $E_{\bullet}=\left(E_{i}\right)_{i \in[1, n-1]}$ in $F l_{n}$. We can identify $F l_{n}$ with the homogeneous space $\mathrm{GL}_{n}(\mathrm{k}) / B_{n}$ by associating the matrix $M=\left(u_{1}, \ldots, u_{n}\right)$ to a flag $U$. where $\left\{u_{j}\right\}_{j \in[1, i]}$ is a basis of $U_{i}$.

There is an isomorphism of graded rings ([6, Theorem 1.1])

$$
\begin{equation*}
\Omega^{*}\left(F l_{n}\right) \cong \mathbb{L}\left[x_{1}, \ldots, x_{n}\right] / \mathbb{S}_{n} \tag{4}
\end{equation*}
$$

sending $c_{1}\left(\left(\mathcal{U}_{i}^{(n)} / \mathcal{U}_{i-1}^{(n)}\right)^{\vee}\right)$ to $x_{i}$, where $\mathbb{S}_{n}$ is the ideal generated by the homogeneous symmetric polynomials in $x_{1}, \ldots, x_{n}$ of strictly positive degree.

Let $\iota_{n}: F l_{n} \hookrightarrow F l_{n+1}$ be the embedding induced by the canonical inclusion $E_{n} \hookrightarrow E_{n+1}$. We have $\iota_{n}^{*} \mathcal{U}_{i}^{(n+1)}=\mathcal{U}_{i}^{(n)}$ for all $i \in[1, n]$ and $\iota_{n}^{*} \mathcal{U}_{n+1}^{(n+1)}=\mathcal{E}_{n+1}$. As
a consequence, under the isomorphism (4), the pullback map $\iota_{n}^{*}: \Omega^{*}\left(F l_{n+1}\right) \rightarrow$ $\Omega^{*}\left(F l_{n}\right)$ is the natural projection given by setting $x_{n+1}=0$. For each $m \in \mathbb{Z}$, let $\mathcal{R}^{m}$ be the projective limit of $\Omega^{m}\left(F l_{n}\right)$ with respect to $\iota_{n}^{*}$. We define the graded projective limit of $\Omega^{*}\left(F l_{n}\right)$ with respect to $\iota_{n}^{*}$ to be $\mathcal{R}:=\bigoplus_{m \in \mathbb{Z}} \mathcal{R}^{m}$.

In order to give a ring presentation of $\mathcal{R}$, we introduce the following ring of formal power series. Let $x=\left(x_{i}\right)_{i \in \mathbb{Z}_{>0}}$ be a sequence of infinitely many indeterminates. Let $\mathbb{Z}^{\infty}$ be the set of infinite sequence $\mathbf{s}=\left(s_{i}\right)_{i \in \mathbb{Z}_{>0}}$ of nonnegative integers such that all but finitely many $s_{i}$ 's are zero. Let $\mathbb{L}[[x]]^{(m)}$ be the space of formal power series of degree $m \in \mathbb{Z}$. An element $f(x)$ of $\mathbb{L}[[x]]^{(m)}$ is uniquely given as

$$
f(x)=\sum_{\mathbf{s} \in \mathbb{Z}^{\infty}} a_{\mathbf{s}} x^{\mathbf{s}}, \quad a_{\mathbf{s}} \in \mathbb{L}, \quad x^{\mathbf{s}}=\prod_{i=1}^{\infty} x_{i}^{s_{i}}
$$

such that $|\mathbf{s}|+\operatorname{deg} a_{\mathbf{s}}=m$ where $|\mathbf{s}|=\sum_{i=0}^{\infty} s_{i}$ and $\operatorname{deg} a_{\mathbf{s}}$ is the degree of $a_{\mathbf{s}}$ in $\mathbb{L}$. An element $f(x) \in \mathbb{L}[[x]]^{(m)}$ is bounded if $p_{n}(f(x)) \in \mathbb{L}\left[x_{1}, \ldots, x_{n}\right]^{(m)}$, where $p_{n}$ is the substitution of $x_{k}=0$ for all $k>n$ and $\mathbb{L}\left[x_{1}, \ldots, x_{n}\right]^{(m)}$ is the degree $m$ part of $\mathbb{L}\left[x_{1}, \ldots, x_{n}\right]$. Let $\mathbb{L}[[x]]_{b d}^{(m)}$ be the set of all such bounded elements of $\mathbb{L}[[x]]^{m}$. We set

$$
\mathbb{L}[[x]]_{b d}:=\bigoplus_{m \in \mathbb{Z}} \mathbb{L}[[x]]_{b d}^{(m)}
$$

This is a graded sub $\mathbb{L}$-algebra of the ring $\mathbb{L}[[x]]$ of formal power series.
Proposition 2.2 There is an isomorphism of graded $\mathbb{L}$-algebras

$$
\mathcal{R} \cong \mathbb{L}[[x]]_{b d} / \mathbb{S}_{\infty}
$$

where $\mathbb{S}_{\infty}$ is the ideal of $\mathbb{L}[[x]]_{b d}$ generated by symmetric functions in $x$ of strictly positive degree.
Proof Let $m \in \mathbb{Z}$. The projections $p_{n}: \mathbb{L}[[x]]_{b d}^{(m)} \rightarrow \mathbb{L}\left[x_{1}, \ldots, x_{n}\right]^{(m)}$ for $n>0$ induce a surjective homomorphism

$$
\Phi: \mathbb{L}[[x]]_{b d}^{(m)} \rightarrow \lim _{n \rightarrow \infty} \mathbb{L}\left[x_{1}, \ldots, x_{n}\right]^{(m)}
$$

sending $f(x)$ to $\left\{p_{n}(f(x))\right\}_{n \in \mathbb{Z}>0}$. It is also easy to see that $\Phi$ is injective, and thus an isomorphism. Moreover, $p_{n}$ 's induce surjections

$$
\mathbb{L}[[x]]_{b d}^{(m)} \cap \mathbb{S}_{\infty} \rightarrow \mathbb{L}\left[x_{1}, \ldots, x_{n}\right]^{(m)} \cap \mathbb{S}_{n}, \quad n>0
$$

inducing a bijection

$$
\mathbb{L}[[x]]_{b d}^{(m)} \cap \mathbb{S}_{\infty} \cong \lim _{n \rightarrow \infty}\left(\mathbb{L}\left[x_{1}, \ldots, x_{n}\right]^{(m)} \cap \mathbb{S}_{n}\right)
$$

Thus we obtain the isomorphism

$$
\bigoplus_{m \in \mathbb{Z}} \mathbb{L}[[x]]_{b d}^{(m)} /\left(\mathbb{L}[[x]]_{b d}^{(m)} \cap \mathbb{S}_{\infty}\right) \cong \bigoplus_{m \in \mathbb{Z}} \lim _{n \rightarrow \infty} \mathbb{L}\left[x_{1}, \ldots, x_{n}\right]^{(m)} /\left(\mathbb{L}\left[x_{1}, \ldots, x_{n}\right]^{(m)} \cap \mathbb{S}_{n}\right)
$$

which is the desired one.
Definition 2.3 An element in $\mathcal{R}^{i}$ is a sequence $\left(\alpha_{n}\right)_{n \in \mathbb{Z}_{>0}}$ such that $\alpha_{n} \in \Omega^{i}\left(F l_{n}\right)$ and $\iota_{n}^{*}\left(\alpha_{n+1}\right)=\alpha_{n}$ for all $n>0$. An element of $\mathcal{R}$ is a finite linear combinations of such sequences and we call it a stable class.

Remark 2.4 In order to specify an element of $\mathcal{R}^{i}$, we only need to provide $\alpha_{i}$ for all $i \geq N$ for some fixed integer $N$. In fact, for $i<N$ the elements $\alpha_{i}$ can be obtained from $\alpha_{N}$ by applying the projections $\iota_{n}^{*}$.

### 2.3 Divided Difference Operators

Let $W_{n}$ be the Weyl group of $\mathrm{GL}_{n}(\mathrm{k})$. The maximal torus $T_{n}$ and the Borel subgroup $B_{n}$ define a system of simple reflections $s_{1}, \ldots, s_{n-1} \in W_{n}$ and we can identify $W_{n}$ with the symmetric group $S_{n}$ in $n$ letters, where each $s_{i}$ corresponds to the transposition of the letters $i$ and $i+1$. We denote the length of $w$ by $\ell(w)$.

For each $i \in[1, n-1]$, the divided difference operator $\partial_{i}$ is an operator on $\Omega^{*}\left(F l_{n}\right)$ defined as follows. Let $F l_{n}^{(i)}$ be the partial flag variety consisting of flags of the form $U_{1} \subset \cdots \subset U_{i-1} \subset U_{i+1} \subset \cdots \subset U_{n-1}$ with $\operatorname{dim} U_{k}=k$. Denote the canonical projection $F l_{n} \rightarrow F l_{n}^{(i)}$ by $p_{i}$. Then define $\partial_{i}:=p_{i *} \circ p_{i}^{*}$. It is known from [6] that under the presentation (4), we have

$$
\begin{equation*}
\partial_{i}(f(x))=\left(\mathrm{id}+s_{i}\right) \frac{f(x)}{F\left(x_{i}, \chi\left(x_{i+1}\right)\right)}=\frac{f(x)}{F\left(x_{i}, \chi\left(x_{i+1}\right)\right)}+\frac{s_{i} f(x)}{F\left(x_{i+1}, \chi\left(x_{i}\right)\right)} . \tag{5}
\end{equation*}
$$

Lemma 2.5 The pullback $\iota_{n}^{*}: \Omega^{*}\left(F l_{n+1}\right) \rightarrow \Omega^{*}\left(F l_{n}\right)$ commutes with $\partial_{i}$ for all $i \in$ [1, $n-1]$. In particular, this shows that $\partial_{i}$ can be defined in the projective limit $\mathcal{R}$ and it is given by the formula (5).

Proof For each $i \in[1, n-1], \iota_{n}$ and $p_{i}$ form a fiber diagram

and, since they are transverse, we have $\iota_{n}^{*} \circ p_{i *}=p_{i *} \circ \iota_{n}^{*}$. Thus $\iota_{n}^{*} \circ \partial_{i}=\iota_{n}^{*} \circ p_{i *} \circ$ $p_{i}^{*}=p_{i *} \circ \iota_{n}^{*} \circ p_{i}^{*}=p_{i *} \circ p_{i}^{*} \circ \iota_{n}^{*}=\partial_{i} \circ \iota_{n}^{*}$.

For a permutation $w \in W_{n}$, let $\dot{X}_{w}^{(n)}=B_{n} \cdot w\left(E_{\mathbf{\bullet}}\right)$ be the Bruhat cell associated to $w$ in $F l_{n}$, where $w\left(E_{\bullet}\right)$ is the flag consisting of $w\left(E_{i}\right)=\left\langle e_{w(1)}, \ldots, e_{w(i)}\right\rangle$ for
each $i \in[1, n-1]$. The Schubert varieties $X_{w}^{(n)}$ are the closures of the Bruhat cells: $X_{w}^{(n)}:=\overline{B_{n} \cdot w\left(E_{\bullet}\right)}$. The opposite Schubert varieties are defined via $X_{(n)}^{w}=w_{0}$. $X_{w_{0} w}^{(n)}$, where $w_{0}=w_{0}^{(n)}$ is the longest element of $W_{n}$. As an orbit closure, we have $X_{(n)}^{w}=\overline{B_{n}^{-} \cdot w\left(E_{\bullet}\right)}$ where $B_{n}^{-}:=w_{0} B_{n} w_{0}$ is the opposite Borel subgroup of lower triangular matrices.

Remark 2.6 The fundamental class $\left[X_{(n)}^{w}\right]$ of $X_{(n)}^{w}$ is well-defined in the Chow ring of $F l_{n}$. Those classes are stable along pullbacks, i.e., $\iota_{n}^{*}\left[X_{(n+1)}^{w}\right]=\left[X_{(n)}^{w}\right]$ in $C H^{*}\left(F l_{n}\right)$ where $w \in S_{n}$ is regarded as an element of $S_{n+1}$ under the natural embed$\operatorname{ding} S_{n} \subset S_{n+1}$. As it is well-known, its stable limit can be identified with the Schubert polynomial of Lascoux-Schützenberger [14]. It is also worth mentioning that the Schubert classes admit the following compatibility with divided difference operators, reflected on the definition of Schubert polynomials: for each $i \in[1, n-1]$, we have

$$
\partial_{i}\left[X_{(n)}^{w}\right]= \begin{cases}{\left[X_{(n)}^{w s_{i}}\right]} & \ell\left(w s_{i}\right)=\ell(w)+1 \\ 0 & \text { otherwise }\end{cases}
$$

### 2.4 Some Facts on Permutations and Reduced Words

We conclude this section by fixing notations for reduced words and showing a few lemmas and a proposition that will be used in the rest of the paper.

We denote by $\underline{W}_{n}$ the set of words in $s_{1}, \ldots, s_{n-1}$ : an element of $\underline{W}_{n}$ will be written as a finite sequence $s_{i_{1}} \ldots s_{i_{r}}$, while the empty word is denoted by 1 . The length of a word $\underline{w}=s_{i_{1}} \ldots s_{i_{r}}$ is the number $r$ of the letters $s_{i}$ 's in $\underline{w}$ and we denote it by $\ell(\underline{w})$. For a word $\underline{w} \in \underline{W}_{n}$, we denote the corresponding permutation by $w \in W_{n}$. Let $W_{n}^{i}$ be the subgroup of $W_{n}$ generated by all simple reflections $s_{j}$ with $j \neq i$ and $\underline{W}_{n}^{i}$ the corresponding set of words. In particular, we can identify $W_{n}$ with $W_{n+1}^{n}$ and $\underline{W}_{n}$ with $\underline{W}_{n+1}^{n}$.

We denote the Bruhat order in $W_{n}$ by $\leq$, i.e., $w \leq v$ if and only if every reduced word for $v$ contains a subword which is a reduced word for $w$.

We denote by $c^{(n)}$ the Coxeter element $s_{1} \ldots s_{n}$ of $W_{n+1}$. It has a unique reduced word $\underline{c}^{(n)}=s_{1} \ldots s_{n}$. Note that $\underline{c}^{(n)} \underline{c}^{(n-1)} \ldots \underline{c}^{(1)}$ is a reduced word for the longest element $w_{0}^{(n+1)}$ of $W_{n+1}$.

Lemma 2.7 If $\underline{c}^{(n)} \underline{v} \in \underline{W}_{n+1}$ is a reduced word, then $\underline{v}$ is a reduced word in $\underline{W}_{n}$.
Proof There exists a reduced word $\underline{u}$ such that $\underline{c}^{(n)} \underline{v} \underline{u}=\underline{w}_{0}^{(n+1)}$ is a reduced word for the longest element $w_{0}^{(n+1)} \in W_{n+1}$. Since $v u=\left(c^{(n)}\right)^{-1} w_{0}^{(n+1)}=w_{0}^{(n)}$, we have $v u \in W_{n}$ so that any reduced word of $v u$ lies in $\underline{W}_{n}$ and in particular $\underline{v} \underline{u}$ is a reduced word in $\underline{W}_{n}$. Thus $\underline{v}$ is a reduced word in $\underline{W}_{n}$.
Lemma 2.8 If $\underline{v} \in \underline{W}_{n}$ is a reduced word, then $\underline{c}^{(n)} \underline{v} \in \underline{W}_{n+1}$ is a reduced word. In particular, if $v=w_{0}^{(n)} w$ for some $w \in W_{n}$, then $\underline{\underline{c}}^{(n)} \underline{v}$ is a reduced word for $w_{0}^{(n+1)} w$.

Proof There exists a reduced word $\underline{u}$ such that $\underline{v} \underline{u}$ is a reduced word for $w_{0}^{(n)}$. Then $\underline{c}^{(n)} \underline{v} \underline{u}$ is a reduced word for $w_{0}^{(n+1)}$. This implies that $\underline{c}^{(n)} \underline{v}$ is a reduced word.

Proposition 2.9 Let $w \in W_{n+1}$ such that $c:=c^{(n)} \leq w$. Every reduced word $\underline{w} \in$ $\underline{W}_{n+1}$ for $w$ decomposes, modulo commuting relations, as $\underline{w}=\underline{u} \underline{c} \underline{v}$ with $\underline{u} \in \underline{W}_{n+1}^{1}$ and $\underline{v} \in \underline{W}_{n}$.

Proof In this proof, all the equalities of words are modulo commuting relations. Since $c \leq w$, the definition of the Bruhat order implies that $\underline{w}$ contains as a subword $\underline{c}$, the unique reduced word of $c$. We choose such a subword by selecting the first occurrence of $s_{1}$, the first occurrence of $s_{2}$ after the chosen $s_{1}$ and so on. We thus have a decomposition

$$
\underline{w}=\underline{w}_{1} s_{1} \underline{w}_{2} s_{2} \underline{w}_{3} \ldots s_{n-1} \underline{w}_{n} s_{n} \underline{w}_{n+1}
$$

with $\underline{w}_{i} \in \underline{W}_{n+1}^{i}$ for $i \in[1, n+1]$. We have $\underline{w}_{i}=\underline{v}_{i} \underline{u}_{i}$ for $i \in[2, n]$, where $\underline{v}_{i}$ is a word in the $s_{k}$ 's for $k \in[1, i-1]$ and $\underline{u}_{i}$ is a word in the $s_{k}$ 's for $k \in[i+1, n]$. Note in particular that we have $\underline{u}_{n}=1$. Observing that $\underline{v}_{i} \underline{u}_{j}=\underline{u}_{j} \underline{v}_{i}$ and $s_{i-1} \underline{u}_{j}=\underline{u}_{j} s_{i-1}$ for $2 \leq i \leq j \leq n$, we thus obtain

$$
\underline{w}^{w}=\underline{w}_{1}\left(\underline{u}_{2} \underline{u}_{3} \ldots \underline{u}_{n}\right)\left(s_{1} \underline{v}_{2} s_{2} \underline{v}_{3} \ldots \underline{v}_{n} s_{n}\right) \underline{w}_{n+1} .
$$

For each $i \in[2, n]$, we claim that the word $\underline{v}_{i}$ does not contain $s_{i-1}$, i.e., it is a word in the $s_{k}$ 's for $1 \leq k \leq i-2$. We prove the claim by induction on $i$. First of all, it is easy to see that $\underline{v}_{2}$ is an empty word since it is a word in $s_{1}$ only, and there is $s_{1}$ on the left of $\underline{v}_{2}$ in the reduced word $\underline{w}$. Assume that the claim holds for $i \leq k$. We have

$$
\underline{w}^{\underline{w}}=\underline{w}_{1}\left(\underline{u}_{2} \underline{u}_{3} \ldots \underline{u}_{n}\right)\left(s_{1} s_{2} \ldots s_{k}\right)\left(\underline{v}_{2} \underline{v}_{3} \ldots \underline{v}_{k}\right)\left(\underline{v}_{k+1} s_{k+1} \ldots \underline{v}_{n} s_{n}\right) \underline{w}_{n+1} .
$$

Since $s_{1} s_{2} \ldots s_{k} \underline{v}_{1} \ldots \underline{v}_{k} \underline{v}_{k+1} \in \underline{W}_{k+1}$ is reduced, Lemma 2.7 implies that $\underline{v}_{1} \ldots$ $\underline{v}_{k+1} \in \underline{W}_{k}$ and, in particular, we find that $\underline{v}_{k+1}$ doesn't contain $s_{k}$. Thus the claim holds. Now by moving all $\underline{v}_{i}$ to the right using commuting relations, we obtain

$$
\underline{w}=\underline{w}_{1} \underline{u}_{2} \underline{u}_{3} \ldots \underline{u}_{n} s_{1} s_{2} \ldots s_{n} \underline{v}_{2} \underline{v}_{3} \ldots \underline{v}_{n} \underline{w}_{n+1} .
$$

Using Lemma 2.7 again, we obtain $\underline{v}_{2} \underline{v}_{3} \ldots \underline{v}_{n} \underline{w}_{n+1} \in \underline{W}_{n}$. This proves the proposition since $\underline{w}_{1} \underline{u}_{2} \underline{u}_{3} \ldots \underline{u}_{n} \in \underline{W}_{n+1}^{1}$.

## 3 Stable Bott-Samelson Classes

In this section, we introduce stable Bott-Samelson classes in the limit $\mathcal{R}$ of $\Omega^{*}\left(F l_{n}\right)$. We also compute some of those classes explicitly in the case of infinitesimal cohomology.

### 3.1 The Stability of Bott-Samelson Classes

A Schubert variety is, in general, normal and Cohen-Macaulay, and has rational singularities. There exists several resolutions of singularities for it. We will be interested in the so-called Bott-Samelson resolutions.

We set $F_{i}^{(n)}:=\left\langle e_{n}, \ldots, e_{n+1-i}\right\rangle$ and denote the trivial bundle with fiber $F_{i}^{(n)}$ by $\mathcal{F}_{i}^{(n)}$.

Definition 3.1 For a reduced word $\underline{v}=s_{i_{1}} \ldots s_{i_{r}} \in \underline{W}_{n}$, the Bott-Samelson variety $Y_{\underline{v}}^{(n)}$ is a subvariety of $\left(F l_{n}\right)^{r}$ defined as follows:

$$
Y_{\underline{v}}^{(n)}=\left\{\left(U_{\bullet}^{[0]}, U_{\bullet}^{[1]}, \ldots, U_{\bullet}^{[r]}\right) \in\left(F l_{n}\right)^{r} \mid U_{i}^{[k-1]}=U_{i}^{[k]}, \forall k=[1, r], \forall i \in[1, n-1] \backslash\left\{i_{k}\right\}\right\},
$$

where $U_{\bullet}^{[0]}=F_{\bullet}^{(n)}$. If there is no confusion, we will sometimes write $Y_{\underline{v}}$ for $Y_{\underline{v}}^{(n)}$.
Remark 3.2 In Definition 4.1 we will give another equivalent construction (denoted $X_{\underline{w}}$ ) of the Bott-Samelson resolutions.

It is well-known (cf.[5]) that $Y_{\underline{v}}$ is a smooth projective variety of dimension $r$. Let $\pi_{n}:\left(F l_{n}\right)^{r} \rightarrow F l_{n}$ be the projection to the $r$-th component. If $w \in W_{n}$ and $v=w_{0}^{(n)} w$, the projection $\pi_{n}$ induces a birational map $Y_{\underline{v}} \rightarrow X^{w}$, which we refer to as a Bott-Samelson resolution of $X^{w} \subset F l_{n}$.

Theorem 3.3 Let $\underline{v} \in \underline{W}_{n}$ be a reduced word. There is a fiber diagram

and we have $\iota_{n}^{*}\left(\left[Y_{\underline{c}^{(n)} \underline{v}} \rightarrow F l_{n+1}\right]\right)=\left[Y_{\underline{v}} \rightarrow F l_{n}\right]$.
Furthermore, let $\underline{c}^{[n+m]}:=\underline{c}^{(n+m-1)} \ldots \underline{c}^{(n+1)} \underline{c}^{(n)}$ where $\underline{c}^{[n]}=1$, then the sequence

$$
\left[Y_{\underline{c}^{[n+m]} \underline{v}} \rightarrow F l_{n+m}\right], \quad m \geq 0
$$

defines a stable class in $\mathcal{R}$, which we call a stable Bott-Samelson class associated to $w$ and denote by $\mathrm{BS} \frac{v}{w}$ if $v=w_{0}^{(n)} w$.

Proof First we note that, by definition, an element of $Y_{\underline{v}}^{(n)}$ can be specified by a sequence of subspaces $\left(V_{1}, \ldots, V_{r}\right)$ where $V_{k}=U_{i_{k}}^{[k]}$. We show that the map $\tilde{\iota}_{n}$ : $Y_{\underline{v}}^{(n)} \rightarrow Y_{\underline{c}^{(n)} \underline{v}}^{(n+1)}$ defined by

$$
\tilde{\iota}_{n}\left(V_{1}, V_{2}, \ldots, V_{r}\right):=\left(F_{1}^{(n)}, \ldots, F_{n}^{(n)}, V_{1}, \ldots, V_{r}\right)
$$

gives the desired fiber diagram. If we write an element of $Y_{\underline{c}^{(n)} \underline{\underline{v}}}^{(n+1)}$ as

$$
\left(A_{\bullet}^{[1]}, \ldots, A_{\bullet}^{[n]}, B_{\bullet}^{[1]}, \ldots, B_{\bullet}^{[r]}\right)
$$

it suffices to show that $A_{k}^{[k]}=F_{k}^{(n)}$ for all $k \in[1, n]$ over the image of $F l_{n}$. Suppose that $B_{\bullet}^{[r]}$ is in the image of $F l_{n}$, then $B_{n}^{[r]}=E_{n}$. Since $i_{1}, \ldots, i_{r} \in[1, n-1]$, we have $A_{n}^{[n]}=E_{n}=F_{n}^{(n)}$. We use backward induction on $k$ with the base case being $k=n$. Assume $A_{k+1}^{[k+1]}=F_{k+1}^{(n)}$. We then have

$$
A_{k}^{[k]} \subset F_{k+1}^{(n+1)} \cap A_{k+1}^{[k+1]}=F_{k+1}^{(n+1)} \cap F_{k+1}^{(n)}=F_{k}^{(n)}
$$

For the latter claim, we use the identity $\iota_{n}^{*} \pi_{n+1 *}=\pi_{n *} \overbrace{n}^{*}$ (see [16, p. 144 (BM2)]). We get
$\iota_{n}^{*}\left[Y_{\underline{\mathcal{C}^{(n)} \underline{v}}} \rightarrow F l_{n+1}\right]=\iota_{n}^{*} \pi_{n+1 *}\left(1_{Y_{\underline{c^{(n) \underline{v}}}}}\right)=\pi_{n *} \tilde{i}_{n}^{*}\left(1_{Y_{\underline{\underline{C}}(n) \underline{\underline{v}}}}\right)=\pi_{n *}\left(1_{Y_{\underline{v}}}\right)=\left[Y_{\underline{v}} \rightarrow F l_{n}\right]$.
This completes the proof of the claim.
Remark 3.4 We sometimes denote $\mathrm{BS} \frac{v}{w}$ by $\mathrm{BS}_{\bar{w}}^{v}(x)$ in order to stress that we regard it as an element of $\mathbb{L}[x]_{b d} / \mathbb{S}_{\infty}$ under the identification in Proposition 2.2.

The following compatibility of Bott-Samelson classes with divided difference operators was established in [6].

Lemma 3.5 For a reduced word $\underline{v}=s_{i_{1}} \ldots s_{i_{r}} \in \underline{W}_{n}$, and $k \in[1, n-1]$, we have

$$
\partial_{i}\left[Y_{\underline{v}} \rightarrow F l_{n}\right]= \begin{cases}{\left[Y_{\underline{v} s_{i}} \rightarrow F l_{n}\right]} & \text { if } \underline{v} s_{i} \text { is a reduced word } \\ 0 & \text { otherwise }\end{cases}
$$

Since, as explained in Sect. 2.2, the divided difference operators commute with the pullbacks $\iota_{n}^{*}$, we obtain the next corollary.

Corollary 3.6 Let $w \in W_{n}$ and set $v=w_{0}^{(n)} w$. Let $\underline{v}$ be a reduced word of $v$. For any $i \in \mathbb{Z}_{>0}$, we have

$$
\partial_{i} \mathrm{BS}{ }_{w}^{v}= \begin{cases}\mathrm{BS}_{w s_{i}}^{\frac{v s_{i}}{v}} & \text { if } \ell\left(w s_{i}\right)<\ell(w) \\ 0 & \text { otherwise }\end{cases}
$$

Remark 3.7 In the connective $K$-theory of $F l_{n}$, the class $\left[Y_{\underline{v}} \rightarrow F l_{n}\right]$ coincides with the class of the opposite Schubert variety $X^{w}$, provided that $\bar{v}=w_{0}^{(n)} w$. Its associated class in the projective limit is represented by the Grothendieck polynomial $\mathfrak{G}_{w}(x)$ associated to $w$.

### 3.2 Examples in Infinitesimal Cohomology

Throughout this section we will consider infinitesimal cohomology instead of algebraic cobordism. In combination with Proposition 2.2, the use of this simpler theory will allow us to perform an explicit computation of the stable Bott-Samelson classes in terms of power series in $x$.

As in Sect. 2.1, the formal group law and its formal inverse for the infinitesimal cohomology $I_{2}^{*}$ are given by

$$
F_{I_{2}}(x, y)=x \boxplus y=(x+y)(1+\gamma x y), \quad \chi_{I_{2}}(x)=\boxminus x=-x
$$

with $\gamma^{2}=0$. We denote $I_{2}^{*}(\mathrm{pt})=\mathbb{Z}[\gamma] /\left(\gamma^{2}\right)$ by $\mathbb{I}$. As explained in Sect. 2.1, we have

$$
I_{2}^{*}\left(F l_{n}\right) \cong \mathbb{I}\left[x_{1}, \ldots, x_{n}\right] / \mathbb{S}_{n}
$$

where $\mathbb{S}_{n}$ is the ideal generated by the homogeneous symmetric polynomials of strictly positive degree in $x_{1}, \ldots, x_{n}$. We set

$$
\mathcal{R}_{\mathbb{I}}:=\mathcal{R} \otimes_{\mathbb{L}} \mathbb{I}=\mathbb{I}[x]_{b d} / \mathbb{S}_{\infty}
$$

By specialising (5) to this particular case we obtain that on $\mathcal{R}_{\mathbb{I}}$ the divided difference operator $\partial_{i}$ is given by

$$
\partial_{i} f=\frac{f-s_{i} f}{x_{i}-x_{i+1}} \cdot\left(1+\gamma x_{i} x_{i+1}\right), \quad f \in \mathcal{R}_{\mathbb{I}}
$$

Remark 3.8 (1) If $f$ is symmetric in $x_{i}$ and $x_{i+1}$, then $\partial_{i}(f g)=f \partial_{i} g$ for all $g \in$ $\mathcal{R}_{\mathbb{I}}$.
(2) If $|i-j| \geq 2$, then $\partial_{i} \partial_{j}=\partial_{j} \partial_{i}$.

For a reduced word $\underline{v}=s_{i_{1}} \ldots s_{i_{r}}$, let $\partial_{\underline{v}}:=\partial_{i_{r}} \ldots \partial_{i_{1}}$. Recall that $\underline{c}^{(n)}=s_{1} \ldots s_{n}$.
Lemma 3.9 For $n \geq 1$, we have

$$
\partial_{\underline{c}^{(n)}}\left(x_{1}^{n} x_{2}^{n-1} \ldots x_{n}\right)=\left(x_{1}^{n-1} x_{2}^{n-2} \ldots x_{n-1}\right)\left(1+\gamma e_{2}\left(x_{1}, \ldots, x_{n+1}\right)\right)
$$

Proof First we observe that

$$
\begin{equation*}
\partial_{k}\left(x_{k}\left(1+\gamma e_{2}\left(x_{1}, \ldots, x_{k}\right)\right)\right)=1+\gamma e_{2}\left(x_{1}, \ldots, x_{k+1}\right) \tag{6}
\end{equation*}
$$

which can be shown by a straightforward computation using the identities

$$
\begin{aligned}
e_{2}\left(x_{1}, \ldots, x_{k}\right) & =e_{2}\left(x_{1}, \ldots, x_{k-1}\right)+x_{k} e_{1}\left(x_{1}, \ldots, x_{k-1}\right) \\
e_{2}\left(x_{1}, \ldots, x_{k+1}\right) & =e_{2}\left(x_{k}, x_{k+1}\right)+e_{2}\left(x_{1}, \ldots, x_{k-1}\right)+e_{1}\left(x_{1}, \ldots, x_{k-1}\right) e_{1}\left(x_{k}, x_{k+1}\right)
\end{aligned}
$$

Now we prove the formula by induction on $n$. The case $n=1$ is obvious. If $n>1$, by induction hypothesis, we have

$$
\partial_{n} \ldots \partial_{1}\left(x_{1}^{n} x_{2}^{n-1} \ldots x_{n}\right)=\left(x_{1}^{n-1} x_{2}^{n-2} \ldots x_{n-1}\right) \partial_{n}\left(x_{n}\left(1+e_{2}\left(x_{1}, \ldots, x_{n}\right)\right)\right)
$$

Thus the claim follows from Eq. 6.
Lemma 3.10 Modulo $\mathbb{S}_{N}$, we have

$$
\sum_{k=n+1}^{N-1} e_{2}\left(x_{1}, \ldots, x_{k}\right)=-\sum_{i=n+1}^{N-1}(i-n) x_{i+1} e_{1}\left(x_{1}, \ldots, x_{i}\right) .
$$

Proof Let us begin by recalling the following identity of elementary symmetric polynomials

$$
e_{2}\left(x_{1}, \ldots, x_{N}\right)=e_{2}\left(x_{1}, \ldots, x_{k}\right)+\sum_{i=k}^{N-1} x_{i+1} e_{1}\left(x_{1}, \ldots, x_{i}\right)
$$

Thus modulo $\mathbb{S}_{N}$, it follows that

$$
\sum_{k=n+1}^{N-1} e_{2}\left(x_{1}, \ldots, x_{k}\right)=-\sum_{k=n+1}^{N-1} \sum_{i=k}^{N-1} x_{i+1} e_{1}\left(x_{1}, \ldots, x_{i}\right)=-\sum_{i=n+1}^{N-1} \sum_{k=n+1}^{i} x_{i+1} e_{1}\left(x_{1}, \ldots, x_{i}\right)
$$

The right hand side is the desired formula.
For $w_{0}^{(n)} \in S_{n}$ the corresponding Schubert variety $X_{(n)}^{w_{n}^{(n)}}$ in $F l_{n}$ is a point, and so is the unique Bott-Samelson variety $Y_{1}^{(n)}$. In $I_{2}^{*}\left(F l_{n}\right)$ we have

$$
\left[Y_{1}^{(n)} \rightarrow F l_{n}\right]=x_{1}^{n-1} x_{2}^{n-2} \ldots x_{n-1}
$$

The stable Bott-Samelson class $\mathrm{BS}_{w_{0}^{(n)}}^{1}$ introduced in Theorem 3.3 is given by the sequence

$$
\left[Y_{\underline{c}^{(N-1)} \ldots \underline{c}^{(n)}}^{(N)} \rightarrow F l_{N}\right] \in I_{2}^{*}\left(F l_{N}\right), \quad N \geq n
$$

By Lemmas 3.9 and 3.10, we can identify a formal power series representing this class in the ring $\mathcal{R}_{\mathbb{I}}$ as follows.

Theorem 3.11 In $\mathcal{R}_{\mathbb{I}}$, we have

$$
\begin{equation*}
\mathrm{BS}_{w_{0}^{(n)}}^{1}(x)=\left(\prod_{i=1}^{n-1} x_{i}^{n-i}\right)\left(1-\gamma \sum_{i=n+1}^{\infty}(i-n) x_{i+1} e_{1}\left(x_{1}, \ldots, x_{i}\right)\right) . \tag{7}
\end{equation*}
$$

Proof In view of Lemma 3.5, we can compute the class $\left[Y_{\underline{c}^{(N-1)} \ldots \underline{c}^{(n)}}^{(N)} \rightarrow F l_{N}\right]$ via divided difference operators:

$$
\begin{aligned}
{\left[Y_{\underline{\underline{c}}^{(N-1)} \ldots \underline{c}^{(n)}}^{(N)} \rightarrow F l_{N}\right] } & =\partial_{\underline{c}^{(n)}} \ldots \partial_{\underline{c}^{(N-1)}}\left[Y_{1}^{(N)} \rightarrow F l_{N}\right] \\
& =\partial_{\underline{c}^{(n)}} \ldots \partial_{\underline{c}^{(N-1)}}\left(x_{1}^{N-1} x_{2}^{N-2} \ldots x_{N-1}\right)
\end{aligned}
$$

Consecutive applications of Lemma 3.9 give

$$
\left[Y_{\underline{c}^{(N-1)} \ldots \underline{c}^{(n)}}^{(N)} \rightarrow F l_{N}\right]=\left(\prod_{i=1}^{n-1} x_{i}^{n-i}\right)\left(1+\gamma \sum_{k=n+1}^{N-1} e_{2}\left(x_{1}, \ldots, x_{k}\right)\right) .
$$

By Lemma 3.10 we can rewrite this expression (modulo $\mathbb{S}_{N}$ ) as:

$$
\begin{equation*}
\left[Y_{\underline{c}^{(N-1)} \ldots \underline{c}^{(n)}}^{(N)} \rightarrow F l_{N}\right]=\left(\prod_{i=1}^{n-1} x_{i}^{n-i}\right)\left(1-\gamma \sum_{i=n+1}^{N-1}(i-n) x_{i+1} e_{1}\left(x_{1}, \ldots, x_{i}\right)\right) \tag{8}
\end{equation*}
$$

The right hand side of (7) is well-defined as an element of $\mathbb{I}[x]_{b d}$ and it projects to (8) for all $N \geq n$. This completes the proof.

In view of Corollary 3.6, all stable Bott-Samelson classes can be computed from (7) by applying divided difference operators. More precisely, pick $w \in W_{n}$ and $\underline{v} \in \underline{W}_{n}$ such that $v=w_{0}^{(n)} w$. Then

$$
\mathrm{BS}{\underset{w}{v}}_{\underline{v}}=\partial_{\underline{v}} \mathrm{BS}_{w_{0}^{(n)}}^{1} .
$$

Moreover, since the second factor of (7) is symmetric in $x_{1}, \ldots, x_{n}$, one simply has to identify $\partial_{\underline{v}}\left(x_{1}^{n-1} x_{2}^{n-2} \ldots x_{n-1}\right)$. That is,

$$
\mathrm{BS} \mathrm{~S}_{w}^{v}=\mathcal{B}_{n}(x) \partial_{\underline{v}}\left(x_{1}^{n-1} x_{2}^{n-2} \ldots x_{n-1}\right)
$$

where we denote

$$
\mathcal{B}_{n}(x):=1-\gamma \sum_{i=n+1}^{\infty}(i-n) x_{i+1} e_{1}\left(x_{1}, \ldots, x_{i}\right)
$$

Based on this, we will now obtain explicit closed formulas for the power series representing the stable Bott-Samelson classes associated to dominant permutations.

Definition 3.12 For a permutation $w \in S_{n}$, consider a $n \times n$ grid with dots in the boxes $(i, w(i))$. The diagram of $w$ is the set of boxes that remain after deleting boxes weakly east and south of each dot. A permutation $w$ is called dominant if its diagram is located at the NW corner of the grid, and coincides with a Young diagram of a partition $\lambda=\left(\lambda_{1}, \ldots, \lambda_{r}\right)$ with $\lambda_{i} \leq n-i$. For a given such partition $\lambda$, there is
a unique dominant permutation $w_{\lambda} \in S_{n}$. For example, the longest element $w_{0}^{(n)}$ is dominant and its associated partition is $\rho:=(n-1, n-2, \ldots, 2,1)$.

Let $T$ be the standard tableau of $\rho$, i.e., the fillings of the boxes of the $i$-th row of $T$ are all $i$. One places $\lambda$ at the NW corner of $\rho$ with its boxes shaded. We order the antidiagonals starting from the the inner ones to the outer ones, i.e., the $i$-th anti-diagonal consists of boxes at $(a, b)$ with $a+b=n+2-i$. Let $m$ be the biggest number such that the $m$-th anti-diagonal contains unshaded boxes. Let $\underline{v}^{(i)}, 1 \leq i \leq m$ be the reduced word obtained by reading the numbers in the $i$-th anti-diagonal. Then $\underline{v}:=\underline{v}^{(1)} \ldots \underline{v}^{(m)}$ is a reduced word of $v=w_{0}^{(n)} w_{\lambda}$. For each $i$, let $\mathbf{x}_{k}^{(i)}\left(k=1, \ldots, a_{i}\right)$ be the orbits of $v^{(i)}$ in $\left\{x_{1}, \ldots, x_{n}\right\}$ with cardinality greater than 1.

Theorem 3.13 Let $w_{\lambda} \in S_{n}$ be the dominant permutation associated to the partition $\lambda=\left(\lambda_{1}, \ldots, \lambda_{r}\right)$. Let $\underline{v}:=\underline{v}^{(1)} \ldots \underline{v}^{(m)}$ be the reduced word of $v=w_{0}^{(n)} w_{\lambda}$ constructed in Definition 3.12. We have

$$
\operatorname{BS} \bar{w}_{\lambda}^{v}=x_{1}^{\lambda_{1}} \ldots x_{r}^{\lambda_{r}}\left(1+\gamma\left(\sum_{i=1}^{m} \sum_{k=1}^{a_{i}} e_{2}\left(\mathbf{x}_{k}^{(i)}\right)\right)\right) \mathcal{B}_{n}(x) .
$$

Proof We prove the formula by induction on $m$. If $m=1$, then by Lemma 3.9 we have

$$
\partial_{\underline{v}^{(1)}}\left(x_{1}^{n-1} x_{2}^{n-2} \ldots x_{n-1}\right)=x_{1}^{\lambda_{1}} \ldots x_{r}^{\lambda_{r}}\left(1+\gamma\left(\sum_{k=1}^{a_{1}} e_{2}\left(\mathbf{x}_{k}^{(1)}\right)\right)\right) .
$$

Now, let $m>1$. By the induction hypothesis, we have

$$
\partial_{\underline{v}^{(m)}} \ldots \partial_{\underline{v}^{(1)}}\left(x_{1}^{n-1} x_{2}^{n-2} \ldots x_{n-1}\right)=\partial_{\underline{v}^{(m)}}\left(x_{1}^{\lambda_{1}^{\prime}} \ldots x_{r}^{\lambda_{r}^{\prime}}\left(1+\gamma\left(\sum_{i=1}^{m-1} \sum_{k=1}^{a_{i}} e_{2}\left(\mathbf{x}_{k}^{(i)}\right)\right)\right)\right),
$$

where $\lambda^{\prime}=\lambda \cup(n-m, n-m-1, \ldots, 2,1$,$) . Since the unshaded boxes form a$ skew shape $\rho / \lambda$, it follows that $\underline{v}^{(m)}$ stabilizes the second factor, allowing it pass through $\partial_{\underline{v}^{(m)}}$ :

$$
\partial_{\underline{v}^{(m)}} \ldots \partial_{\underline{v}^{(1)}}\left(x_{1}^{n-1} x_{2}^{n-2} \ldots x_{n-1}\right)=\left(1+\gamma\left(\sum_{i=1}^{m-1} \sum_{k=1}^{a_{i}} e_{2}\left(\mathbf{x}_{k}^{(i)}\right)\right)\right) \cdot \partial_{\underline{v}^{(m)}}\left(x_{1}^{\lambda_{1}^{\prime}} \ldots x_{r}^{\lambda_{r}^{\prime}}\right) .
$$

Now the desired formula follows again from Lemma 3.9.
Example 3.14 Consider $w_{\lambda}=(53124) \in S_{5}$ where $\lambda=(4,2)$.

| 1 | 1 | 1 | 1 |
| :--- | :--- | :--- | :--- |
| 2 | 2 | 2 |  |
| 3 | 3 |  |  |
| 4 |  |  |  |
|  |  |  |  |
|  |  |  |  |

The reduced word $\underline{v}$ of $v=w_{0}^{(5)} w_{\lambda}$ is $\underline{v}=\left(s_{2} s_{3} s_{4}\right)\left(s_{3}\right)$. We have

$$
\left(\partial_{3}\right)\left(\partial_{4} \partial_{3} \partial_{2}\right)\left(x_{1}^{4} x_{2}^{3} x_{3}^{2} x_{4}\right)=\left(1+\gamma e_{2}([2,5]) \cdot \partial_{3}\left(x_{1}^{4} x_{2}^{2} x_{3}\right)=x_{1}^{4} x_{2}^{2}\left(1+\gamma\left(e_{2}^{[2,5]}+e_{2}^{[3,4]}\right)\right) .\right.
$$

Example 3.15 Consider $w_{\lambda}=(45123) \in S_{5}$ where $\lambda=(3,3)$.

| 1 | 1 | 1 | 1 |
| :---: | :---: | :---: | :---: |
| 2 | 2 | 2 |  |
| 3 | 3 |  |  |
| 4 |  |  |  |

The reduced word of $v=w_{0}^{(5)} w_{\lambda}$ is $\underline{v}=\left(s_{1} s_{3} s_{4}\right)\left(s_{3}\right)$. We have

$$
\left(\partial_{3}\right)\left(\partial_{4} \partial_{3} \partial_{1}\right)\left(x_{1}^{4} x_{2}^{3} x_{3}^{2} x_{4}\right)=\left(1+\gamma\left(e_{2}^{[1,2]}+e_{2}^{[3,5]}\right)\right) \partial_{3}\left(x_{1}^{3} x_{2}^{3} x_{3}\right)=x_{1}^{3} x_{2}^{3}\left(1+\gamma\left(e_{2}^{[1,2]}+e_{2}^{[3,5]}+e_{2}^{[3,4]}\right)\right) .
$$

Example 3.16 Consider $w_{\lambda}=(563412) \in S_{6}$ where $\lambda=(4,4,2,2)$.

| 1 | 1 | 1 | 1 | 1 |
| :---: | :---: | :---: | :---: | :---: |
| 2 | 2 | 2 | 2 |  |
| 3 | 3 | 3 |  |  |
| 4 | 4 |  |  |  |
| 5 |  |  |  |  |

The reduced word of $v=w_{0}^{(6)} w_{\lambda}$ is $\underline{v}=s_{1} s_{3} s_{5}$. We have

$$
\partial_{5} \partial_{3} \partial_{1}\left(x_{1}^{5} x_{2}^{4} x_{3}^{3} x_{4}^{2} x_{5}\right)=\left(x_{1} x_{2}\right)^{4}\left(x_{3} x_{4}\right)^{2}\left(1+\gamma\left(e_{2}^{[1,2]}+e_{2}^{[3,4]}+e_{2}^{[5,6]}\right)\right)
$$

## 4 Restriction of Bott-Samelson Classes

In this section we generalise the restriction formula in Theorem 3.3 of the previous section. Namely, we will prove a formula for the product of the cobordism class of any Bott-Samelson resolution with the class $\left[F l_{n-1} \rightarrow F l_{n}\right.$ ]. In order to simplify the proof we will use another equivalent definition of Bott-Samelson resolutions.

### 4.1 Bott-Samelson Resolution Revisited

In this section, we provide another construction of the Bott-Samelson variety $X_{\underline{w}}$ associated to a word $\underline{w}$ by viewing it as a configuration space. This description will be better suited for our purposes.

Definition 4.1 Let $\underline{w}=s_{i_{1}} \ldots s_{i_{r}}$ be a word in $\underline{W}_{n+1}$.
(1) For $a \in[1, n+1]$, define $\mathrm{LO}_{\underline{w}}(a)$, the last occurence of $a$ in $\underline{w}$, by

$$
\mathrm{LO}_{\underline{w}}(a)=\sup \left\{k \in[1, r] \mid i_{k}=a\right\} .
$$

Note that if the above set is empty, then $\mathrm{LO}_{\underline{w}}(a)=-\infty$.
(2) If $\mathrm{LO}_{\underline{w}}(a)=-\infty$, then we set $V_{\mathrm{LO}_{\underline{w}}(a)}=\left\langle e_{1}, \ldots, e_{a}\right\rangle$.
(3) For $k \in[1, r]$, define $\underline{w}[k]=s_{i_{1}} \ldots s_{i_{k}}$.
(4) For $k \in[1, r]$, the left and the right predecessors of $k$ in $\underline{w}$, denoted $\mathrm{LP}_{\underline{w}}$ and $\mathrm{RP}_{\underline{w}}$, are defined as:

$$
\begin{aligned}
& \mathrm{LP}_{\underline{w}}(k):=\mathrm{LO}_{\underline{w}[k]}\left(i_{k}-1\right), \\
& \operatorname{RP}_{\underline{w}}(k):=\mathrm{LO}_{\underline{w}[k]}\left(i_{k}+1\right) .
\end{aligned}
$$

Remark 4.2 Note that by (2) and (4) in the above definition, we have:

$$
\begin{aligned}
V_{\mathrm{LP}_{\underline{w}}(k)} & =\left\langle e_{1}, \ldots, e_{i_{k}-1}\right\rangle \text { if } \mathrm{LP}_{\underline{w}}(k)=-\infty \text { and } \\
V_{\mathrm{RP}_{\underline{w}}}(k) & =\left\langle e_{1}, \ldots, e_{i_{k}+1}\right\rangle \text { if } \mathrm{RP}_{\underline{w}}(k)=-\infty .
\end{aligned}
$$

Definition 4.3 Given a word $\underline{w}=s_{i_{1}} \ldots s_{i_{r}}$, define the Bott-Samelson variety $X_{\underline{w}}$ as follows:

$$
X_{\underline{w}}=\left\{\left(V_{k}\right)_{k \in[1, r]} \mid \operatorname{dim} V_{k}=i_{k} \text { and } V_{\mathrm{LP}_{\underline{w}}(k)} \subset V_{k} \subset V_{\operatorname{RP}_{\underline{w}}(k)}\right\} .
$$

Define a morphism $\pi_{\underline{w}}: X_{\underline{w}} \rightarrow F l_{n+1}$ by $\pi_{\underline{w}}\left(\left(V_{k}\right)_{k \in[1, r]}\right)=\left(V_{\mathrm{LO}_{\underline{w}}(a)}\right)_{a \in[1, n]}$. If $\underline{w}$ is reduced, then the map $\pi_{\underline{w}}$ is a proper birational morphism from $X_{\underline{w}}$ onto the Schubert variety $X_{w}$. In this reduced case, we often call $X_{\underline{w}}$ together with the map $\pi_{\underline{w}}$ a BottSamelson resolution.

Remark 4.4 There is a natural isomorphism between our two construction of the Bott-Samelson variety $X_{\underline{w}}$ and $Y_{\underline{w}}$. It is given by $\left(V_{k}\right)_{k \in[1, r]} \mapsto\left(U_{\bullet}^{(k)}\right)_{k \in[1, r]}$ with $U_{a}^{(k)}=V_{\mathrm{LO}_{\underline{w}[\mid]}(a)}$ for all $k \in[1, r]$ and $a \in[1, n]$.

Recall the following well known fact on Bott-Samelson varieties.
Lemma 4.5 The Bott-Samelson variety does not depend on the choice of a word modulo commuting relations. More precisely, if $\underline{w}=\underline{v}$ modulo commuting relations, then there is an isomorphism $f_{\underline{w}, \underline{v}}: X_{\underline{w}} \rightarrow X_{\underline{v}}$ such that the following diagram is commutative:


Proof It is enough to prove this result in the case in which $\underline{w}$ and $\underline{v}$ are obtained from each other by a unique commuting relation. The result then follows by induction on the number of commuting relations. Assume therefore that $\underline{w}=\underline{u}_{1} s_{a} s_{b} \underline{u}_{2}$ and $\underline{v}=$ $\underline{u}_{1} s_{b} s_{a} \underline{u}_{2}$ with $|a-b| \geq 2$. Let $r_{i}=\ell\left(\underline{u}_{i}\right)$ for $i \in[1,2]$ and set $r=\ell(\underline{w})=\ell(\underline{v})=$ $r_{1}+r_{2}+2$. Define the map $f_{\underline{w}, \underline{v}}: X_{\underline{w}} \rightarrow X_{\underline{v}}$ by $f_{\underline{w}, \underline{v}}\left(\left(V_{k}\right)_{k \in[1, r]}\right)=\left(W_{k}\right)_{k \in[1, r]}$ and $f_{\underline{\underline{v}}, \underline{w}}\left(\left(W_{k}\right)_{k \in[1, r]}\right)=\left(V_{k}\right)_{k \in[1, r]}$ with $W_{k}=V_{k}$ for $k \notin\left\{r_{1}+1, r_{1}+2\right\}$ and $W_{r_{1}+\epsilon}=$ $V_{r_{1}+3-\epsilon}$ for $\epsilon \in\{1,2\}$. These maps are inverses of each other and we only need to check that they indeed map $X_{\underline{w}}$ to $X_{\underline{v}}$ and $X_{\underline{v}}$ to $X_{\underline{w}}$ respectively. By symmetry, we only need to check this for $f_{\underline{w}, \underline{v}}$.

We prove that given $\left(V_{k}\right)_{k \in[1, r]} \in X_{\underline{w}}$, the condition $W_{\mathrm{LP}_{\underline{v}}(k)} \subset W_{k}$ is satisfied for all $k$. The other inclusion $W_{k} \subset W_{\mathrm{RP}_{\underline{v}}(k)}$ is obtained by similar arguments. First note that we have the following relations:

$$
\begin{array}{ll}
\mathrm{LP}_{\underline{v}}(k)=\mathrm{LP}_{\underline{w}}(k) & \text { for } k, \mathrm{LP}_{\underline{w}}(k) \notin\left\{r_{1}+1, r_{1}+2\right\} \\
\mathrm{LP}_{\underline{v}}(k)=r_{1}+3-\epsilon & \text { for } \mathrm{LP}_{\underline{w}}(k)=r_{1}+\epsilon \text { and } \epsilon \in\{1,2\} \\
\mathrm{LP}_{\underline{v}}\left(r_{1}+\epsilon\right)=\mathrm{LP}_{\underline{w}}\left(r_{1}+3-\epsilon\right) & \text { for } \epsilon \in\{1,2\}
\end{array}
$$

For $k$ such that $k, \mathrm{LP}_{\underline{w}}(k) \notin\left\{r_{1}+1, r_{1}+2\right\}$, we have $W_{\mathrm{LP}_{\underline{v}}(k)}=W_{\mathrm{LP}_{\underline{w}}(k)}=V_{\mathrm{LP}_{\underline{w}}(k)} \subset$ $V_{k}=W_{k}$. For $\mathrm{LP}_{\underline{w}}(k)=r_{1}+\epsilon$ with $\epsilon \in\{1,2\}$, note that $k \notin\left\{r_{1}+1, r_{1}+2\right\}$ thus we have $W_{\mathrm{LP}_{\underline{v}}(k)}=W_{r_{1}+3-\epsilon}=V_{r_{1}+\epsilon}=V_{\mathrm{LP}_{\underline{w}}(k)} \subset V_{k}=W_{k}$. Finally, for $k=r_{1}+\epsilon$ with $\epsilon \in\{1,2\}$, note that $\operatorname{LP}_{\underline{w}}\left(r_{1}+3-\epsilon\right) \notin\left\{r_{1}+1, r_{1}+2\right\}$ thus we have $W_{\mathrm{LP}_{\underline{v}}(k)}=W_{\mathrm{LP}_{\underline{v}}\left(r_{1}+\epsilon\right)}=W_{\mathrm{LP}_{\underline{p}}\left(r_{1}+3-\epsilon\right)}=V_{\mathrm{LP}_{\underline{w}}\left(r_{1}+3-\epsilon\right)} \subset V_{r_{1}+3-\epsilon}=W_{r_{1}+\epsilon}=W_{k}$.
Furthermore we have:

$$
\begin{array}{ll}
\operatorname{LO}_{\underline{w}}(a)=\operatorname{LO}_{\underline{v}}(a) & \text { if } \mathrm{LO}_{\underline{w}}(a) \notin\left\{r_{1}+1, r_{1}+2\right\} \\
\operatorname{LO}_{\underline{w}}(a)=\mathrm{LO}_{\underline{v}}(a)+3-\epsilon & \text { if } \mathrm{LO}_{\underline{w}}(a)=r_{1}+\epsilon \text { for } \epsilon \in\{1,2\},
\end{array}
$$

so we easily see that we have $\pi_{\underline{v}} \circ f_{\underline{w}, \underline{v}}=\pi_{\underline{w}}$ and $\pi_{\underline{w}} \circ f_{\underline{v}, \underline{w}}=\pi_{\underline{v}}$.

### 4.2 Fiber Product with a Subflag

We now prove a fiber product formula for Bott-Samelson resolutions.
Define

$$
\begin{equation*}
\mathbf{F}_{n}=\left\{U_{\bullet} \in F l_{n+1} \mid U_{n}=\left\langle e_{2}, \ldots, e_{n+1}\right\rangle=F_{n}^{(n+1)}\right\} \tag{9}
\end{equation*}
$$

We can easily see that $\mathbf{F}_{n}$ coincides with the opposite Schubert variety $X^{c}$ in $F l_{n+1}$ where $c:=c^{(n)}=s_{1} \ldots s_{n}$ is the Coxeter element. Therefore, from a well-known
fact, we have that for $w \in W_{n+1}$

$$
\begin{equation*}
X_{w} \cap \mathbf{F}_{n} \neq \emptyset \text { if and only if } c \leq w . \tag{10}
\end{equation*}
$$

Definition 4.6 For $\underline{u} \in \underline{W}_{n+1}^{1}$ with $\underline{u}=s_{i_{1}} \ldots s_{i_{r}}$, define $c^{-1}(\underline{u}) \in \underline{W}_{n+1}^{n}$ by

$$
c^{-1}(\underline{u})=s_{c^{-1}\left(i_{1}\right)} \ldots s_{c^{-1}\left(i_{r}\right)}=s_{i_{1}-1} \ldots s_{i_{r}-1},
$$

where we observe that for each $k \in[2, n+1]$ one has $c^{-1}(k)=k-1 \in[1, n]$.
Definition 4.7 Denote by $c$ the isomorphism $c: \mathbf{k}^{n+1} \rightarrow \mathbf{k}^{n+1}$ defined by $c\left(e_{i}\right)=$ $e_{c(i)}$ for all $i \in[1, n+1]$.

1. The map $c$ induces an isomorphism $c: F l_{n} \rightarrow \mathbf{F}_{n} \subset F l_{n+1}$.
2. For $\underline{w} \in W_{n+1}^{n}$ with $\ell(\underline{w})=r$, define

$$
c\left(X_{\underline{w}}\right)=\left\{\left(c\left(V_{k}\right)\right)_{k \in[1, r]} \mid\left(V_{k}\right)_{k \in[1, r]} \in X_{\underline{w}}\right\}
$$

and $c\left(\pi_{\underline{w}}\right): c\left(X_{\underline{w}}\right) \rightarrow c\left(X_{w}\right) \subset c\left(F l_{n}\right)=\mathbf{F}_{n}$, so that the following diagram is commutative:


Lemma 4.8 For $\underline{u} \in \underline{W}_{n+1}^{1}$ and $\underline{v} \in \underline{W}_{n+1}^{n}$, define $\underline{w}=\underline{u} \underline{c} \underline{v}$ and $\underline{w}^{\prime}=c^{-1}(\underline{u}) \underline{v}$. Let $r_{1}=\ell(\underline{u})$ and $r_{2}=\ell(\underline{v})$. For $\nsupseteq \in[1, n-1]$, we have the following equalities:
(1) If $k \leq r_{1}$, then $\mathrm{LO}_{\underline{w^{\prime}}[k]}(\ngtr)=\mathrm{LO}_{\underline{w}[k]}(\ngtr+1)$;
(2) If $k \geq r_{1}+1$ and $\mathrm{LO}_{\underline{v}\left[k-r_{1}\right]}(\ngtr) \neq-\infty$, then $\mathrm{LO}_{\underline{w}^{\prime}[k]}(\ngtr)=\mathrm{LO}_{\underline{w}[k+n]}(\ngtr)-n=$ $\mathrm{LO}_{\underline{v}\left[k-r_{1}\right]}(\ngtr)+r_{1} \geq r_{1}+1$;
(3) If $k \geq r_{1}+1$ and $\mathrm{LO}_{\underline{v}\left[k-r_{1}\right]}(\ngtr)=-\infty$, then $\mathrm{LO}_{\underline{w^{\prime}}[k]}(\ngtr)=\mathrm{LO}_{\underline{u}}(\ngtr+1)=$ $\operatorname{RP}_{\underline{w}[k+n]}\left(r_{1}+\ngtr\right) \leq r_{1}$ and $\mathrm{LO}_{\underline{w}[k+n]}(\ngtr)=r_{1}+\ngtr$.

Proof If $k \leq r_{1}$, then $\underline{w}^{\prime}[k]=c^{-1}(\underline{w}[k])$ and the result follows. Assume now that $k \geq$ $r_{1}+1$. If $\mathrm{LO}_{\underline{v}\left[k-r_{1}\right]}(\ngtr) \neq-\infty$, then the last occurence of $\ngtr$ in $\underline{w}^{\prime}[k]$ is obtained at a letter of $\underline{v}\left[k-r_{1}\right]$ so that $\mathrm{LO}_{w u^{\prime}[k]}(\ngtr)=\mathrm{LO}_{\underline{v}\left[k-r_{1}\right]}(\ngtr)+r_{1}=\mathrm{LO}_{\underline{w}[k+n]}(\ngtr)-n \geq$ $r_{1}$. If $\mathrm{LO}_{\underline{v}\left[k-r_{1}\right]}(\ngtr)=-\infty$, the last occurence of $\ngtr$ in $\underline{w}^{\prime}[k]$ is the last occurence of $\ngtr$ in $c^{-1}(\underline{u})$ and we get the equalities $\mathrm{LO}_{\underline{w}^{\prime}[k]}(\ngtr)=\mathrm{LO}_{c^{-1}(\underline{u})}(\ngtr)=\mathrm{LO}_{\underline{u}}(\ngtr+1) \leq r_{1}$. On the other hand, the last occurence of $\ngtr$ in $\underline{w}[k+n]$ is obtained as the $\ngtr$-th letter in $\underline{c}$ thus $\mathrm{LO}_{\underline{w}[k+n]}(\ngtr)=r_{1}+\not \nsupseteq$. This also explains the equality $\mathrm{LO}_{\underline{u}}(\ngtr+1)=$ $\mathrm{RP}_{\underline{w}[k+n]}\left(r_{1}+\ngtr\right)$.

Corollary 4.9 For $\underline{u} \in \underline{W}_{n+1}^{1}$ and $\underline{v} \in \underline{W}_{n+1}^{n}$, define $\underline{w}=\underline{u c} \underline{v}$ and $\underline{w}^{\prime}=c^{-1}(\underline{u}) \underline{v}$. Let $r_{1}=\ell(\underline{u})$ and $r_{2}=\ell(\underline{v})$. We have the following alternatives:
(1) If $\mathrm{LO}_{\underline{v}}(\ngtr) \neq-\infty$, then $\mathrm{LO}_{\underline{w^{\prime}}}(\ngtr)=\mathrm{LO}_{\underline{w}}(\ngtr)-n=\mathrm{LO}_{\underline{v}}(\ngtr)+r_{1} \geq r_{1}+1$;
(2) $\operatorname{If} \mathrm{LO}_{\underline{v}}^{-}(\ngtr)=-\infty$, then $\mathrm{LO}_{\underline{w^{\prime}}}(\ngtr)=\mathrm{LO}_{\underline{u}}(\ngtr+1)=\mathrm{RP}_{\underline{w}}\left(r_{1}+\ngtr\right)$ and $\mathrm{LO}_{\underline{w}}(\ngtr)$ $=r_{1}+\ngtr$.

Proof Apply the previous lemma with $k=r_{1}+r_{2}$.
Corollary 4.10 For $\underline{u} \in \underline{W}_{n+1}^{1}$ and $\underline{v} \in \underline{W}_{n+1}^{n}$, define $\underline{w}=\underline{u c \underline{v}}$ and $\underline{w}^{\prime}=c^{-1}(\underline{u}) \underline{v}$. Let $r_{1}=\ell(\underline{u})$ and $r_{2}=\ell(\underline{v})$. We set $\operatorname{LP}_{\underline{x}}(-\infty)=-\infty$ and $\mathrm{RP}_{\underline{x}}(-\infty)=-\infty$ for any word $x$.
A. We have the following formulas for $\mathrm{LP}_{\underline{w}}$ and $\mathrm{LP}_{\underline{w}^{\prime}}$ :

1. If $a \in\left[1, r_{1}\right]$, then $\operatorname{LP}_{\underline{w^{\prime}}}(a)=\operatorname{LP}_{\underline{w}}(a) \leq r_{1}$.
2. If $a>r_{1}$ and $\operatorname{LP}_{\underline{v}}\left(a-r_{1}\right) \neq \infty$, then $\operatorname{LP}_{\underline{w^{\prime}}}(a)=\operatorname{LP}_{\underline{w}}(a+n)-n>r_{1}$.
3. If $a>r_{1}$ and $\operatorname{LP}_{\underline{v}}\left(a-r_{1}\right)=-\infty$, then $\operatorname{LP}_{\underline{w}}(a+n) \in\left[r_{1}+1, r_{1}+n\right] \cup\{-\infty\}$ and $\mathrm{LP}_{\underline{w}^{\prime}}(a)=\operatorname{RP}_{\underline{w}}\left(\operatorname{LP}_{\underline{w}}(a+n)\right) \leq r_{1}$.
B. We have the following formulas for $\mathrm{RP}_{\underline{w}}$ and $\mathrm{RP}_{\underline{w}^{\prime}}$ :
4. If $a \in\left[1, r_{1}\right]$, then $\operatorname{RP}_{\underline{w^{\prime}}}(a)=\operatorname{RP}_{\underline{w}}(a) \leq r_{1}$.
5. If $a>r_{1}$ and $\mathrm{RP}_{\underline{v}}\left(a-r_{1}\right) \neq \infty$, then $\mathrm{RP}_{\underline{w}^{\prime}}(a)=\mathrm{RP}_{\underline{w}}(a+n)-n>r_{1}$.
6. If $a>r_{1}$ and $\mathrm{RP}_{\underline{v}}\left(a-r_{1}\right)=-\infty$, then $\mathrm{RP}_{\underline{w}^{\prime}}(a)=\mathrm{RP}_{\underline{w}}\left(\operatorname{RP}_{\underline{w}}(a+n)\right) \leq r_{1}$ and $\operatorname{RP}_{\underline{w}}(a+n) \in\left[r_{1}+1, r_{1}+n\right] \cup\{-\infty\}$.
C. If $k \in\left[r_{1}+1, r_{1}+n\right]$, then $\mathrm{LP}_{\underline{w}}(k) \in\left[r_{1}+1, r_{1}+n\right] \cup\{-\infty\}$ Furthermore, for $k \in\left[r_{1}+1, r_{1}+n\right]$, we have the following formulas for $\mathrm{LP}_{\underline{w}}$ and $\mathrm{RP}_{\underline{w}}$ :
7. If $\mathrm{RP}_{\underline{w}}(k)>\mathrm{RP}_{\underline{w}}\left(\mathrm{LP}_{\underline{w}}(k)\right)$, then $\mathrm{RP}_{\underline{w}}\left(\mathrm{LP}_{\underline{w}}(k)\right)=\mathrm{LP}_{\underline{w}}\left(\mathrm{RP}_{\underline{w}}(k)\right)$.
8. If $\mathrm{RP}_{\underline{w}}(k)<\mathrm{RP}_{\underline{w}}\left(\mathrm{LP}_{\underline{w}}(k)\right)$, then $\mathrm{RP}_{\underline{w}}\left(\mathrm{RP}_{\underline{w}}\left(\mathrm{LP}_{\underline{w}}(k)\right)\right)=\mathrm{RP}_{\underline{w}}(k)$.

Proof Write $\underline{u}=s_{i_{1}} \ldots s_{i_{r_{1}}}$ and $\underline{v}=s_{i_{r_{1}+n+1}} \ldots s_{i_{r_{1}+n+r_{2}}}$ so that $w=s_{i_{1}} \ldots s_{i_{r}}$ with $r=$ $r_{1}+r_{2}+n$ and $s_{i_{r_{1}+k}}=s_{k}$ for $k \in[1, n]$. We have $\underline{w}^{\prime}=s_{j_{1}} \ldots s_{j_{r_{1}+r+2}}$ with

$$
j_{k}= \begin{cases}i_{k}-1 & \text { for } k \in\left[1, r_{1}\right] \\ i_{k+n} & \text { for } k \in\left[r_{1}+1, r_{1}+r_{2}\right]\end{cases}
$$

A.1. We have $\operatorname{LP}_{\underline{w^{\prime}}}(a)=\operatorname{LO}_{\underline{w}^{\prime}[a]}\left(j_{a}-1\right)=\mathrm{LO}_{\underline{w}[a]}\left(\left(j_{a}-1\right)+1\right)=\mathrm{LO}_{\underline{w}[a]}\left(j_{a}\right)$ $=\mathrm{LO}_{\underline{w}[a]}\left(i_{a}-1\right)=\mathrm{LP}_{\underline{w}}(a)$. By definition $\mathrm{LP}_{\underline{w^{\prime}}}(a)<a \leq r_{1}$.
 $n)-n \geq r_{1}+1$.
A.3. We have $\mathrm{LP}_{\underline{w^{\prime}}}(a)=\mathrm{LO}_{\underline{w^{\prime}}[a]}\left(j_{a}-1\right)=\operatorname{RP}_{\underline{w}[a+n]}\left(r_{1}+j_{a}-1\right) \leq r_{1}$ and $r_{1}+$ $j_{a}-1=\mathrm{LO}_{\underline{w}[a+n]}\left(j_{a}-1\right)$. $\overline{\mathrm{We}}$ get $\mathrm{LP}_{\underline{w}^{\prime}}(a)=\mathrm{RP}_{\underline{w}[a+n]}\left(\mathrm{LO}_{\underline{w}[a+n]}\left(j_{a}-1\right)\right)=$ $\operatorname{RP}_{\underline{w}[a+n]}\left(\mathrm{LO}_{\underline{w}[a+n]}\left(i_{a+n}-1\right)\right)=\operatorname{RP}_{\underline{w}[a+n]}\left(\operatorname{LP}_{\underline{w}}(a+n)\right)$.

For B.1. and B.2. use the proof of A.1. and A. 2 with RP in place of LP.
B.3. We have $\mathrm{RP}_{\underline{w^{\prime}}}(a)=\mathrm{LO}_{\underline{w^{\prime}}[a]}\left(j_{a}+1\right)=\mathrm{RP}_{\underline{w}[a+n]}\left(r_{1}+j_{a}+1\right) \leq r_{1}$ and $r_{1}+$ $j_{a}+1=\operatorname{LO}_{\underline{w}[a+n]}\left(j_{a}+1\right)$. We get $\operatorname{LP}_{\underline{w}^{\prime}}(a)=\operatorname{RP}_{\underline{w}[a+n]}\left(\operatorname{LO}_{\underline{w}[a+n]}\left(j_{a}+1\right)\right)=$ $\operatorname{RP}_{\underline{w}[a+n]}\left(\mathrm{LO}_{\underline{w}[a+n]}\left(i_{a+n}+1\right)\right)=\operatorname{RP}_{\underline{w}[a+n]}\left(\mathrm{RP}_{\underline{w}}(a+n)\right)$.
C. If $k \in\left[r_{1}+1, r_{1}+n\right]$, then the $k$-th letter of $\underline{w}$ is the $\left(k-r_{1}\right)$-th letter of $\underline{c}$ and we have $\mathrm{LP}_{\underline{w}}(k)=\mathrm{LP}_{\underline{c}}\left(k-r_{1}\right)=k-r_{1}-1$ for $k>r_{1}+1$ and $\mathrm{LP}_{\underline{w}}\left(r_{1}+1\right)=$
$-\infty$. Note that, if any of the two quatities $\mathrm{RP}_{\underline{w}}(k)$ or $\mathrm{RP}_{\underline{w}}\left(\mathrm{LP}_{\underline{w}}(k)\right)$ is finite, we have $\mathrm{RP}_{\underline{w}}(k) \neq \mathrm{RP}_{\underline{w}}\left(\mathrm{LP}_{\underline{w}}(k)\right)$.
C.1. Assume $\mathrm{RP}_{\underline{w}}(k)>\operatorname{RP}_{\underline{w}}\left(\mathrm{LP}_{\underline{w}}(k)\right)$. If $k=r_{1}+1$, then $i_{k}=1$ and $\mathrm{LP}_{\underline{w}}(k)=$ $-\infty$ thus $\operatorname{RP}_{\underline{w}}\left(\operatorname{LP}_{\underline{w}}(k)\right)=-\infty$. Furthermore, $i_{\mathrm{RP}_{\underline{w}}(k)}=i_{k}+1=2$ and since $\underline{u} \in$ $\underline{W}_{n+1}^{1}$ we have $\operatorname{LP}_{\underline{w}}\left(\operatorname{RP}_{\underline{w}}(k)\right)=\mathrm{LO}_{\underline{u}}(2)=-\infty$. Assume now $k>r_{1}+1$, then $\mathrm{LP}_{\underline{w}}(k)=k-1$, thus $\mathrm{RP}_{\underline{w}}\left(\mathrm{LP}_{\underline{w}}(k)\right)=\mathrm{RP}_{\underline{w}}(k-1)=\mathrm{LO}_{\underline{u}}\left(i_{k-1}+1\right)=\mathrm{LO}_{\underline{u}}\left(i_{k}\right)$. We have $i_{\mathrm{RP}_{\underline{w}}(k)}=i_{k}+1$ thus $\mathrm{LP}_{\underline{w}}\left(\mathrm{RP}_{\underline{w}}(k)\right)=\mathrm{LO}_{\underline{w}\left[\mathrm{RP}_{\underline{w}}(k)\right]}\left(i_{\mathrm{RP}_{\underline{w}}(k)}-1\right)=$ $\mathrm{LO}_{\underline{w}\left[\mathrm{RP}_{\underline{w}}(k)\right]}\left(i_{k}\right)$. Since $\mathrm{RP}_{\underline{w}}(k)<k$, we have that $\underline{w}\left[\mathrm{RP}_{\underline{w}}(k)\right]$ is a subword of $\underline{u}$ thus $\mathrm{LO}_{\underline{w}\left[\mathrm{RP}_{\underline{w}}(k)\right]}\left(i_{k}\right) \leq \mathrm{LO}_{\underline{u}}^{\underline{\underline{u}}}\left(i_{k}\right)=\mathrm{RP}_{\underline{w}}\left(\mathrm{LP}_{\underline{w}}(k)\right)$. On the other hand, since $\mathrm{RP}_{\underline{w}}(k)>$ $\operatorname{RP}_{\underline{w}}\left(\operatorname{LP}_{\underline{w}}(k)\right)$, we have $\mathrm{LO}_{\underline{w}\left[\mathrm{RP}_{\underline{w}}(k)\right]}\left(i_{k}\right) \geq \mathrm{LO}_{\underline{w}\left[\mathrm{RP}_{\underline{w}}\left(\mathrm{LP}_{\underline{w}}(k)\right]\right]}\left(i_{k}\right)=\mathrm{RP}_{\underline{w}}\left(\mathrm{LP}_{\underline{w}}(k)\right)$. The last equality holds since $i_{\mathrm{RP}_{\underline{w}}\left(\mathrm{LP}_{\underline{w}}(k)\right)}=i_{k}$.
C.2. Assume $\mathrm{RP}_{\underline{w}}(k)<\mathrm{RP}_{\underline{w}}\left(\mathrm{LP}_{\underline{w}}(k)\right)$. This implies $\mathrm{LP}_{\underline{w}}(k) \neq-\infty$ thus $k>r_{1}+1$. We have $\mathrm{RP}_{\underline{w}}(k)=\mathrm{LO}_{\underline{w}[k]}\left(i_{k}+1\right)$ and $\mathrm{RP}_{\underline{w}}\left(\mathrm{RP}_{\underline{w}}\left(\mathrm{LP}_{\underline{w}}(k)\right)\right)=$ $\mathrm{LO}_{\underline{w}\left[\mathrm{RP}_{\underline{w}}\left(\mathrm{LP}_{\underline{w}}(k)\right)\right]}\left(i_{k}+1\right)$. Since $\operatorname{RP}_{\underline{w}}\left(\operatorname{LP}_{\underline{w}}(k)\right)<k$, then $\underline{w}\left[\mathrm{RP}_{\underline{w}}\left(\mathrm{LP}_{\underline{w}} \overline{(k)}\right)\right]$ is a subword of $\underline{w}[k]$ thus $\mathrm{LO}_{\underline{w}\left[\mathrm{RP}_{\underline{w}}\left(\mathrm{LP}_{\underline{w}}(k)\right)\right]}\left(i_{k}+1\right) \leq \mathrm{LO}_{\underline{w}[k]}\left(i_{k}+1\right)$. On the other hand, since $\operatorname{RP}_{\underline{w}}(k)<\operatorname{RP}_{\underline{w}}\left(\mathrm{LP}_{\underline{w}}^{-}(k)\right)$, we have $\mathrm{LO}_{\underline{w}\left[\mathrm{RP}_{\underline{w}}\left(\mathrm{LP}_{\underline{w}}(k)\right)\right]}\left(i_{k}+1\right) \geq \mathrm{LO}_{\mathrm{RP}_{\underline{w}}[k]}$ $\left(i_{k}+1\right)$.

Theorem 4.11 Let $\underline{w}$ be a reduced word and $w \in W$ the associated element.
(1) If $w \nsupseteq c$, then $X_{\underline{w}} \times_{F l_{n+1}} \mathbf{F}_{n}$ is empty.
(2) If $w \geq c$, there exist $\underline{u} \in \underline{W}_{n+1}^{1}$ and $\underline{v} \in \underline{W}_{n+1}^{n}$ such that $\underline{w}=\underline{u} \underline{c} \underline{v}$ modulo commuting relations and we have an isomorphism of $F l_{n+1}$-varieties $X_{\underline{w}} \times F l_{n+1}$ $\mathbf{F}_{n} \simeq c\left(X_{c^{-1}(\underline{u}) \underline{v}}\right)$.

Proof (1) is clear since the condition implies $X^{w} \cap \mathbf{F}_{n}=\emptyset$. We prove (2). By Proposition 2.9, we can write $\underline{w}=\underline{u} \underline{c} \underline{v}$. Let $\ell(w)=r$, and $\ell(u)=r_{1}, \ell(v)=r_{2}$ so that $r=r_{1}+r_{2}+n$. Since the obvious inclusion $\mathbf{F}_{n} \hookrightarrow F l_{n+1}$ is a closed embedding, we can view $X_{\underline{w}} \times{ }_{F l_{n+1}} \mathbf{F}_{n}$ as the closed subvariety of $X_{\underline{w}}$ given as follows:

$$
X_{\underline{w}} \times_{F l_{n+1}} \mathbf{F}_{n}=\left\{\left(V_{k}\right)_{k \in[1, r]} \in X_{\underline{w}} \mid V_{r_{1}+n}=\left\langle e_{2}, \ldots, e_{n+1}\right\rangle\right\} .
$$

Define the map $f: X_{\underline{w}} \times{ }_{F l_{n+1}} \mathbf{F}_{n} \rightarrow c\left(X_{c^{-1}(\underline{u} \underline{v}}\right)$ by $f\left(\left(V_{k}\right)_{k \in[1, r]}\right)=\left(H_{a}\right)_{a \in\left[r_{1}+r_{2}\right]}$ with

$$
H_{a}= \begin{cases}V_{a} \cap\left\langle e_{2}, \ldots, e_{n+1}\right\rangle & \text { for } a \in\left[1, r_{1}\right] \\ V_{a+n} & \text { for } a \in\left[r_{1}+1, r_{1}+r_{2}\right] .\end{cases}
$$

Define the map $g: c\left(X_{c^{-1}(\underline{u}) \underline{v}}\right) \rightarrow X_{\underline{w}} \times_{F l_{n+1}} \mathbf{F}_{n}$ by $g\left(\left(H_{a}\right)_{a \in\left[r_{1}+r_{2}\right]}\right)=\left(V_{k}\right)_{k \in[1, r]}$ with

$$
V_{k}= \begin{cases}H_{k}+\left\langle e_{1}\right\rangle & \text { for } k \in\left[1, r_{1}\right] \\ V_{\mathrm{RP}_{\underline{w}}(k)} \cap\left\langle e_{2}, \ldots, e_{n+1}\right\rangle & \text { for } k \in\left[r_{1}+1, r_{1}+n\right] \\ H_{k-n} & \text { for } k \in\left[r_{1}+n+1, r\right] .\end{cases}
$$

Set $\underline{w}^{\prime}=c^{-1}(\underline{u}) \underline{v} \in \underline{W}_{n+1}^{n}$. Write $\underline{u}=s_{i_{1}} \ldots s_{i_{r_{1}}}$ and $\underline{v}=s_{i_{r_{1}+n+1}} \ldots s_{i_{r_{1}+n+r_{2}}}$ so that $w=s_{i_{1}} \ldots s_{i_{r}}$ with $s_{i_{r_{1}+k}}=s_{k}$ for $k \in[1, n]$. We have $\underline{w}^{\prime}=s_{j_{1}} \ldots s_{j_{r_{1}+r+2}}$ with

$$
j_{k}= \begin{cases}i_{k}-1 & \text { for } k \in\left[1, r_{1}\right] \\ i_{k+n} & \text { for } k \in\left[r_{1}+1, r_{1}+r_{2}\right] .\end{cases}
$$

We first prove that these maps are well defined. We start with $f$. Note that if $\left(V_{k}\right)_{k \in[1, r]} \in X_{\underline{w}} \times_{F l_{n+1}} \mathbf{F}_{n}$, then

$$
\begin{array}{ll}
V_{k} \supset\left\langle e_{1}\right\rangle & \text { for } k \in\left[1, r_{1}\right] \\
V_{k} \subset\left\langle e_{2}, \ldots, e_{n+1}\right\rangle & \text { for } k \in\left[r_{1}+1, r_{1}+r_{2}-\right. \\
V_{k}=V_{\mathrm{RP}_{\underline{w}}(k)}\left(\left\langle e_{2}, \ldots, e_{n+1}\right\rangle\right. & \text { for } k \in\left[r_{1}+1, r_{1}+n\right] .
\end{array}
$$

Indeed, since $\underline{u} \in \underline{W}_{1}$, we have $\mathrm{LP}_{\underline{w}}(k)=-\infty$ for any $k \in\left[1, r_{1}\right]$ with $i_{k}=2$ implying our first claim. Furthermore, since $\underline{v} \in \underline{W}_{n+1}$, by the same type of arguments we have the equality $V_{r_{1}+n}=\left\langle e_{2}, \ldots, e_{n+1}\right\rangle$ proving the second claim. In particular for $k \in\left[r_{1}+1, r_{1}+n\right]$, we have $V_{k} \subset V_{\mathrm{RP}_{\underline{w}}(k)} \cap\left\langle e_{2}, \ldots, e_{n+1}\right\rangle$ and $\left\langle e_{1}\right\rangle \subset V_{\mathrm{RP}_{\underline{w}}(k)}$. Since $\operatorname{dim} V_{k}=\operatorname{dim} V_{\mathrm{RP}_{w}(k)}-1$ this proves the last equality.

We check that $\operatorname{dim} \bar{H}_{a}=j_{a}$. For $a \in\left[1, r_{1}\right]$, we have $\operatorname{dim} H_{a}=\operatorname{dim}\left(V_{a} \cap\right.$ $\left.\left\langle e_{2}, \ldots, e_{n+1}\right\rangle\right)=\operatorname{dim} V_{a}-1=i_{a}-1=j_{a}$, where the second equality holds since $\underline{u} \in \underline{W}_{n+1}^{1}$, therefore $\left\langle e_{1}\right\rangle \subset V_{a}$. For $a \in\left[r_{1}+1, r_{1}+r_{2}\right]$, we have $\operatorname{dim} H_{a}=$ $\operatorname{dim} V_{a+n}=i_{a+n}=j_{a}$.

We now check the inclusions $H_{\mathrm{LP}_{w^{\prime}}(a)} \subset H_{a} \subset H_{\mathrm{RP}_{w^{\prime}}(a)}$. We start with the inclusions $H_{\mathrm{LP}_{\underline{w}^{\prime}}(a)} \subset H_{a}$. For $a \in\left[1, r_{1}\right]$, then $\operatorname{LP}_{\underline{w}^{\prime}}(a)=\mathrm{LP}_{\underline{w}}(a) \leq r_{1}$ and we have $H_{\mathrm{LP}_{w^{\prime}}(a)}=H_{\mathrm{LP}_{\underline{w}}(a)}=V_{\mathrm{LP}_{\underline{w}}(a)} \cap\left\langle e_{2}, \ldots, e_{n+1}\right\rangle \subset V_{a} \cap\left\langle e_{2}, \ldots, e_{n+1}\right\rangle=$ $H_{a}$. For $a \in\left[r_{1}+1, r_{1}+r_{2}\right]$ and $\mathrm{LP}_{\underline{v}}\left(a-r_{1}\right) \neq-\infty$, then $\mathrm{LP}_{\underline{w}^{\prime}}(a)=\mathrm{LP}_{\underline{w}}(a+$ $n)-n \geq r_{1}+1$ and we have $H_{\mathrm{LP}_{w^{\prime}}(a)}=H_{\mathrm{LP}_{\underline{w}}(a+n)-n}=V_{\mathrm{LP}_{\underline{w}}(a+n)} \subset V_{a+n}=H_{a}$. For $a \in\left[r_{1}+1, r_{1}+r_{2}\right]$ and $\operatorname{LP}_{\underline{v}}(a)=-\infty$, then $\operatorname{LP}_{\underline{w^{\prime}}}(a)=\operatorname{RP}_{\underline{w}}\left(\mathrm{LP}_{\underline{w}}(a+\right.$ $n)) \leq r_{1}$ and $\operatorname{LP}_{\underline{w}}(a+n) \in\left[r_{1}+1, r_{1}+n\right] \cup\{-\infty\}$. If $\operatorname{LP}_{\underline{w}}(a+n)=-\infty$, then $\mathrm{LP}_{\underline{w}^{\prime}}(a)=-\infty$ and $H_{\mathrm{LP}_{w^{\prime}}(a)}=0$, so the inclusion holds. Otherwise, we have $H_{\mathrm{LP}_{w^{\prime}}(a)}=H_{\operatorname{RP}\left(\mathrm{LP}_{\underline{w}}(a)\right)}=V_{\operatorname{RP}\left(\mathrm{LP}_{\underline{w}}(a+n)\right)} \cap\left\langle e_{2}, \ldots, e_{n+1}\right\rangle=V_{\mathrm{LP}_{\underline{w}}(a+n)} \subset V_{a}=H_{a}$.

We prove the inclusions $\bar{H}_{k} \subset H_{\mathrm{RP}_{\underline{w}^{\prime}}(a)}$. For $a \in\left[1, r_{1}\right]$, then $\mathrm{RP}_{\underline{w}^{\prime}}(a)=$ $\mathrm{RP}_{\underline{w}}(a) \leq r_{1}$ and we have $H_{a}=V_{a} \cap\left\langle e_{2}, \ldots, e_{n+1}\right\rangle \subset V_{\mathrm{RP}_{\underline{w}}(a)} \cap\left\langle e_{2}, \ldots, e_{n+1}\right\rangle=$ $H_{\mathrm{RP}_{\underline{w}}(a)}=H_{\mathrm{RP}_{w^{\prime}}(a)}$. For $a \in\left[r_{1}+1, r_{1}+r_{2}\right]$ and $\operatorname{RP}_{\underline{v}}\left(a-r_{1}\right) \neq-\infty$, then $\mathrm{RP}_{\underline{w}^{\prime}}(a)=\mathrm{RP}_{\underline{w}}(a+n)-n \geq r_{1}+1$ and we have $H_{a}=V_{a+n} \subset V_{\mathrm{RP}_{\underline{w}}(a+n)}=$ $H_{\mathrm{RP}_{\underline{w}}(a+n)-n}=H_{\mathrm{RP}_{\underline{w}^{\prime}}(a)}$. For $a \in\left[r_{1}+1, r_{1}+r_{2}\right]$ and $\operatorname{RP}_{\underline{v}}(a)=-\infty$, then $\operatorname{RP}_{\underline{w^{\prime}}}(a)=\mathrm{RP}_{\underline{w}}\left(\mathrm{RP}_{\underline{w}}(a+n)\right) \leq r_{1}$ and $\mathrm{RP}_{\underline{w}}(a+n) \in\left[r_{1}+1, r_{1}+n\right] \cup\{-\infty\}$. If $\operatorname{RP}_{\underline{w}}(a+n)=-\infty$, then $\mathrm{RP}_{\underline{w}^{\prime}}(a)=-\infty$ and $H_{\mathrm{LP}_{w^{\prime}}(a)}=\left\langle e_{1}, \ldots, e_{n+1}\right\rangle$, so the inclusion holds. Otherwise, we have $H_{a}=V_{a+n} \subset V_{\mathrm{RP}_{\underline{w}}(a+n)}=V_{\mathrm{RP}_{\left(\mathrm{RP}_{\underline{w}}(a+n)\right)} \cap} \cap$ $\left\langle e_{2}, \ldots, e_{n+1}\right\rangle=H_{\mathrm{RP}_{\underline{w}}\left(\operatorname{RP}_{\underline{w}}(a+n)\right)}=H_{\mathrm{RP}_{w^{\prime}}(a)}$.

We now prove that $g$ is well defined. Note that for all $a$, we have $H_{a} \subset$ $\left\langle e_{2}, \ldots, e_{n+1}\right\rangle$. We first check the equalities $\operatorname{dim} V_{k}=i_{k}$ for all $k \in[1, r]$. For $k \in\left[1, r_{1}\right]$, we have $\operatorname{dim} V_{k}=\operatorname{dim}\left(H_{k}+\left\langle e_{1}\right\rangle\right)=\operatorname{dim} H_{k}+1=j_{k}+1=i_{k}$, where the second equality holds since $H_{k} \subset\left\langle e_{2}, \ldots, e_{n+1}\right\rangle$. For $k \in\left[r_{1}+1, r_{1}+n\right]$, we have $\operatorname{dim} V_{k}=\operatorname{dim}\left(V_{\mathrm{RP}_{\underline{w}}(k)} \cap\left\langle e_{2}, \ldots, e_{n+1}\right\rangle\right)=\operatorname{dim} V_{\mathrm{RP}_{\underline{w}}(k)}-1=\left(i_{k}+1\right)-1=$ $i_{k}$, where the second equality holds since $\mathrm{RP}_{\underline{w}}(k) \leq r_{1}$ for $k \in\left[r_{1}+1, r_{1}+n\right]$ thus $\left\langle e_{1}\right\rangle \subset V_{\mathrm{RP}_{\underline{w}}(k)}$. For $k \in\left[r_{1}+n+1, r\right]$, we have $\operatorname{dim} V_{k}=\operatorname{dim} H_{k-n}=j_{k-n}=i_{k}$.

We now check, for $k \in[1, r]$, the inclusions $V_{\mathrm{LP}_{\underline{w}}(k)} \subset V_{k} \subset V_{\mathrm{RP}_{\underline{w}}(k)}$. We start with the inclusions $V_{\mathrm{LP}_{\underline{w}}(k)} \subset V_{k}$. For $k \in\left[1, r_{1}\right]$, we have $\mathrm{LP}_{\underline{w}}(k)=\overline{\mathrm{L}}_{\underline{w^{\prime}}}(k) \leq r_{1}$ thus $V_{\mathrm{LP}_{\underline{w}}(k)}=V_{\mathrm{LP}_{w^{\prime}}(k)}=H_{\mathrm{LP}_{w^{\prime}}(k)}+\left\langle e_{1}\right\rangle \subset H_{k}+\left\langle e_{1}\right\rangle=V_{k}$. For $k \in\left[r_{1}+1, r_{1}+\right.$ $n]$ and $\mathrm{RP}_{\underline{w}}(k)>\mathrm{RP}_{\underline{w}}\left(\mathrm{LP}_{\underline{w}} \overline{(k)}\right)$, we have $\mathrm{RP}_{\underline{w}}\left(\mathrm{LP}_{\underline{w}}(k)\right)=\mathrm{LP}_{\underline{w}}\left(\mathrm{RP}_{\underline{w}}(k)\right)$ and $\mathrm{LP}_{\underline{w}}(k) \in\left[r_{1}+1, r_{1}+n\right] \cup\{-\infty\}$. If $\mathrm{LP}_{\underline{w}}(k)=-\infty$, then $V_{\mathrm{LP}_{\underline{w}}(k)}=0$ and the inclusion holds. Otherwise, we have $V_{\mathrm{LP}_{w}(k)}=V_{\mathrm{RP}_{w}\left(\mathrm{LP}_{w}(k)\right)} \cap\left\langle e_{2}, \ldots, e_{n+1}\right\rangle=$ $V_{\mathrm{LP}_{\underline{w}}\left(\mathrm{RP}_{\underline{w}}(k)\right)} \cap\left\langle e_{2}, \ldots, e_{n+1}\right\rangle \subset V_{\mathrm{RP}_{\underline{w}}(k)} \cap\left\langle e_{2}, \ldots, e_{n+1}\right\rangle=\bar{V}_{k} . \quad$ For $\quad k \in\left[r_{1}+1\right.$, $\left.r_{1}+n\right]$ and $\operatorname{RP}_{\underline{w}}(k)<\operatorname{RP}_{\underline{w}}\left(\mathrm{LP}_{\underline{w}}(k)\right)$, we have $\operatorname{RP}_{\underline{w}}\left(\operatorname{RP}_{\underline{w}}\left(\operatorname{LP}_{\underline{w}}(k)\right)\right)=\operatorname{RP}_{\underline{w}}(k) \leq$ $r_{1}$. We have $V_{\mathrm{LP}_{\underline{w}}(k)}=V_{\operatorname{RP}_{\underline{w}}\left(\mathrm{LP}_{\underline{w}}(k)\right)} \cap\left\langle e_{2}, \ldots, e_{n+1}\right\rangle \subset V_{\operatorname{RP}_{\underline{w}}\left(\operatorname{RP}_{\underline{w}}\left(\mathrm{LP}_{\underline{w}}(k)\right)\right)} \cap\left\langle e_{2}, \ldots\right.$, $\left.e_{n+1}\right\rangle \subset V_{\operatorname{RP}_{\underline{w}}(k)} \cap\left\langle e_{2}, \ldots, e_{n+1}\right\rangle=V_{K}$. For $k \geq r_{1}+n$ and $\overline{\mathrm{LP}_{\underline{v}}}\left(k-n-r_{1}\right) \neq$ $-\infty$, we have $\mathrm{LP}_{\underline{w}}(k)-n=\mathrm{LP}_{\underline{w^{\prime}}}(k-n) \geq r_{1}+n+1$ thus $V_{\mathrm{LP}_{\underline{w}}(k)}=H_{\mathrm{LP}_{\underline{w}}(k)-n}=$ $H_{\mathrm{LP}_{w^{\prime}}(k-n)} \subset H_{k-n}=V_{k}$. For $k \geq r_{1}+n$ and $\mathrm{LP}_{\underline{v}}\left(k-n-r_{1}\right)=-\infty$, we have $\operatorname{RP}_{\underline{w}}\left(\mathrm{LP}_{\underline{w}}(k)\right)=\mathrm{LP}_{\underline{w}^{\prime}}(k-n) \leq r_{1} \quad$ and $\quad \mathrm{LP}_{\underline{w}}(k) \in\left[r_{1}+1, r_{1}+n\right] \cup\{-\infty\}$. If $\mathrm{LP}_{\underline{w}}^{-}(k)=-\infty$, then $V_{\mathrm{LP}_{\underline{w}}(k)}=0$ and the inclusion holds. Otherwise, we have $V_{\mathrm{LP}_{\underline{w}}(k)}=V_{\mathrm{RP}_{\underline{w}}\left(\mathrm{LP}_{\underline{w}}(k)\right)} \cap\left\langle e_{2}, \ldots, e_{n+1}\right\rangle=V_{\mathrm{LP}_{\underline{w^{\prime}}}(k-n)} \cap\left\langle e_{2}, \ldots, e_{n+1}\right\rangle=\left(H_{\mathrm{LP}_{\underline{w^{\prime}}}(k-n)}\right.$ $\left.+\left\langle e_{1}\right\rangle\right) \cap\left\langle e_{2}, \ldots, e_{n+1}\right\rangle=H_{\mathrm{LP}_{\underline{w}^{\prime}}(k-n)} \subset H_{k-n}=V_{k}$.

We finish with the inclusions $V_{k} \subset V_{\mathrm{RP}_{\underline{w}}(k)}$. For $k \in\left[1, r_{1}\right]$, we have $\mathrm{RP}_{\underline{w}}(k)=$ $\mathrm{RP}_{\underline{w^{\prime}}}(k) \leq r_{1}$ thus $V_{k}=H_{k}+\left\langle e_{1}\right\rangle \subset H_{\mathrm{RP}_{\underline{w}^{\prime}}(k)}+\left\langle e_{1}\right\rangle=V_{\mathrm{RP}_{\underline{w^{\prime}}}(k)}=V_{\mathrm{RP}_{\underline{w}}(k)}$. For $k \in$ $\left[r_{1}+1, r_{1}+n\right]$, we have $V_{k}=V_{\mathrm{RP}_{w}(k)} \cap\left\langle e_{2}, \ldots, e_{n+1}\right\rangle \subset V_{\mathrm{RP}_{w}(k)}$. For $k \geq r_{1}+n$ and $\mathrm{RP}_{\underline{v}}\left(k-n-r_{1}\right) \neq-\infty$, we have $\mathrm{RP}_{\underline{w}}(k)-n=\mathrm{RP}_{\underline{w^{\prime}}}(k-n) \geq r_{1}+n+1$ thus $V_{k}=H_{k-n} \subset H_{\mathrm{RP}_{\underline{w^{\prime}}}(k-n)}=H_{\mathrm{RP}_{\underline{w}}(k)-n}=V_{\mathrm{RP}_{\underline{w}}(k)}$. For $k \geq r_{1}+n$ and $\mathrm{RP}_{\underline{v}}(k-$ $\left.n-r_{1}\right)=-\infty$, we have $\operatorname{RP}_{\underline{w}}\left(\operatorname{RP}_{\underline{w}}(\bar{k})\right)=\operatorname{RP}_{\underline{w^{\prime}}}(\bar{k}-n) \leq r_{1}$ and $\mathrm{RP}_{\underline{w}}(k) \in\left[r_{1}+\right.$ $\left.1, r_{1}+n\right] \cup\{-\infty\}$. If $\operatorname{RP}_{\underline{w}}(k)=-\infty$, then $V_{\mathrm{LP}_{\underline{w}}(k)}=\left\langle e_{1}, \ldots, e_{n+1}\right\rangle$ and the inclusion holds. Otherwise, we have $V_{k}=H_{k-n} \subset H_{\mathrm{RP}_{\underline{w}^{\prime}}(k-n)}=V_{\mathrm{RP}_{\underline{w}^{\prime}}(k-n)}=$ $V_{\mathrm{RP}_{\underline{w}}\left(\operatorname{RP}_{\underline{w}}(k)\right)}$ But since $V_{k}=H_{k-n} \subset\left\langle e_{2}, \ldots, e_{n+1}\right\rangle$, we get $V_{k} \subset V_{\mathrm{RP}_{\underline{w}}\left(\operatorname{RP}_{\underline{w}}(k)\right)} \cap$ $\left\langle e_{2}, \ldots, e_{n+1}\right\rangle=V_{\mathrm{RP}_{\underline{w}}(k)}$ where the last equality holds since $\mathrm{RP}_{\underline{w}}(k) \in\left[r_{1}+1, r_{1}+\right.$ $n]$.

Now we prove that $f$ and $g$ are inverse to each other and that $\pi_{\underline{w}}=\pi_{\underline{w^{\prime}}} \circ c^{-1} \circ f$. We first prove that $g \circ f$ is the identity. Write $g \circ f\left(\left(V_{k}\right)_{k \in[1, r]}\right)=\left(V_{k}^{\prime}\right)_{k \in[1, r]}$ and $f\left(\left(V_{k}\right)_{k \in[1, r]}\right)=\left(H_{a}\right)_{a \in\left[1, r_{1}+r_{2}\right]}$. For $k \in\left[1, r_{1}\right]$, we have $V_{k}^{\prime}=H_{k}+\left\langle e_{1}\right\rangle=$ $\left(V_{k} \cap\left\langle e_{2}, \ldots, e_{n+1}\right\rangle\right)+\left\langle e_{1}\right\rangle$. But for such $k$, we have $\left\langle e_{1}\right\rangle \subset V_{k}$, this implies $\left(V_{k} \cap\left\langle e_{2}, \ldots, e_{n+1}\right\rangle\right)+\left\langle e_{1}\right\rangle=V_{k}$. For $k \in\left[r_{1}+1, r_{1}+n\right]$, we proceed by induction on $k$ and remark that $\mathrm{RP}_{\underline{w}}(k)<k$. We have $V_{k}^{\prime}=V_{\mathrm{RP}_{\underline{w}}(k)}^{\prime} \cap\left\langle e_{2}, \ldots, e_{n+1}\right\rangle=$ $V_{\mathrm{RP}_{w}(k)} \cap\left\langle e_{2}, \ldots, e_{n+1}\right\rangle=V_{k}$. For $k \geq r_{1}+n+1$, we have $V_{k}^{\prime}=H_{k-n}=V_{k}$.

Next we prove that $f \circ g$ is the identity. Write $f \circ g\left(\left(H_{a}\right)_{a \in\left[1, r_{1}+r_{2}\right]}\right)=$ $\left(H_{a}^{\prime}\right)_{a \in\left[1, r_{1}+r_{2}\right]}$ and $g\left(\left(H_{a}\right)_{a \in\left[1, r_{1}+r_{2}\right]}\right)=\left(V_{k}\right)_{k \in[1, r]}$. For $a \in\left[1, r_{1}\right]$, we have $H_{a}^{\prime}=$ $V_{a} \cap\left\langle e_{2}, \ldots, e_{n+1}\right\rangle=\left(H_{a}+\left\langle e_{1}\right\rangle\right) \cap\left\langle e_{2}, \ldots, e_{n+1}\right\rangle$. But for such $a$, we have $H_{a} \subset$ $\left\langle e_{2}, \ldots, e_{n+1}\right\rangle$ and this implies $\left(H_{a}+\left\langle e_{1}\right\rangle\right) \cap\left\langle e_{2}, \ldots, e_{n+1}\right\rangle=H_{a}$. For $a \in\left[r_{1}+\right.$ $\left.1, r_{1}+r_{2}\right]$, we have $H_{a}^{\prime}=V_{a+n}=H_{a}$.

Finally we check that $\pi_{\underline{w}}=\pi_{\underline{w}^{\prime}} \circ c^{-1} \circ f$. Write $f\left(\left(V_{k}\right)_{k \in[1, r]}\right)=\left(H_{a}\right)_{a \in\left[r_{1}+r_{2}\right]}$, $\pi_{\underline{w}}\left(\left(V_{k}\right)_{k \in[1, r]}\right)=\left(U_{\ngtr}\right)_{\ngtr \in[1, n]}$ and $\pi_{\underline{w}^{\prime}}\left(\left(H_{a}\right)_{a \in\left[1, r_{1}+r_{2}\right]}\right)=\left(U_{\ngtr}^{\prime}\right)_{\ngtr \in[1, n]}$. We need to prove that $U_{\ngtr}=U_{\nexists}^{\prime}$ for all $\not \equiv \in[1, n]$. If $\mathrm{LO}_{\underline{v}}(\nsupseteq) \neq-\infty$, we have $\mathrm{LO}_{\underline{w}^{\prime}}(\ngtr)=$ $\mathrm{LO}_{\underline{w}}(\nsupseteq)-n \geq r_{1}+1$. We get $U_{\ngtr}^{\prime}=H_{\mathrm{LO}_{\underline{w}^{\prime}}(\ngtr)}=V_{\mathrm{LO}_{\underline{w}^{\prime}}(\nsupseteq)+n}=V_{\mathrm{LO}_{\underline{w}}(\ngtr)}=U_{\ngtr}$. If
$\mathrm{LO}_{\underline{v}}(\nsupseteq)=-\infty$, we have $\mathrm{LO}_{\underline{w}^{\prime}}(\nsupseteq)=\mathrm{LO}_{\underline{u}}(\ngtr+1)=\mathrm{RP}_{\underline{w}}\left(r_{1}+\nsupseteq\right) \leq r_{1}$ and $\mathrm{LO}_{\underline{w}}(\nsupseteq)=$ $r_{1}+\not \not \neq$. We get $U_{\ngtr}^{\prime}=H_{\mathrm{LO}_{\underline{w}^{\prime}}(\ngtr)}=H_{\mathrm{RP}_{\underline{w}}\left(r_{1}+\ngtr\right)}^{-}=V_{\mathrm{RP}_{\underline{w}}\left(r_{1}+\ngtr\right)} \cap\left\langle e_{2} ; \ldots, e_{n+1}\right\rangle=$ $V_{r_{1}+\neq}=V_{\mathrm{LO}_{\underline{\underline{w}}}(\ngtr)}=U_{\ngtr}$.

### 4.3 A Product Formula in Cobordism

As a consequence of Theorem 4.11 we prove a product formula in the algebraic cobordism $\Omega^{*}\left(F l_{n+1}\right)$.

Corollary 4.12 Let $\underline{w}$ be a reduced word and $w \in W$ the associated element.

1. If $w \nsupseteq c$, then $\left[X_{\underline{w}}\right] \cdot\left[\mathbf{F}_{n}\right]=0$ in $\Omega^{*}\left(F l_{n+1}\right)$.
2. If $w \geq c$, there exist $\underline{u} \in \underline{W}_{n+1}^{1}$ and $\underline{v} \in \underline{W}_{n+1}^{n}$ such that $\underline{w}=\underline{u} \underline{c} \underline{v}$ modulo commuting relations and we have

$$
\left[X_{\underline{w}}\right] \cdot\left[\mathbf{F}_{n}\right]=\left[X_{c^{-1}(\underline{u}) \underline{v}}\right] .
$$

$$
\text { in } \Omega^{*}\left(F l_{n+1}\right)
$$

Proof The product $\left[X_{\underline{w}}\right] \cdot\left[\mathbf{F}_{n}\right]$ is given by pulling back the exterior product $X_{\underline{w}} \times$ $\mathbf{F}_{n} \rightarrow F l_{n+1} \times F l_{n+1}$ along the diagonal map $\Delta: F l_{n+1} \rightarrow F l_{n+1} \times F l_{n+1}$, see [16, Remark 4.1.14]. We thus have $\left[X_{\underline{w}}\right] \cdot\left[\mathbf{F}_{n}\right]=\Delta^{*}\left[X_{\underline{w}} \times \mathbf{F}_{n} \rightarrow F l_{n+1} \times F l_{n+1}\right]$. Applying [16, Corollary 6.5.5.1], we get $\Delta^{*}\left[X_{\underline{w}} \times \mathbf{F}_{n} \rightarrow F l_{n+1} \times F l_{n+1}\right]=$ $\left[X_{\underline{w}} \times_{F l_{n+1}} \mathbf{F}_{n}\right]$ in $\Omega^{*}(X)$. The result follows since $\left[c\left(X_{c^{-1}(\underline{u}) \underline{v}}\right)\right]=\left[X_{c^{-1}(\underline{u}) \underline{v}}\right]$.

As a special case, we recover the restriction formula in Theorem 3.3 as a product formula.

Corollary 4.13 Let $\underline{w}=\underline{c} \underline{v}$ be a reduced word with $\underline{v} \in \underline{W}_{n+1}^{n}$. Then we have the following formula in $\Omega^{*}\left(F l_{n+1}\right)$ :

$$
\left[X_{\underline{w}}\right] \cdot\left[\mathbf{F}_{n}\right]=\left[X_{\underline{v}}\right] .
$$

By reversing the order of the simple reflection $s_{1}, \ldots, s_{n}$ (or equivalently by conjugating with the element $w_{0}$ ) we also obtain the following results in $\Omega^{*}\left(F l_{n+1}\right)$ :

Proposition 4.14 Let $\underline{w} \in \underline{W}_{n+1}$ be a reduced word and $\underline{c}^{\prime}:=s_{n} \ldots s_{1}$ a Coxeter element. Let $\mathbf{F}_{n}^{\prime}=\left\{U_{\bullet} \in F l_{n+1} \mid U_{1}=\left\langle e_{n+1}\right\rangle\right\}$. Then $w \nsupseteq c^{\prime}$ is equivalent to $X_{w} \cap$ $\mathbf{F}_{n}^{\prime}=\emptyset$ and we have the following alternatives:
(1) If $w \nsupseteq c^{\prime}$, then $\left[X_{\underline{w}}\right] \cdot\left[\mathbf{F}_{n}^{\prime}\right]=0$ in $\Omega^{*}\left(F l_{n+1}\right)$.
(2) If $w \geq c^{\prime}$, then, modulo commuting relations, we have $\underline{w}=\underline{u} \underline{c^{\prime}} \underline{v}$ with $\underline{u} \in \underline{W}_{n+1}^{1}$ and $\underline{v} \in \underline{W}_{n+1}^{n}$. Furthermore, we have

$$
\left[X_{\underline{w}}\right] \cdot\left[\mathbf{F}_{n}^{\prime}\right]=\left[X_{c^{\prime-1}(\underline{u}) \underline{v}}\right]
$$

in $\Omega^{*}\left(F l_{n+1}\right)$.
In particular, if $\underline{w}=\underline{c}^{\prime} \underline{v}$ is reduced with $\underline{v} \in W_{n+1}^{1}$, then we have the following formula in $\Omega^{*}\left(F l_{n+1}\right)$ :

$$
\left[X_{\underline{w}}\right] \cdot\left[\mathbf{F}_{n}^{\prime}\right]=\left[X_{\underline{v}}\right] .
$$

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# Residue Mirror Symmetry for Grassmannians 

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#### Abstract

Motivated by recent works on localizations in A-twisted gauged linear sigma models, we discuss a generalization of toric residue mirror symmetry to complete intersections in Grassmannians.


Keywords A-twisted gauged linear sigma models • Quasimaps • Toric residue mirror symmetry

## 1 Introduction

A-twisted gauged linear sigma models are 2-dimensional topological field theories introduced by Witten [69]. An A-twisted gauged linear sigma model is specified by a reductive algebraic group $G$ (or its compact real form) called the gauge group, an affine space $W$ with a linear action of $G \times \mathbb{G}_{m}$ called the matter, and an element $\xi$ of the dual $\mathfrak{z}^{*}$ of the center of the Lie algebra of $G$ called the Fayet-Iliopoulos param-

[^31]eter. The weights of the $\mathbb{G}_{m}$-actions are called $R$-charges. One can also introduce a superpotential in the theory, which is a $G$-invariant function on $W$ of R-charge 2. The correlation functions, which are quantities of primary interest, do not depend on the potential.

An A-twisted gauged linear sigma model with a suitable Fayet-Iliopoulos parameter is expected to be equivalent to the topological sigma model whose target is the classical vacuum subspace of the symplectic reduction $W / /{ }_{\xi} G$. This comes from a stronger expectation that the low-energy limit of a gauged linear sigma model should give the non-linear sigma model whose target is the symplectic reduction $W / /{ }_{\xi} G$.

A prototypical example is the case $G=\mathbb{G}_{m}$ and $W=\mathbb{A}^{6}$, with the action

$$
\begin{equation*}
G \times \mathbb{G}_{m} \ni(\alpha, \beta):\left(z_{1}, \ldots, z_{5}, P\right) \mapsto\left(\alpha z_{1}, \ldots, \alpha z_{5}, \alpha^{-5} \beta^{2} P\right) \tag{1}
\end{equation*}
$$

and a potential $P f$, which is the product of the variable $P$ and a homogeneous polynomial $f$ in $z_{1}, \ldots, z_{5}$ of degree 5 . The symplectic reduction $W / / \xi G$ for the positive $\xi$ gives the total space of the bundle $\mathcal{O}_{\mathbb{P}^{4}}(-5)$. The R-charge of the $P$ field indicates that the target space should be considered not as a manifold but as a supermanifold, where the parity of the fiber is odd.

One candidate for a mathematical theory of A-twisted gauged linear sigma models is symplectic vortex invariants [19, 20, 35, 59, 60, 74] and their generalizations incorporating potentials [28,65]. Another candidate is quasimap theory, which is an intersection theory on moduli spaces of maps to the quotient stacks [ $W / G$ ]. A review of the latter theory, with historical remarks and extensive references, can be found in [21]. These two approaches should be related by Hitchin-Kobayashi correspondence for vortices [13, 58, 68].

When the gauge group is abelian, quasimap theory as a mathematical theory of Atwisted gauged linear sigma models goes back to [52]. The relation with the Yukawa coupling of the mirror is formulated as toric residue mirror conjecture in $[3,4]$ and proved in [11, 46, 63, 64].

Quasimap theory in the special case of projective hypersurfaces is also studied in the insightful paper [34], where a heuristic relation with semi-infinite homologies of loop spaces is discussed. This eventually leads to Givental's proof [31] of classical mirror symmetry [16] for the quintic 3-fold. This has been extended to toric complete intersections in [32].

The correlation functions of A-twisted gauged linear sigma models in the cases when gauge groups are not necessarily abelian are computed in [9, 24] using supersymmetric localization of path integrals. The result is given in terms of JeffreyKirwan residues, and reproduces the results of [52] in abelian cases.

The aim of this paper is twofold. One is to give an expository account of quasimap theory and its relation to other subjects such as instantons and integrable systems. The other is to formulate Conjecture 5 , which states that the correlation function defined in (5) in terms of residues coincides with the generating function of quasimap invariants defined in (190), and prove it for Grassmannians in Sect. 12. This can be considered as a generalization of toric residue mirror symmetry to Grassmannians. We also show
in Sect. 8 that a slightly weakened version of toric residue mirror conjecture follows from Givental's mirror theorem. Nothing else in this paper is new.

This paper is organized as follows: In Sect. 2, we recall the description of correlation functions of A-twisted gauged linear sigma models given in [9, 24]. In Sect. 3, we recall the definition of the quasimap spaces $\mathbf{Q}\left(\mathbb{P}^{n-1} ; d\right)$. They are compactifications of the spaces of holomorphic maps of degrees $d$ from $\mathbb{P}^{1}$ to $\mathbb{P}^{n-1}$, and play an essential role in Givental's homological geometry [33, 34]. In Sect. 4, we recall toric residue mirror symmetry for Calabi-Yau complete intersections in projective spaces. In Sect.5, we discuss quasimap invariants of concave bundles. In Sect. 6, we recall classical mirror symmetry for toric hypersurfaces proved in [32]. The exposition in Sect. 6 follows [43] closely. In Sect. 7, we briefly recall the definition of quasimap spaces for toric varieties due to [52]. In Sect. 8, we show that a slightly weakened version of toric residue mirror conjecture for CY hypersurfaces follows from classical mirror symmetry. In Sect. 9, we recall a theorem of Martin which relates integration on a symplectic quotient by a compact Lie group to that on the quotient by a maximal torus. In Sect. 10, we recall the definition of quasimap spaces to GIT quotients, which are called quasimap graph spaces in [22]. The quasimap spaces come with the universal $G$-bundle and the canonical virtual fundamental classes, which allow us to define numerical invariants. We formulate Conjecture 5, which states that correlation functions of A-twisted gauged linear sigma models given in (5) are generating functions of quasimap invariants. There is a natural $\mathbb{G}_{m}$-action on the quasimap graph space coming from the $\mathbb{G}_{m}$-action on the domain curve. There is a distinguished connected component of the fixed locus of this action, which is used to define the $I$-function. In Sect. 11, the quasimap spaces and the $I$-functions for Grassmannians are recalled from [10]. In Sect. 12, we prove Conjecture 5 for Grassmannians. For this purpose, we introduce abelianized quasimap spaces for Grassmannians, which allows us to relate quasimap invariants for Grassmannians with correlations function in (5). In Sect. 13, we discuss the relation between gauged linear sigma models and Bethe ansatz following [56]. In Sects. 14-16, we recall the relations of quasimaps with instantons, monopoles, and vortices respectively.

## 2 Correlation Functions of A-Twisted Gauged Linear Sigma Models

## 2.1

Let $G$ be a reductive algebraic group of rank $r$, and $W$ be a representation of $G \times \mathbb{G}_{m}$. The center of $G$ and its Lie algebra will be denoted by $Z(G)$ and $\mathfrak{z}$. Fix a maximal torus $T$ of $G$, and let $\mathfrak{t}$ be its Lie algebra. The set of roots, its subset of positive roots, and the Weyl group will be denoted by $\Delta, \Delta_{+}$, and $\mathscr{W}:=N(T) / T$. Let $W=\bigoplus_{i=1}^{N} W_{i}$ be the weight space decomposition of $W$ with respect to the action of $T \times \mathbb{G}_{m}$. The weight of $W_{i}$ will be denoted by $\left(\rho_{i}, r_{i}\right) \in \mathfrak{t}^{\vee} \oplus \mathbb{Z}$, and $r_{i}$ will be
called the $R$-charge. If $W$ admits an action of another torus $H$ commuting with the action of $G \times \mathbb{G}_{m}$, then one can introduce the twisted mass $\lambda \in \mathfrak{h}$ in the theory, which corresponds to the equivariant parameter for the $H$-action. The $T \times \mathbb{G}_{m} \times H$ weight of $W_{i}$ will be denoted by $\left(\rho_{i}, r_{i}, \nu_{i}\right) \in \mathfrak{t}^{\vee} \oplus \mathbb{Z} \oplus \mathfrak{h}^{\vee}$. We also introduce the complexified Fayet-Illiopoulos parameter $t^{\prime} \in \mathfrak{z}^{\vee} \otimes_{\mathbb{R}} \mathbb{C}$, which corresponds to the complexified Kähler form of the symplectic quotient. Here, we save the unprimed symbol $t$ for the indeterminate in the generating function of quasimap invariants (see (22), (60), (85) and (231)).

For $d \in \mathfrak{t}$ and $t^{\prime} \in \mathfrak{z}^{\vee}$, the composition of the surjection $\mathfrak{t}^{\vee} \rightarrow \mathfrak{z}^{\vee}$ dual to the inclu$\operatorname{sion} \mathfrak{z} \hookrightarrow \mathfrak{t}$ and the evaluation $\mathfrak{t}^{\vee} \times \mathfrak{t} \rightarrow \mathbb{C}$ will be denoted by $t^{\prime} \cdot d$ or $t^{\prime}(d)$.

## 2.2

For $d \in \mathfrak{t}$ and $x \in \mathfrak{t}$, let

$$
\begin{equation*}
Z_{d}(x):=Z_{d}^{\mathrm{vec}}(x) Z_{d}^{\mathrm{mat}}(x) \tag{2}
\end{equation*}
$$

be the product of

$$
\begin{equation*}
Z_{d}^{\mathrm{vec}}(x):=\prod_{\alpha \in \Delta_{+}}(-1)^{\alpha(d)+1} \alpha^{2}(x) \tag{3}
\end{equation*}
$$

and

$$
\begin{equation*}
Z_{d}^{\mathrm{mat}}(x):=\prod_{i=1}^{N}\left(\rho_{i}(x)+\nu_{i}(\lambda)\right)^{r_{i}-\rho_{i}(d)-1} \tag{4}
\end{equation*}
$$

Here the superscripts 'vec' and 'mat' stands for the vector multiplet and the matter chiral multiplet respectively. According to [9, 24], the correlation function of a $\mathscr{W}-$ invariant polynomial $P(x) \in \mathbb{C}[\mathfrak{t}]^{\mathscr{W}}$ on a 2 -sphere is given, up to sign introduced by hand, by

$$
\begin{equation*}
\langle P(x)\rangle_{\mathrm{GLSM}}=\frac{1}{|\mathscr{W}|} \sum_{d \in \mathrm{P}^{\vee}} e^{t^{\prime} \cdot d} \mathrm{JK}_{\mathfrak{c}}\left(Z_{d}(x) P(x)\right) \tag{5}
\end{equation*}
$$

Here $\mathrm{P}^{\vee}$ is the coweight lattice of $G$ and $\mathrm{JK}_{c}$ is the Jeffrey-Kirwan residue defined in [63, Sect. 2] (cf. also [15]). The cone $\mathfrak{c} \subset \mathfrak{z}^{\vee}$ is the ample cone of the GIT quotient determined by the Fayet-Iliopoulos parameter $\eta$.

## 2.3

One can introduce a variable $\mathbf{z}$ associated with the background value of an auxiliary gauge field in the gravity multiplet. This corresponds to the equivariant parameter for the $\mathbb{G}_{m}$-action on the domain curve. This turns (2) into the product

$$
\begin{equation*}
Z_{d}(x ; \mathbf{z}):=Z_{d}^{\mathrm{vec}}(x ; \mathbf{z}) Z_{d}^{\mathrm{mat}}(x ; \mathbf{z}) \tag{6}
\end{equation*}
$$

of

$$
Z_{d}^{\mathrm{vec}}(x ; \mathbf{z}):=\prod_{\alpha \in \Delta_{+}}(-1)^{\alpha(d)+1} \alpha(x) \alpha(x+d \mathbf{z})
$$

and

$$
\begin{equation*}
Z_{d}^{\mathrm{mat}}(x ; \mathbf{z}):=\prod_{i=1}^{N} \frac{\prod_{l=-\infty}^{-1}\left(\rho_{i}(x)+\nu_{i}(\lambda)-\left(l+\frac{r_{i}}{2}\right) \mathbf{z}\right)}{\prod_{l=-\infty}^{\rho_{i}(d)-r_{i}}\left(\rho_{i}(x)+\nu_{i}(\lambda)-\left(l+\frac{r_{i}}{2}\right) \mathrm{z}\right)}, \tag{8}
\end{equation*}
$$

and the correlation function of $P(x) \in \mathbb{C}[\mathfrak{t}]^{\mathscr{W}}$ is given by

$$
\begin{equation*}
\langle P(x)\rangle_{\mathrm{GLSM}}^{H \times \mathbb{G}_{m}}=\frac{1}{|\mathscr{W}|} \sum_{d \in \mathrm{P}^{\vee}} e^{t^{\prime} \cdot d} \mathrm{JK}_{\mathfrak{c}}\left(Z_{d}(x ; \mathbf{z}) P(x)\right) . \tag{9}
\end{equation*}
$$

## 2.4

Another quantity of interest is the effective twisted superpotential on the Coulomb branch, or the effective potential for short. It is defined as the sum

$$
\begin{equation*}
W_{\mathrm{eff}}\left(x ; t^{\prime}\right):=W_{\mathrm{FI}}\left(x ; t^{\prime}\right)+W_{\mathrm{vec}}(x)+W_{\mathrm{mat}}(x) \tag{10}
\end{equation*}
$$

of the Fayet-Illiopoulos term

$$
\begin{equation*}
W_{\mathrm{FI}}\left(x ; t^{\prime}\right):=t^{\prime} \cdot x, \tag{11}
\end{equation*}
$$

the vector multiplet term

$$
\begin{equation*}
W_{\mathrm{vec}}(x):=-\pi \sqrt{-1} \sum_{\alpha \in \Delta^{+}} \alpha(x), \tag{12}
\end{equation*}
$$

and the matter term

$$
\begin{equation*}
W_{\mathrm{mat}}(x):=-\sum_{i=1}^{N}\left(\rho_{i}(x)+\nu_{i}(\lambda)\right)\left(\log \left(\rho_{i}(x)+\nu_{i}(\lambda)\right)-1\right) . \tag{13}
\end{equation*}
$$

## 3 Quasimap Spaces for Projective Spaces

## 3.1

A holomorphic map $u: \mathbb{P}^{1} \rightarrow \mathbb{P}^{n-1}$ of degree $d$ is given by a collection $\left(u_{i}\left(z_{1}, z_{2}\right)\right)_{i=1}^{n}$ of $n$ homogeneous polynomials of degree $d$ satisfying the following condition:

$$
\begin{equation*}
\text { There exists no }\left(z_{1}, z_{2}\right) \in \mathbb{A}^{2} \backslash\{0\} \text { such that } u\left(z_{1}, z_{2}\right)=0 \in \mathbb{A}^{n} . \tag{14}
\end{equation*}
$$

Two collections $\left(u_{i}\left(z_{1}, z_{2}\right)\right)_{i=1}^{n}$ and $\left(u_{i}^{\prime}\left(z_{1}, z_{2}\right)\right)_{i=1}^{n}$ define the same map if and only if there exists $\alpha \in \mathbb{G}_{m}$ such that $u_{i}\left(z_{1}, z_{2}\right)=\alpha u_{i}^{\prime}\left(z_{1}, z_{2}\right)$ for all $i \in\{1, \ldots, n\}$. It follows that the space

$$
\begin{equation*}
\mathcal{M}\left(\mathbb{P}^{n-1} ; d\right):=\left\{u: \mathbb{P}^{1} \rightarrow \mathbb{P}^{n-1} \mid \operatorname{deg} u=d\right\} \tag{15}
\end{equation*}
$$

of holomorphic maps of degree $d$ from $\mathbb{P}^{1}$ to $\mathbb{P}^{n-1}$ can be compactified to the projective space of dimension $n(d+1)-1$, whose homogeneous coordinate is given by the coefficients $\left(a_{i j}\right)_{i, j}$ of the collection $\left(u_{i}\left(z_{1}, z_{2}\right)\right)_{i=1}^{n}$ of homogeneous polynomials of degree $d$;

$$
\begin{equation*}
u_{i}\left(z_{1}, z_{2}\right)=\sum_{j=0}^{d} a_{i j} z_{1}^{j} z_{2}^{d-j}, \quad i=1, \ldots, n \tag{16}
\end{equation*}
$$

This compactification is called the quasimap space and denoted by $\mathbf{Q}\left(\mathbb{P}^{n-1} ; d\right)$. An element of the quasimap space is called a quasimap.

## 3.2

A point $\left[z_{1}: z_{2}\right] \in \mathbb{P}^{1}$ is a base point (or singularity) of a quasimap $u$ if $u\left(z_{1}, z_{2}\right)=0$. A quasimap is a genuine map outside of the base locus. If the degree of the base locus is $d^{\prime}$, then a quasimap can be considered as a genuine map of degree $d-d^{\prime}$. However, it is more convenient to think of a quasimap as a morphism to the quotient stack $\left[\mathbb{A}^{n} / \mathbb{G}_{m}\right]$. By definition, a morphism from $\mathbb{P}^{1}$ to $\left[\mathbb{A}^{n} / \mathbb{G}_{m}\right]$ is a principal $\mathbb{G}_{m^{-}}$ bundle $P$ over $\mathbb{P}^{1}$ and a $\mathbb{G}_{m}$-equivariant morphism $\tilde{u}: P \rightarrow \mathbb{A}^{n}$. It is a quasimap if the generic point of $P$ is mapped to the semi-stable locus $\mathbb{A}^{n} \backslash\{0\}$.

## 3.3

Let $x \in H^{2}\left(\mathbf{Q}\left(\mathbb{P}^{n-1} ; d\right) ; \mathbb{Z}\right)$ be the ample generator of the cohomology ring of $\mathbf{Q}\left(\mathbb{P}^{n-1} ; d\right) \cong \mathbb{P}^{n(d+1)-1}$, so that

$$
\begin{equation*}
H^{*}\left(\mathbf{Q}\left(\mathbb{P}^{n-1} ; d\right) ; \mathbb{Z}\right) \cong \mathbb{Z}[x] /\left(x^{n(d+1)}\right) \tag{17}
\end{equation*}
$$

Given a polynomial $P(x) \in \mathbb{C}[x]$, we set

$$
\begin{equation*}
\langle P(x)\rangle_{\mathbb{P}^{n-1}}:=\sum_{d=0}^{\infty} q^{d}\langle P(x)\rangle_{\mathbb{P}^{n-1}, d} \in \mathbb{C} \llbracket q \rrbracket, \tag{18}
\end{equation*}
$$

where

$$
\begin{equation*}
\langle P(x)\rangle_{\mathbb{P}^{n-1}, d}:=\int_{\mathbf{Q}_{\left(\mathbb{P}^{n-1} ; d\right)}} P(x) \tag{19}
\end{equation*}
$$

is the integration over the quasimap space. It follows from

$$
\left\langle x^{k}\right\rangle_{\mathbb{P}^{n-1}, d}= \begin{cases}1 & k=n(d+1)-1  \tag{20}\\ 0 & \text { otherwise }\end{cases}
$$

that

$$
\left\langle x^{k}\right\rangle_{\mathbb{P}^{n-1}}= \begin{cases}q^{d} & k=n(d+1)-1 \text { for some } d \in \mathbb{Z}^{\geq 0}  \tag{21}\\ 0 & \text { otherwise }\end{cases}
$$

## 3.4

If we set $G:=\mathbb{G}_{m}$ and $W:=\mathbb{C}^{n}$ with the action $G \times \mathbb{G}_{m} \ni(\alpha, \beta):\left(w_{1}, \ldots, w_{n}\right) \mapsto$ $\left(\alpha w_{1}, \ldots, \alpha w_{n}\right)$, then we have $Z_{d}^{\text {vec }}(x)=1$ and $Z_{d}^{\text {mat }}(x)=\left(x^{-d-1}\right)^{n}$, so that (5) gives the same result as (21) under the identification

$$
\begin{equation*}
q=e^{t^{\prime}} \tag{22}
\end{equation*}
$$

## 3.5

The small quantum cohomology of $\mathbb{P}^{n-1}$ is the free $\mathbb{C} \llbracket q \rrbracket$-module

$$
\begin{equation*}
\mathrm{QH}\left(\mathbb{P}^{n-1}\right):=H^{*}\left(\mathbb{P}^{n-1} ; \mathbb{C} \llbracket q \rrbracket\right) \tag{23}
\end{equation*}
$$

equipped with multiplication given by

$$
\begin{equation*}
x^{i} \circ x^{j}:=\sum_{k=0}^{n} \sum_{d=0}^{\infty} q^{d}\left\langle I_{0,3, d}\right\rangle\left(x^{i}, x^{j}, x^{k}\right) x^{n-k-1} \tag{24}
\end{equation*}
$$

Here

$$
\begin{equation*}
\left\langle I_{0,3, d}\right\rangle(a, b, c):=\int_{\left[\overline{\mathcal{M}}_{0,3}\left(\mathbb{P}^{n-1} ; d\right)\right]^{\text {virt }}} \operatorname{ev}_{1}^{*} a \cup \operatorname{ev}_{2}^{*} b \cup \operatorname{ev}_{3}^{*} c \tag{25}
\end{equation*}
$$

is the 3-point Gromov-Witten invariant. It is an associative commutative deformation of the classical cohomology ring;

$$
\begin{equation*}
\mathrm{QH}\left(\mathbb{P}^{n-1}\right) /(q) \cong H^{*}\left(\mathbb{P}^{n-1} ; \mathbb{C}\right) \tag{26}
\end{equation*}
$$

Since the virtual dimension of the moduli space of stable maps is given by

$$
\begin{equation*}
\text { virt. } \operatorname{dim} \overline{\mathcal{M}}_{g, k}(X ; d)=(1-g)(\operatorname{dim} X-3)+\left\langle c_{1}(X), d\right\rangle+k \tag{27}
\end{equation*}
$$

in general, one has

$$
\begin{equation*}
\text { virt.dim} \overline{\mathcal{M}}_{0,3}\left(\mathbb{P}^{n-1} ; d\right)=n d+n-1 \tag{28}
\end{equation*}
$$

The 3-point Gromov-Witten invariant in (24) is non-zero only if

$$
\begin{equation*}
\text { virt. } \operatorname{dim} \overline{\mathcal{M}}_{0,3}\left(\mathbb{P}^{n-1} ; d\right)=i+j+k \tag{29}
\end{equation*}
$$

Since $0 \leq i, j, k \leq n-1$, one has (29) only if $d=0, i+j+k=n-1$ or $d=1$, $i+j+k=2 n-1$. This shows that $x^{i} \circ x^{j}=x^{i+j}$ for $i+j \leq n-1$. Since there is a unique line passing through two points on $\mathbb{P}^{n-1}$ in general position, and this line intersects a hyperplane at one point, one has $x \circ x^{n-1}=q$. Hence the ring structure of the quantum cohomology of $\mathbb{P}^{n-1}$ is given by

$$
\begin{equation*}
\mathrm{QH}\left(\mathbb{P}^{n-1}\right) \cong(\mathbb{C} \llbracket q \rrbracket)[x] /\left(x^{n}-q\right) \tag{30}
\end{equation*}
$$

We write the ring homomorphism $\mathbb{C}[x] \rightarrow \mathrm{QH}\left(\mathbb{P}^{n-1}\right)$ sending $x$ to $x$ as $P(x) \mapsto$ $\stackrel{\circ}{P}(x)$.

Theorem 1 For any $P(x) \in \mathbb{C}[x]$, one has

$$
\begin{equation*}
\langle P(x)\rangle_{\mathbb{P}^{n-1}}=\int_{\mathbb{P}^{n-1}} \stackrel{\circ}{P}(x) . \tag{31}
\end{equation*}
$$

Proof Since both sides of (31) are linear in $P(x) \in \mathbb{C}[x]$, it suffices to show

$$
\begin{equation*}
\left\langle x^{k}\right\rangle_{\mathbb{P}^{n-1}}=\int_{\mathbb{P}^{n-1}} x^{\circ k} \tag{32}
\end{equation*}
$$

for any $k \in \mathbb{N}$, which is obvious from (21) and (30).
Theorem 1 is equivalent to the Vafa-Intriligator formula [42, 67]:
Corollary 2 (Vafa-Intriligator formula for projective spaces) For any $P(x) \in \mathbb{C}[x]$, one has

$$
\begin{equation*}
\int_{\mathbb{P}^{n-1}} \stackrel{\circ}{P}(x)=\frac{1}{n} \sum_{\lambda^{n}=q} \frac{P(\lambda)}{\lambda^{n-1}}, \tag{33}
\end{equation*}
$$

where the sum is over $\lambda \in \mathbb{C} \llbracket q^{1 / n} \rrbracket$ satisfying $\lambda^{n}=q$.
Proof Since the integration over the projective space can be written by residue as

$$
\begin{equation*}
\int_{\mathbb{P}^{r-1}} x^{k}=\delta_{r-1, k}=\operatorname{Res} \frac{x^{k} d x}{x^{r}}, \tag{34}
\end{equation*}
$$

one has

$$
\begin{align*}
\int_{\mathbb{P}^{n-1}} \stackrel{\circ}{P}(x) & =\langle P(x)\rangle_{\mathbb{P}^{n-1}}  \tag{35}\\
& =\sum_{d=0}^{\infty} q^{d} \int_{\mathbf{Q}\left(\mathbb{P}^{n-1} ; d\right)} P(x)  \tag{36}\\
& =\sum_{d=0}^{\infty} q^{d} \operatorname{Res} \frac{P(x) d x}{x^{n(d+1)}}  \tag{37}\\
& =\operatorname{Res} \frac{x^{-n} P(x)}{1-q x^{-n}}  \tag{38}\\
& =\operatorname{Res} \frac{P(x)}{x^{n}-q}  \tag{39}\\
& =\frac{1}{n} \sum_{\lambda^{n}=q} \frac{P(\lambda)}{\lambda^{n-1}}, \tag{40}
\end{align*}
$$

and (33) is proved.

## 3.6

The projective space $\mathbb{P}^{n-1}$ has a natural action of $\mathrm{GL}_{n}$, which restricts to the action of the diagonal maximal torus $H$. The equivariant cohomology is defined as the ordinary cohomology $H_{H}^{*}\left(\mathbb{P}^{n-1}\right):=H^{*}\left(\mathbb{P}_{H}^{n-1}\right)$ of the Borel construction $\mathbb{P}_{H}^{n-1}:=\mathbb{P}^{n-1} \times_{H}$ $E H$, where $E H$ is the product of $n$ copies of the total space of the tautological bundle $\mathcal{O}_{\mathbb{P}^{\infty}}(-1)$ over $B \mathbb{G}_{m}=\mathbb{P}^{\infty}$. It follows that $\mathbb{P}_{H}^{n-1}$ is the projectivization $\mathbb{P}(\mathcal{E})$ of the vector bundle $\mathcal{E}:=\bigoplus_{i=1}^{n} \pi_{i}^{*} \mathcal{O}_{\mathbb{P} \infty}(-1)$ of rank $n$ over $\left(\mathbb{P}^{\infty}\right)^{n}$. A standard result on the cohomology of a projective bundle (see e.g. [37, p. 606]) shows that $H^{*}\left(\mathbb{P}_{H}^{n-1}\right)$ is generated over $H_{H}^{*}(\mathrm{pt})=H^{*}\left(\left(\mathbb{P}^{\infty}\right)^{n}\right) \cong \mathbb{C}\left[\lambda_{1}, \ldots, \lambda_{n}\right]$ by $x:=-c_{1}\left(\mathcal{O}_{\mathbb{P}(\mathcal{E})}(-1)\right)$ with one relation

$$
\begin{equation*}
(-x)^{n}-c_{1}(\mathcal{E})(-x)^{n-1}+c_{2}(\mathcal{E})(-x)^{n-2}+\cdots+(-1)^{n} c_{n}(\mathcal{E})=0 \tag{41}
\end{equation*}
$$

Since $c_{i}(\mathcal{E})=(-1)^{i} \sigma_{i}\left(\lambda_{1}, \ldots, \lambda_{n}\right)$, one obtains

$$
\begin{equation*}
H_{H}^{*}\left(\mathbb{P}^{n-1}\right) \cong \mathbb{C}\left[x, \lambda_{1}, \ldots, \lambda_{n}\right] / \prod_{i=1}^{n}\left(x-\lambda_{i}\right) \tag{42}
\end{equation*}
$$

The $H$-fixed locus $\left(\mathbb{P}^{n-1}\right)^{H}$ consists of $n$ points $\left\{p_{i}\right\}_{i=1}$, where $p_{i}$ is the point $\left[z_{1}\right.$ : $\left.\cdots: z_{n}\right] \in \mathbb{P}^{n-1}$ with $z_{i}=1$ and $z_{j}=0$ for $i \neq j$. Since the tautological bundle $\mathcal{O}_{\mathbb{P}(\mathcal{E})}(-1)$ restricts to $\pi_{i}^{*} \mathcal{O}_{\mathbb{P} \infty}(-1)$ on $\left(p_{i}\right)_{T}=\left(\mathbb{P}^{\infty}\right)^{n}$, one has

$$
\begin{equation*}
\iota_{i}^{*} x=\lambda_{i} \tag{43}
\end{equation*}
$$

The push-forward

$$
\begin{equation*}
\int_{\mathbb{P}^{n-1}}^{H}: H_{H}^{*}\left(\mathbb{P}^{n-1}\right) \rightarrow H_{H}^{*}(\mathrm{pt}) \cong \mathbb{C}\left[\lambda_{1}, \ldots, \lambda_{n}\right] \tag{44}
\end{equation*}
$$

along the natural map $\left(\mathbb{P}^{n-1}\right)_{H} \rightarrow(\mathrm{pt})_{H} \cong B H$ is called the equivariant integration. The localization theorem [2] shows

$$
\begin{align*}
\int_{\mathbb{P}^{n-1}}^{H} P(x) & =\sum_{i=1}^{n} \frac{\iota_{i}^{*} P(x)}{\operatorname{Eul}^{H}\left(N_{p_{i} / \mathbb{P}^{n-1}}\right)}  \tag{45}\\
& =\sum_{i=1}^{n} \frac{P\left(\lambda_{i}\right)}{\prod_{j \neq i}\left(\lambda_{i}-\lambda_{j}\right)} \\
& =\operatorname{Res} \frac{P(x) d x}{\prod_{i=1}^{n}\left(x-\lambda_{i}\right)}
\end{align*}
$$

for any $P(x) \in H_{H}^{*}\left(\mathbb{P}^{n-1}\right)$.

## 3.7

The quasimap space $\mathbf{Q}\left(\mathbb{P}^{n-1} ; d\right)$ has a natural action of $H \times \mathbb{G}_{m}$ given by

$$
\begin{equation*}
H \times \mathbb{G}_{m} \ni\left(\alpha_{1}, \ldots, \alpha_{n}, \beta\right):\left(u_{i}\left(z_{1}, z_{2}\right)\right)_{i=1}^{n} \mapsto\left(\alpha_{i} u_{i}\left(z_{1}, \beta z_{2}\right)\right)_{i=1}^{n} \tag{46}
\end{equation*}
$$

The equivariant cohomology of $\mathbf{Q}\left(\mathbb{P}^{n-1} ; d\right)$ with respect to this torus action is given by

$$
\begin{equation*}
H_{H \times \mathbb{G}_{m}}^{*}\left(\mathbf{Q}\left(\mathbb{P}^{n-1} ; d\right) ; \mathbb{C}\right) \cong \mathbb{C}\left[x, \lambda_{1}, \ldots, \lambda_{n}, \mathbf{z}\right] /\left(\prod_{i=1}^{n} \prod_{j=0}^{d}\left(x-\lambda_{i}-j \mathbf{z}\right)\right) \tag{47}
\end{equation*}
$$

The $H \times \mathbb{G}_{m}$-equivariant integration

$$
\begin{equation*}
\langle-\rangle_{\mathbb{P}^{n-1}, d}^{H \times \mathbb{G}_{m}}: H_{H \times \mathbb{G}_{m}}^{*}\left(\mathbf{Q}\left(\mathbb{P}^{n-1} ; d\right) ; \mathbb{C}\right) \rightarrow H^{*}\left(B\left(H \times \mathbb{G}_{m}\right) ; \mathbb{C}\right) \tag{48}
\end{equation*}
$$

is given by

$$
\begin{align*}
\langle P(x)\rangle_{\mathbb{P}^{n-1}, d}^{H \times \mathbb{G}_{m}} & =\operatorname{Res} \frac{P(x) d x}{\prod_{i=1}^{n} \prod_{j=0}^{d}\left(x-\lambda_{i}-j z\right)}  \tag{49}\\
& =\sum_{i=1}^{n} \sum_{j=0}^{d} \frac{P\left(\lambda_{i}+j \mathbf{z}\right)}{\prod_{(k, l) \neq(i, j)}\left(\left(\lambda_{i}+j \mathbf{z}\right)-\left(\lambda_{k}+l \mathbf{z}\right)\right)} \tag{50}
\end{align*}
$$

The $H \times \mathbb{G}_{m}$-equivariant correlator is given by

$$
\begin{equation*}
\langle P(x)\rangle_{\mathbb{P}^{n-1}}^{H \times \mathbb{G}_{m}}:=\sum_{d=0}^{\infty} q^{d}\langle P(x)\rangle_{\mathbb{P}^{n-1}, d}^{H \times \mathbb{G}_{m}} \tag{51}
\end{equation*}
$$

The $H$-equivariant correlator $\langle P(x)\rangle_{\mathbb{P}^{n-1}}^{H}$ and the $\mathbb{G}_{m}$-equivariant correlator $\langle P(x)\rangle_{\mathbb{P}_{m-1}}^{\mathbb{G}_{m}}$ are obtained by setting $\mathbf{Z}=0$ and $\boldsymbol{\lambda}=\left(\lambda_{1}, \ldots, \lambda_{n}\right)=0$ respectively.

## 3.8

The fixed point of the $\mathbb{G}_{m}$-action on $\mathbf{Q}\left(\mathbb{P}^{n-1} ; d\right)$ is the disjoint union

$$
\begin{equation*}
\mathbf{Q}\left(\mathbb{P}^{n-1} ; d\right)^{\mathbb{G}_{m}}=\coprod_{i=0}^{d} \mathbf{Q}\left(\mathbb{P}^{n-1} ; d\right)_{i}^{\mathbb{G}_{m}} \tag{5}
\end{equation*}
$$

of $d+1$ connected components
$\mathbf{Q}\left(\mathbb{P}^{n-1} ; d\right)_{i}^{\mathbb{G}_{m}}:=\left\{\left[a_{1} z_{1}^{i} z_{2}^{d-i}, \ldots, a_{n} z_{1}^{i} z_{2}^{d-i}\right] \in \mathbf{Q}\left(\mathbb{P}^{n-1} ; d\right) \mid\left[a_{1}, \ldots, a_{n}\right] \in \mathbb{P}^{n-1}\right\}$.

Each of these connected components is isomorphic to $\mathbb{P}^{n-1}$, and the base locus is $i 0+$ $(d-i) \infty$. The connected component $\mathbf{Q}\left(\mathbb{P}^{n-1} ; d\right)_{0}^{\mathbb{G}_{m}}$ will be denoted by $\mathbf{Q} .\left(\mathbb{P}^{n-1} ; d\right)$. There is a natural map ev: Q. $\left(\mathbb{P}^{n-1} ; d\right) \rightarrow \mathbb{P}^{n-1}$ called the evaluation map, and one has

$$
\begin{equation*}
\mathbf{Q}\left(\mathbb{P}^{n-1} ; d\right)^{\mathbb{G}_{m}} \cong \coprod_{d_{1}+d_{2}=d} \text { Q. }\left(\mathbb{P}^{n-1} ; d_{1}\right) \times_{\mathbb{P}^{n-1}} \mathbf{Q} .\left(\mathbb{P}^{n-1} ; d_{2}\right) \tag{54}
\end{equation*}
$$

The normal bundle of $\mathbf{Q} .\left(\mathbb{P}^{n-1} ; d\right)$ in $\mathbf{Q}\left(\mathbb{P}^{n-1} ; d\right)$ is given by $\mathcal{O}_{\mathbb{P}^{n-1}}(1)^{\oplus n d}$, whose equivariant Euler class is given by

$$
\begin{equation*}
\operatorname{Eul}^{H \times \mathbb{G}_{m}}\left(N_{\mathbf{Q} .\left(\mathbb{P}^{n-1} ; d\right) / \mathbf{Q}\left(\mathbb{P}^{n-1} ; d\right)}\right)=\prod_{i=1}^{n} \prod_{l=1}^{d}\left(x-\lambda_{i}+l \mathbf{z}\right) \tag{55}
\end{equation*}
$$

The equivariant I-function is defined by

$$
\begin{equation*}
I_{\mathbb{P}^{n-1}}^{H}(t ; \mathbf{z}):=e^{t x / z} \sum_{d=0}^{\infty} e^{d t} I_{d}^{H} \tag{56}
\end{equation*}
$$

where

$$
\begin{align*}
I_{d}^{H}(\mathbf{z}) & :=\operatorname{ev}_{*}\left(\frac{1}{\operatorname{Eul}^{H \times \mathbb{G}_{m}}\left(N_{\mathbf{Q} \cdot\left(\mathbb{P}^{n-1} ; d\right) / \mathbf{Q}\left(\mathbb{P}^{n-1} ; d\right)}\right)}\right)  \tag{57}\\
& =\frac{1}{\prod_{i=1}^{n} \prod_{l=1}^{d}\left(x-\lambda_{i}+l \mathbf{z}\right)} \tag{58}
\end{align*}
$$

The non-equivariant $I$-function is defined similarly, and given by setting $\boldsymbol{\lambda}=0$ in (56);

$$
\begin{equation*}
I_{\mathbb{P}^{n-1}}(t ; \mathbf{z}):=e^{t x / z} \sum_{d=0}^{\infty} \frac{e^{d t}}{\prod_{l=1}^{d}(x+l \mathbf{z})^{n}} \tag{59}
\end{equation*}
$$

## 3.9

Let $\left(\mathbb{C} \llbracket e^{t} \rrbracket\right)[t]$ be the polynomial ring in $t$ with the ring $\mathbb{C} \llbracket e^{t} \rrbracket$ of formal power series in $e^{t}$ as a coefficient. The equivariant $I$-function in (56) is an element of $H_{H}^{*}\left(\mathbb{P}^{n-1} ; \mathbb{C}\right) \otimes_{\mathbb{C}}\left(\mathbb{C} \llbracket e^{t} \rrbracket\right)[t]$, and the variable $t$ is related to the variable $q$ appearing in the correlator by

$$
\begin{equation*}
q=e^{t} \tag{60}
\end{equation*}
$$

The equivariant $I$-function can also be considered as a $H_{H}^{*}\left(\mathbb{P}^{n-1} ; \mathbb{C}\right)$-valued analytic function, which is multi-valued as a function of $q$ and single-valued as a function of $t=\log q$.

### 3.10

There is a $\mathbb{G}_{m}$-equivariant evaluation map $\mathrm{ev}_{0}: \mathbf{Q}\left(\mathbb{P}^{n-1} ; d\right) \rightarrow\left[\mathbb{C}^{n} / \mathbb{G}_{m}\right]$ at the point $0 \in \mathbb{P}^{1}$. By abuse of notation, we also let $x$ denote the $\mathbb{G}_{m}$-equivariant Euler class of the line bundle $\operatorname{ev}_{0}^{*}\left(\mathcal{O}_{\left[\mathbb{C}^{n} / \mathbb{G}_{m}\right]}(1)\right)$. Here $\mathcal{O}_{\left[\mathbb{C}^{n} / \mathbb{G}_{m}\right]}(1)$ is the line bundle $\left[\left(\mathbb{C}^{n} \times \mathbb{C}\right) / \mathbb{G}_{m}\right]$ on the quotient stack with weights $((1, \ldots, 1), 1)$.

Let $\iota_{i}: \mathbf{Q}\left(\mathbb{P}^{n-1} ; d\right)_{i}^{\mathbb{G}_{m}} \rightarrow \mathbf{Q}\left(\mathbb{P}^{n-1} ; d\right)$ be the inclusion of the $i$ th connected component (53). Since $\iota_{i}^{*}(x)=x+i \mathbf{Z}$ (under the identification $\left.\mathbf{Q}\left(\mathbb{P}^{n-1} ; d\right)_{i}^{\mathbb{G}_{m}}=\mathbb{P}^{n-1}\right)$ and

$$
\begin{equation*}
\frac{1}{\operatorname{Eul}^{\mathbb{G}_{m}}\left(N_{\mathbf{Q}\left(\mathbb{P}^{n-1} ; d\right)_{i}^{\mathbb{G}}} / \mathbf{Q}\left(\mathbb{P}^{n-1} ; d\right)\right.} ⿵ 冂=I_{i}(\mathbf{z}) \cup I_{d-i}(-\mathrm{z}) \tag{61}
\end{equation*}
$$

localization with respect to the $\mathbb{G}_{m}$-action shows that

$$
\begin{align*}
\sum_{d=0}^{\infty} e^{d \tau}\left\langle\left. e^{(t-\tau) x / \mathbf{z}}\right|_{\mathbb{P}^{n-1}, d} ^{\mathbb{G}_{m}}\right. & =\sum_{d=0}^{\infty} e^{d \tau} \sum_{i=0}^{d} \int_{\mathbf{Q}\left(\mathbb{P}^{n-1} ; d\right)_{i}^{G_{m}}} \frac{\iota_{i}^{*}\left(e^{(t-\tau) x / \mathbf{z}}\right)}{\operatorname{Eul}^{\mathbb{G}_{m}}\left(N_{\mathbf{Q}\left(\mathbb{P}^{n-1} ; d\right)_{i}^{\mathbb{G}_{m}} / \mathbf{Q}\left(\mathbb{P}^{n-1} ; d\right)}\right)} \\
& =\sum_{d=0}^{\infty} e^{d \tau} \sum_{i=0}^{d} \int_{\mathbb{P}^{n-1}} e^{(t-\tau)(x+i \mathbf{z}) / \mathbf{z}} \cup I_{i}(\mathbf{z}) \cup I_{d-i}(-\mathbf{z}) \\
& =\sum_{d=0}^{\infty} \sum_{i=0}^{d} \int_{\mathbb{P}^{n-1}} e^{t x / \mathbf{z}} e^{t i} I_{i}(\mathbf{z}) \cup e^{-\tau x / \mathbf{z}} e^{(d-i) \tau} I_{d-i}(-\mathbf{z}) \\
& =\int_{\mathbb{P}^{n-1}} I_{\mathbb{P}^{n-1}}(t ; \mathbf{z}) \cup I_{\mathbb{P}^{n-1}}(\tau ;-\mathbf{z}) . \tag{62}
\end{align*}
$$

The factorization of the $H \times \mathbb{G}_{m}$-equivariant correlator is proved similarly as

$$
\begin{aligned}
& \sum_{d=0}^{\infty} e^{d \tau}\left\langle\left. e^{(t-\tau) x / z}\right|_{\mathbb{P}^{n-1}, d} ^{H \times \mathbb{G}_{m}}\right. \\
& \quad=\sum_{d=0}^{\infty} \operatorname{Res} \frac{e^{d \tau} e^{(t-\tau) x / z} d x}{\prod_{i=1}^{n} \prod_{l=0}^{d}\left(x-\lambda_{i}-l \mathbf{z}\right)} \\
& \quad=\sum_{d=0}^{\infty} \sum_{m=0}^{d} \sum_{j=1}^{n} \operatorname{Res}_{x=\lambda_{j}+m z} \frac{e^{d \tau} e^{(t-\tau) x / z} d x}{\prod_{i=1}^{n} \prod_{l=0}^{d}\left(x-\lambda_{i}-l \mathbf{z}\right)}
\end{aligned}
$$

$$
\begin{aligned}
& =\sum_{d=0}^{\infty} \sum_{m=0}^{d} \sum_{j=1}^{n} \operatorname{Res}_{x=\lambda_{j}} \frac{e^{d \tau} e^{(t-\tau) x / z} e^{(t-\tau) m} d x}{\prod_{i=1}^{n} \prod_{l=0}^{d}\left(x-\lambda_{i}-(l-m) \mathbf{z}\right)} \\
& =\sum_{d=0}^{\infty} \sum_{m=0}^{d} \sum_{j=1}^{n} \operatorname{Res}_{x=\lambda_{j}} \frac{e^{t x / z} e^{m t}}{\prod_{i=1}^{n} \prod_{l=1}^{m}\left(x-\lambda_{i}+l \mathbf{z}\right)} \frac{e^{-\tau x / z} e^{(d-m) \tau}}{\prod_{i=1}^{n} \prod_{l=1}^{d-m}\left(x-\lambda_{i}-l \mathbf{z}\right)} \frac{d x}{\prod_{i=1}^{n}\left(x-\lambda_{i}\right)} \\
& =\sum_{d=0}^{\infty} \sum_{d^{\prime}=0}^{\infty} \sum_{j=1}^{n} \operatorname{Res}_{x=\lambda_{j}} \frac{e^{t x / z} e^{d t}}{\prod_{i=1}^{n} \prod_{l=1}^{d}\left(x-\lambda_{i}+l \mathbf{z}\right)} \frac{e^{-\tau x / z} e^{d^{\prime} \tau}}{\prod_{i=1}^{n} \prod_{l=1}^{d^{\prime}}\left(x-\lambda_{i}-l \mathbf{z}\right)} \frac{d x}{\prod_{i=1}^{n}\left(x-\lambda_{i}\right)} \\
& =\int_{\mathbb{P}^{n-1}}^{H} I_{\mathbb{P}^{n-1}}^{H}(t ; \mathbf{z}) \cup I_{\mathbb{P}^{n-1}}^{H}(\tau ;-\mathbf{z}) .
\end{aligned}
$$

This can also be regarded as a purely combinatorial proof.

### 3.11

Let ev: $\overline{\mathcal{M}}_{0,1}\left(\mathbb{P}^{n-1} ; d\right) \rightarrow \mathbb{P}^{n-1}$ be the evaluation map from the moduli space of stable maps of genus 0 and degree $d$ with 1 marked point, and $\psi$ be the first Chern class of the line bundle over $\overline{\mathcal{M}}_{0,1}\left(\mathbb{P}^{n-1} ; d\right)$ whose fiber at a stable map $\varphi:(C, x) \rightarrow \mathbb{P}^{n-1}$ is the cotangent line $T_{x}^{*} C$ at the marked point. The equivariant $J$-function [31] is a $H^{*}\left(\mathbb{P}^{n-1} ; \mathbb{C}\right)$-valued hypergeometric series given by

$$
\begin{equation*}
J_{\mathbb{P}^{n-1}}^{H}(t ; \mathbf{z}):=e^{t x / \mathbf{z}} \sum_{d=0}^{\infty} e^{d t} J_{d} \tag{63}
\end{equation*}
$$

where

$$
\begin{equation*}
J_{d}:=\operatorname{ev}_{*}\left(\frac{1}{\mathrm{z}(\mathrm{z}-\psi)}\right) . \tag{64}
\end{equation*}
$$

### 3.12

The graph space is defined by $G\left(\mathbb{P}^{n-1} ; d\right):=\overline{\mathcal{M}}_{0,0}\left(\mathbb{P}^{n-1} \times \mathbb{P}^{1} ;(d, 1)\right)$. The source of any map $\varphi: C \rightarrow \mathbb{P}^{n-1} \times \mathbb{P}^{1}$ in $G\left(\mathbb{P}^{n-1} ; d\right)$ has a distinguished irreducible component $C_{1}$ which maps isomorphically to $\mathbb{P}^{1}$. Let $G\left(\mathbb{P}^{n-1} ; d\right)_{0}$ be the open subspace of $G\left(\mathbb{P}^{n-1} ; d\right)$ consisting of stable maps without irreducible components mapping constantly to $0 \in \mathbb{P}^{1}$. There is a map $\mathrm{ev}_{0}: G\left(\mathbb{P}^{n-1} ; d\right)_{0} \rightarrow \mathbb{P}^{n-1}$ sending $\varphi: C \rightarrow \mathbb{P}^{n-1} \times \mathbb{P}^{1}$ to $\operatorname{pr}_{1} \circ \varphi\left(\left(\operatorname{pr}_{2} \circ \varphi\right)^{-1}(0)\right)$. The fixed locus of the natural $\mathbb{G}_{m}$-action on $G\left(\mathbb{P}^{n-1} ; d\right)_{0}$ can be identified with $\overline{\mathcal{M}}_{0,1}\left(\mathbb{P}^{n-1} ; d\right)$. Since the natural morphism $G\left(\mathbb{P}^{n-1} ; d\right)_{0} \rightarrow \mathbf{Q}\left(\mathbb{P}^{n-1} ; d\right)_{0}$ is a $\mathbb{G}_{m}$-equivariant birational morphism which commutes with the evaluation maps, the push-forwards $J_{d}$ and $I_{d}$ of 1 by ev ${ }_{0}: G\left(\mathbb{P}^{n-1} ; d\right)_{0} \rightarrow \mathbb{P}^{n-1}$ and ev: $\mathbf{Q}\left(\mathbb{P}^{n-1} ; d\right)_{0} \rightarrow \mathbb{P}^{n-1}$ are equal, and hence $I_{\mathbb{P}^{n-1}}(t ; \mathbf{z})=J_{\mathbb{P}^{n-1}}(t ; \mathbf{z})$.

### 3.13

The effective potential (10) is given by

$$
\begin{equation*}
W_{\mathrm{eff}}(x ; t)=t x-\sum_{i=1}^{n}\left(x-\lambda_{i}\right)\left(\log \left(x-\lambda_{i}\right)-1\right) \tag{65}
\end{equation*}
$$

One can easily see

$$
\begin{align*}
& \frac{\partial W_{\mathrm{eff}}}{\partial x}=t-\sum_{i=1}^{n} \log \left(x-\lambda_{i}\right),  \tag{66}\\
& e^{\partial_{x} W_{\mathrm{eff}}}=q \prod_{i=1}^{n}\left(x-\lambda_{i}\right)^{-1} \tag{67}
\end{align*}
$$

so that

$$
\begin{equation*}
\langle P(x)\rangle_{\mathbb{P}^{n-1}}^{H}=\operatorname{Res} \frac{P(x) d x}{\prod_{i=1}^{n}\left(x-\lambda_{i}\right)\left(1-e^{\partial_{x} W_{\mathrm{eff}}}\right)} . \tag{68}
\end{equation*}
$$

Note that the equation

$$
\begin{equation*}
e^{\partial_{x} W_{\mathrm{eff}}}=1 \tag{69}
\end{equation*}
$$

gives the relation

$$
\begin{equation*}
\prod_{i=1}^{n}\left(x-\lambda_{i}\right)=q \tag{70}
\end{equation*}
$$

in the equivariant quantum cohomology of $\mathbb{P}^{n-1}$.

## 4 Projective Complete Intersections

## 4.1

Let $f_{1}\left(w_{1}, \ldots, w_{n}\right), \ldots, f_{r}\left(w_{1}, \ldots, w_{n}\right) \in \mathbb{C}\left[w_{1}, \ldots, w_{n}\right]$ be homogeneous polynomials of degrees $l_{1}, \ldots, l_{r}$ satisfying the Calabi-Yau condition

$$
\begin{equation*}
l_{1}+\cdots+l_{r}=n \tag{71}
\end{equation*}
$$

Assume that $f_{1}, \ldots, f_{r}$ are sufficiently general so that

$$
\begin{equation*}
Y:=\left\{\left[w_{1}, \ldots, w_{n}\right] \in \mathbb{P}^{n-1} \mid f_{1}\left(w_{1}, \ldots, w_{n}\right)=\cdots=f_{r}\left(w_{1}, \ldots, w_{n}\right)=0\right\} \tag{72}
\end{equation*}
$$

is a smooth complete intersection of dimension $n-r-1$, whose Poincaré dual is

$$
\begin{equation*}
v:=\prod_{i=1}^{r}\left(l_{i} x\right) \tag{73}
\end{equation*}
$$

Define the quasimap space $\mathbf{Q}(Y ; d)$ as the subset of $\mathbf{Q}\left(\mathbb{P}^{n-1} ; d\right)$ consisting of $\left[\varphi_{1}\left(z_{1}, z_{2}\right), \ldots, \varphi_{n}\left(z_{1}, z_{2}\right)\right]$ satisfying

$$
\begin{equation*}
f_{i}\left(\varphi_{1}\left(z_{1}, z_{2}\right), \ldots, \varphi_{n}\left(z_{1}, z_{2}\right)\right)=0 \in \mathbb{C}\left[z_{1}, z_{2}\right] \text { for any } i \in\{1, \ldots, r\} \tag{74}
\end{equation*}
$$

Since $f_{i}\left(\varphi_{1}\left(z_{1}, z_{2}\right), \ldots, \varphi_{n}\left(z_{1}, z_{2}\right)\right) \in \mathbb{C}\left[z_{1}, z_{2}\right]$ is a homogeneous polynomial of degree $d l_{i}$ in $z_{1}$ and $z_{2}$, it contains $d l_{i}+1$ terms, each of which is a homogeneous polynomial of degree $l_{i}$ in $\left(a_{k l}\right)_{k, l}$. With this in mind, the Morrison-Plesser class is defined by

$$
\begin{equation*}
\Phi(Y ; d):=\prod_{i=1}^{r}\left(l_{i} x\right)^{l_{i} d} \in H^{*}\left(\mathbf{Q}\left(\mathbb{P}^{n-1} ; d\right) ; \mathbb{Z}\right) \tag{75}
\end{equation*}
$$

so that $\Phi(Y ; d) \cup v$ is the Poincaré dual of $[\mathbf{Q}(Y ; d)]^{\text {virt }} \in H_{*}\left(\mathbf{Q}\left(\mathbb{P}^{n-1} ; d\right) ; \mathbb{Z}\right)$. If we set

$$
\begin{equation*}
\langle P(x)\rangle_{Y, d}:=\int_{\mathbf{Q}\left(\mathbb{P}^{n-1} ; d\right)} P(x) \cup \Phi(Y ; d) \cup v \tag{76}
\end{equation*}
$$

and

$$
\begin{equation*}
\langle P(x)\rangle_{Y}:=\sum_{d=0}^{\infty} q^{d}\langle P(x)\rangle_{Y, d} \tag{77}
\end{equation*}
$$

for $P(x) \in \mathbb{C}[x]$, then we have

$$
\begin{align*}
\left\langle x^{n-r-1}\right\rangle_{Y} & =\sum_{d=0}^{\infty} q^{d} \operatorname{Res} \frac{x^{n-r-1} \Phi(Y, d) v d x}{x^{n(d+1)}}  \tag{78}\\
& =\sum_{d=0}^{\infty} q^{d} \operatorname{Res} \frac{x^{n-r-1} \prod_{i=1}^{r}\left(l_{i} x\right)^{l_{i} d+1} d x}{x^{n(d+1)}}  \tag{79}\\
& =\sum_{d=0}^{\infty} q^{d} \prod_{i=1}^{r}\left(l_{i}\right)^{l_{i} d+1} \tag{80}
\end{align*}
$$

$$
\begin{equation*}
=\frac{\prod_{i=1}^{r} l_{i}}{1-q \prod_{i=1}^{r}\left(l_{i}\right)^{l_{i}}} . \tag{81}
\end{equation*}
$$

## 4.2

The gauged linear sigma model for $Y$ is obtained from the gauged linear sigma model for $\mathbb{P}^{n-1}$ by adding $r$ fields of $G=\mathbb{G}_{m}$-charge $-l_{1}, \ldots,-l_{r}$ and R-charge 2. One has $Z_{d}^{\text {vec }}(x)=1$ and $Z_{d}^{\text {mat }}(x)=\left(x^{-d-1}\right)^{n} \cdot \prod_{i=1}^{r}\left(-l_{i} x\right)^{l_{i} d+1}$ in this case, so that (5) gives

$$
\begin{align*}
\sum_{d=0}^{\infty} e^{t^{\prime} d} \operatorname{Res}\left(x^{-d-1}\right)^{n} \prod_{i=1}^{r}\left(-l_{i} x\right)^{l_{i} d+1} x^{n-r-1} & =\sum_{d=0}^{\infty} e^{t^{\prime} d} \prod_{i=1}^{r}\left(-l_{i}\right)^{l_{i} d+1}  \tag{82}\\
& =\sum_{d=0}^{\infty}(-1)^{\sum_{i=1}^{r} l_{i} d} e^{t^{\prime} d} \prod_{i=1}^{r}\left(l_{i}\right)^{l_{i} d+1}  \tag{83}\\
& =\sum_{d=0}^{\infty}\left((-1)^{n} e^{t^{\prime}}\right)^{d} \prod_{i=1}^{r}\left(l_{i}\right)^{l_{i} d+1} \tag{84}
\end{align*}
$$

which coincides with (78) under the identification

$$
\begin{equation*}
q=(-1)^{n} e^{t^{\prime}} \tag{85}
\end{equation*}
$$

## 4.3

The mirror $\check{Y}$ of $Y$ is a compactification of a complete intersection in $\mathbb{C}^{n}$ defined by

$$
\begin{align*}
\check{f}_{1} & :=1-\left(a_{1} \check{y}_{1}+\cdots+a_{l_{1}} \check{y}_{l_{1}}\right),  \tag{86}\\
\check{f}_{2} & :=1-\left(a_{l_{1}+1} \check{y}_{l_{1}+1}+\cdots+a_{l_{1}+l_{2}} \check{y}_{l_{1}+l_{2}}\right),  \tag{87}\\
& \vdots  \tag{88}\\
\check{f_{r}} & :=1-\left(a_{l_{1}+\cdots+l_{i-1}+1} \check{y}_{l_{1}+\cdots+l_{i-1}+1}+\cdots+a_{n} \check{y}_{n}\right),  \tag{89}\\
\check{f} 0 & :=\check{y}_{1} \cdots \check{y}_{n}-1 . \tag{90}
\end{align*}
$$

The complex structure of $\check{Y}$ depends not on the individual $a_{i}$ but only on $\alpha=a_{1} \cdots a_{n}$. The Yukawa (n-2)-point function is defined by

$$
\begin{equation*}
\mathcal{Y}(\alpha):=\frac{(-1)^{(n-1)(n-2) / 2}}{(2 \pi \sqrt{-1})^{n-1}} \int_{\check{Y}} \Omega \wedge\left(\alpha \frac{\partial}{\partial \alpha}\right)^{n-2} \Omega, \tag{91}
\end{equation*}
$$

where

$$
\begin{equation*}
\Omega:=\operatorname{Res}\left(\frac{d \check{y}_{1} \wedge \cdots \wedge d \check{y}_{n}}{\check{f}_{0} \check{f}_{1} \cdots \check{f}_{r}}\right) \tag{92}
\end{equation*}
$$

is the holomorphic volume form on $\check{Y}$. The computation in [5, Proposition 5.1.2] shows

$$
\begin{equation*}
\mathcal{Y}(\alpha)=\frac{\prod_{i=1}^{r} l_{i}}{1-\alpha \prod_{i=1}^{r}\left(l_{i}\right)^{l_{i}}} \tag{93}
\end{equation*}
$$

which coincides with (81) under the identification $q=\alpha$ of variables;

$$
\begin{equation*}
\mathcal{Y}(\alpha)=\left.\left\langle x^{n-r-1}\right\rangle_{Y}\right|_{q=\alpha} \tag{94}
\end{equation*}
$$

A generalization of (94) to toric complete intersections is toric residue mirror symmetry conjectured in [3, 4] and proved in [11, 46, 63, 64].

## 5 Concave Bundles on Projective Spaces

## 5.1

Let $l_{1}, l_{2}, \cdots, l_{r}$ be positive integers and

$$
\begin{equation*}
Y:=\mathcal{S}^{\operatorname{pec}}{\underset{\mathbb{P}^{n-1}}{ }}\left(\mathcal{S}_{2} m^{*} \mathcal{E}^{\vee}\right) \tag{95}
\end{equation*}
$$

be the total space of the vector bundle associated with the locally free sheaf

$$
\begin{equation*}
\mathcal{E}:=\mathcal{O}_{\mathbb{P}^{n-1}}\left(-l_{1}\right) \oplus \cdots \oplus \mathcal{O}_{\mathbb{P}^{n-1}}\left(-l_{r}\right) \tag{96}
\end{equation*}
$$

on $\mathbb{P}^{n-1}$. Since any holomorphic map from $\mathbb{P}^{1}$ to $Y$ of positive degree $d$ factors through the zero-section $\mathbb{P}^{n-1} \rightarrow Y$, we define the quasimap space to $Y$ as

$$
\begin{equation*}
\mathbf{Q}(Y ; d):=\mathbf{Q}\left(\mathbb{P}^{n-1} ; d\right) \tag{97}
\end{equation*}
$$

## 5.2

To equip $\mathbf{Q}(Y ; d)$ with a natural obstruction theory, we identify $\mathbf{Q}(Y ; d)$ with an open substack of the mapping stack $\operatorname{Map}\left(\mathbb{P}^{1}, \mathcal{Y}\right)$ to the quotient stack

$$
\begin{equation*}
\mathcal{Y}:=\left[\left(\mathbb{A}^{n} \times \mathbb{A}^{r}\right) / \mathbb{G}_{m}\right] \tag{98}
\end{equation*}
$$

of $\mathbb{A}^{n} \times \mathbb{A}^{r}$ by the $\mathbb{G}_{m}$-action given by

$$
\begin{equation*}
\mathbb{G}_{m} \ni \alpha:\left(x_{1}, \ldots, x_{n}, z_{1}, \ldots, z_{r}\right) \mapsto\left(\alpha x_{1}, \ldots, \alpha x_{n}, \alpha^{-l_{1}} z_{1}, \ldots, \alpha^{-l_{r}} z_{r}\right) \tag{99}
\end{equation*}
$$

A morphism $\mathbb{P}^{1} \rightarrow \mathcal{Y}$ consists of a line bundle $\mathcal{L}$ on $\mathbb{P}^{1}$ and sections

$$
\begin{equation*}
\left(\left(\varphi_{i}\right)_{i=1}^{n},\left(\psi_{j}\right)_{j=1}^{r}\right) \in\left(\left(H^{0}(\mathcal{L})\right)^{n} \times \prod_{j=1}^{r} H^{0}\left(\mathcal{L}^{\otimes\left(-l_{j}\right)}\right)\right) \tag{100}
\end{equation*}
$$

whose degree is defined as the degree of $\mathcal{L}$.

## 5.3

Recall from [8, Definition 4.4] that an obstruction theory for a Deligne-Mumford stack $\mathcal{X}$ is a morphism $\phi: E \rightarrow L_{\mathcal{X}}$ from an object $E$ of the derived category of quasicoherent sheaves on $\mathcal{X}$ satisfying
(1) $h^{i}(E) \cong 0$ for $i>0$, and
(2) $h^{i}(E)$ is coherent for $i=0,-1$
to the cotangent complex $L_{\mathcal{X}}$ such that
(1) $h^{0}(\phi)$ is an isomorphism, and
(2) $h^{-1}(\phi)$ is an epimorphism.

It is said to be perfect if $E$ is locally isomorphic to a two-term complex of locally free sheaves of finite rank [8, Definition 5.1].

## 5.4

A perfect obstruction theory produces the virtual fundamental cycle $[\mathcal{X}]^{\text {virt }}$ in the Chow group $A_{\text {virt. } \operatorname{dim} \mathcal{X}}(\mathcal{X})$ of degree

$$
\begin{equation*}
\text { virt. } \operatorname{dim} \mathcal{X}=\operatorname{rank} h^{0}(E)-\operatorname{rank} h^{-1}(E) \tag{101}
\end{equation*}
$$

When $\mathcal{X}$ is a smooth scheme, then the cotangent complex $L_{\mathcal{X}}$ is isomorphic to the sheaf $\Omega_{\mathcal{X}}$ of Kähler differentials, and the virtual fundamental cycle is the Euler class of $h^{-1}(E)$.

## 5.5

The derived mapping stack $\mathbb{R} \operatorname{Map}(\mathcal{S}, \mathcal{T})$ from a proper scheme $\mathcal{S}$ to a derived Artin stack $\mathcal{T}$ is a derived Artin stack (see e.g. [66, Corollary 3.3]) whose tangent complex is given by

$$
\begin{equation*}
T_{\mathbb{R} M a p(\mathcal{S}, \mathcal{T})} \cong \mathbb{R} \pi_{*}\left(\mathbb{L e v}{ }^{*} T_{\mathcal{T}}\right) \tag{102}
\end{equation*}
$$

where

$$
\begin{equation*}
\pi: \mathbb{R} \operatorname{Map}(\mathcal{S}, \mathcal{T}) \times \mathcal{S} \rightarrow \mathbb{R} \operatorname{Map}(\mathcal{S}, \mathcal{T}) \tag{103}
\end{equation*}
$$

is the first projection and

$$
\begin{equation*}
\mathrm{ev}: \mathbb{R} \operatorname{Map}(\mathcal{S}, \mathcal{T}) \times \mathcal{S} \rightarrow \mathcal{T} \tag{104}
\end{equation*}
$$

is the evaluation morphism. It is a derived thickening of the mapping $\operatorname{stack} \operatorname{Map}(\mathcal{S}, \mathcal{T})$, and the pull-back

$$
\begin{equation*}
j^{*}: j^{*} L_{\mathbb{R} \operatorname{Map}(\mathcal{S}, \mathcal{T})} \rightarrow L_{\operatorname{Map}(\mathcal{S}, \mathcal{T})} \tag{105}
\end{equation*}
$$

by the structure morphism

$$
\begin{equation*}
j: \operatorname{Map}(\mathcal{S}, \mathcal{T}) \rightarrow \mathbb{R} \operatorname{Map}(\mathcal{S}, \mathcal{T}) \tag{106}
\end{equation*}
$$

gives an obstruction theory on $\operatorname{Map}(\mathcal{S}, \mathcal{T})$.

## 5.6

The restriction of the natural obstruction theory for $\operatorname{Map}\left(\mathbb{P}^{1}, \mathcal{Y}\right)$ to the open substack $\mathbf{Q}(Y ; d)$ gives an obstruction theory for $\mathbf{Q}(Y ; d)$ with $E=\left.j^{*} L_{\mathbb{R M a p}\left(\mathbb{P}^{1}, \mathcal{Y}\right)}\right|_{\mathbf{Q}(Y ; d)}$ and $\phi=\left.j^{*}\right|_{\mathbf{Q}(Y ; d)}$.

## 5.7

Since $\operatorname{Pic}\left(\mathbb{A}^{n} \times \mathbb{A}^{r}\right)$ is trivial, the Picard group $\operatorname{Pic} \mathcal{Y} \cong \operatorname{Pic} \mathbb{G}_{m}\left(\mathbb{A}^{n} \times \mathbb{A}^{r}\right)$ can be identified with the group of characters of $\mathbb{G}_{m}$, which is non-canonically isomorphic to $\mathbb{Z}$. We fix an isomorphism in such a way that $\bigoplus_{a=0}^{\infty} H^{0}\left(\mathcal{O}_{\mathcal{Y}}(a)\right)$ is the coordinate ring of $\mathbb{A}^{n}$, where $\mathcal{O}_{\mathcal{Y}}(a)$ is the line bundle associated with $a \in \mathbb{Z} \cong$ Pic $\mathcal{Y}$. Since $\mathcal{Y}$ is the quotient stack of $\mathbb{A}^{n} \times \mathbb{A}^{r}$ by the action of $\mathbb{G}_{m}$, the tangent complex $T_{\mathcal{Y}}$ satisfies

$$
\begin{equation*}
\varpi^{*} T_{\mathcal{Y}} \cong \operatorname{Cone}\left(\mathcal{O}_{\mathbb{A}^{n} \times \mathbb{A}^{r}} \otimes \operatorname{Lie}\left(\mathbb{G}_{m}\right) \rightarrow T_{\mathbb{A}^{n} \times \mathbb{A}^{r}}\right) \tag{107}
\end{equation*}
$$

where $\varpi: \mathbb{A}^{n} \times \mathbb{A}^{r} \rightarrow \mathcal{Y}$ is the quotient morphism. This in turn implies that

$$
\begin{equation*}
T_{\mathcal{Y}} \cong \text { Cone }\left(\mathcal{O}_{\mathcal{Y}} \rightarrow \mathcal{O}_{\mathcal{Y}}(1)^{\oplus n} \oplus \bigoplus_{i=1}^{r} \mathcal{O}_{\mathcal{Y}}\left(-l_{i}\right)\right) \tag{108}
\end{equation*}
$$

## 5.8

We write the restriction of the evaluation morphism $\operatorname{Map}\left(\mathbb{P}^{1}, \mathcal{Y}\right) \times \mathbb{P}^{1} \rightarrow \mathcal{Y}$ to the open substack $\mathbf{Q}(Y ; d) \subset \operatorname{Map}\left(\mathbb{P}^{1}, \mathcal{Y}\right)$ as

$$
\begin{equation*}
\mathrm{ev}: \mathbf{Q}(Y ; d) \times \mathbb{P}^{1} \rightarrow \mathcal{Y} \tag{109}
\end{equation*}
$$

again by abuse of notation. We have

$$
\begin{equation*}
\mathrm{ev}^{*} \mathcal{O}_{\mathcal{Y}}(1) \cong \mathcal{O}_{\mathbf{Q}(Y ; d)}(1) \boxtimes \mathcal{O}_{\mathbb{P}^{1}}(d) \tag{110}
\end{equation*}
$$

essentially by definition, where $\mathcal{O}_{\mathbf{Q}(Y ; d)}(1)$ is the ample generator of the Picard group of $\mathbf{Q}(Y ; d) \cong \mathbb{P}^{n(d+1)-1}$. The dual of the natural obstruction theory is given by

$$
\begin{equation*}
\phi^{\vee}: T_{\mathbf{Q}(Y ; d)} \rightarrow E^{\vee}:=\mathbb{R} \pi_{*} \mathrm{ev}^{*} T_{\mathcal{Y}} \tag{111}
\end{equation*}
$$

Note that (the inverse of) the isomorphism $h^{0}\left(\phi^{\vee}\right)$ from $T_{\mathbf{Q}(Y ; d)}$ to

$$
\begin{gather*}
h^{0}\left(E^{\vee}\right) \cong R^{0} \pi_{*} \operatorname{Cone}\left(\mathcal{O}_{\mathbf{Q}(Y ; d) \times \mathbb{P}^{1}} \rightarrow\left(\mathcal{O}_{\mathbf{Q}(Y ; d)}(1) \boxtimes \mathcal{O}_{\mathbb{P}^{1}}(d)\right)^{\oplus n}\right)  \tag{112}\\
\cong \operatorname{Cone}\left(\mathcal{O}_{\mathbf{Q}(Y ; d)} \rightarrow \mathcal{O}_{\mathbf{Q}(Y ; d)}(1)^{\oplus n(d+1)}\right) \tag{113}
\end{gather*}
$$

gives the Euler sequence

$$
\begin{equation*}
0 \rightarrow \mathcal{O}_{\mathbf{Q}(Y ; d)} \rightarrow \mathcal{O}_{\mathbf{Q}(Y ; d)}(1)^{\oplus n(d+1)} \rightarrow T_{\mathcal{Q}(Y ; d)} \rightarrow 0 \tag{114}
\end{equation*}
$$

on $\mathbf{Q}(Y ; d) \cong \mathbb{P}^{n(d+1)-1}$. One has

$$
\begin{align*}
h^{1}\left(E^{\vee}\right) & \cong R^{1} \pi_{*}\left(\bigoplus_{i=1}^{r} \mathcal{O}_{\mathbf{Q}(Y ; d)}\left(-l_{i}\right) \boxtimes \mathcal{O}_{\mathbb{P}^{1}}\left(-l_{i} d\right)\right)  \tag{115}\\
& \cong \bigoplus_{i=1}^{r} \mathcal{O}_{\mathbf{Q}(Y ; d)}\left(-l_{i}\right) \otimes H^{1}\left(\mathcal{O}_{\mathbb{P}^{1}}\left(-l_{i} d\right)\right)  \tag{116}\\
& \cong \bigoplus_{i=1}^{r} \mathcal{O}_{\mathbf{Q}(Y ; d)}\left(-l_{i}\right)^{\oplus\left(l_{i} d-1\right)} \tag{117}
\end{align*}
$$

and $h^{i}\left(E^{\vee}\right) \cong 0$ for $i \neq 0,1$, so that this obstruction theory is perfect. By [8, Proposition 5.6], the resulting virtual fundamental class is given by

$$
\begin{equation*}
[\mathbf{Q}(Y ; d)]^{\mathrm{virt}}=[\mathbf{Q}(Y ; d)] \cap \operatorname{Eul}\left(h^{1}\left(E^{\vee}\right)\right)=[\mathbf{Q}(Y ; d)] \cap \prod_{i=1}^{r}\left(-l_{i} x\right)^{l_{i} d-1} \tag{118}
\end{equation*}
$$

## 5.9

When the degree is zero, the quasimap space $\mathbf{Q}(Y ; 0)$ is naturally isomorphic to $Y$ equipped with the trivial perfect obstruction theory, so that

$$
\begin{equation*}
[\mathbf{Q}(Y ; 0)]^{\text {virt }}=[Y] . \tag{119}
\end{equation*}
$$

### 5.10

For any $P(x) \in \mathbb{C}[x]$, we define

$$
\begin{equation*}
\langle P(x)\rangle_{Y, d}:=\int_{[\mathbf{Q}(Y ; d)]^{\mathrm{yirt}}} P(x) \tag{120}
\end{equation*}
$$

and

$$
\begin{equation*}
\langle P(x)\rangle_{Y}:=\sum_{d=0}^{\infty} q^{d}\langle P(x)\rangle_{Y, d} . \tag{121}
\end{equation*}
$$

It follows that

$$
\begin{align*}
\langle P(x)\rangle_{Y} & =\sum_{d=0}^{\infty} q^{k} \int_{\mathbb{P}^{n(d+1)-1}} P(x) \prod_{i=1}^{r}\left(-l_{i} x\right)^{l_{i} d-1}  \tag{122}\\
& =\sum_{d=0}^{\infty} q^{k} \operatorname{Res} \frac{P(x) \prod_{i=1}^{r}\left(-l_{i} x\right)^{l_{i} d}}{x^{n(d+1)} \prod_{i=1}^{r}\left(-l_{i} x\right)} \tag{123}
\end{align*}
$$

### 5.11

The gauged linear sigma model for $Y$ is obtained from the gauged linear sigma model for $\mathbb{P}^{n-1}$ by adding $r$ fields of $G=\mathbb{G}_{m}$-charge $-l_{1}, \ldots,-l_{r}$ and R-charge 0 . One has $Z_{d}^{\text {vec }}(x)=1$ and $Z_{d}^{\text {mat }}(x)=\left(x^{-d-1}\right)^{n} \cdot \prod_{i=1}^{r}\left(-l_{i} x\right)^{l_{i} d-1}$ in this case, so that (5) gives

$$
\begin{equation*}
\langle P(x)\rangle_{\mathrm{GLSM}}=\sum_{d=0}^{\infty} e^{t^{\prime} d} \operatorname{Res}\left(x^{-d-1}\right)^{n} \prod_{i=1}^{r}\left(-l_{i} x\right)^{l_{i} d-1} P(x), \tag{124}
\end{equation*}
$$

which coincides with (123) under the identification

$$
\begin{equation*}
q=e^{t^{\prime}} \tag{125}
\end{equation*}
$$

### 5.12

If $\left(l_{1}, \ldots, l_{r}\right)$ satisfies the Calabi-Yau condition

$$
\begin{equation*}
l_{1}+\cdots+l_{r}=n \tag{126}
\end{equation*}
$$

then (123) gives

$$
\left\langle x^{k}\right\rangle_{Y}= \begin{cases}\frac{1}{\left(\prod_{i=1}^{r}\left(-l_{i}\right)\right)\left(1-q \prod_{i=1}^{r}\left(-l_{i}\right)^{l_{i}}\right)} & k=n+r  \tag{127}\\ 0 & \text { otherwise }\end{cases}
$$

which matches the Yukawa coupling of the mirror (see e.g. [47, Example 6.15]).

## 6 Classical Mirror Symmetry for Toric Hypersurfaces

## 6.1

Let $\boldsymbol{N}:=\mathbb{Z}^{n}$ be a free abelian group of rank $n$ and $\boldsymbol{M}:=\check{\boldsymbol{N}}:=\operatorname{Hom}(\boldsymbol{N}, \mathbb{Z})$ be the dual group. Let further $(\Delta, \check{\Delta})$ be a polar dual pair of reflexive polytopes in $\boldsymbol{M}$ and $N$.

## 6.2

Recall that the fan polytope of a fan is defined as the convex hull of primitive generators of one-dimensional cones. Let $(\Sigma, \check{\Sigma})$ be a pair of smooth projective fans whose fan polytopes are $\check{\Delta}$ and $\Delta$. The associated toric varieties will be denoted by $X:=X_{\Sigma}$ and $\check{X}:=X_{\check{\Sigma}}$.

## 6.3

The set of primitive generators of one-dimensional cones of the fan $\Sigma$ will be denoted by

$$
\begin{equation*}
B:=\left\{\boldsymbol{b}_{1}, \ldots, \boldsymbol{b}_{m}\right\} \subset N \tag{128}
\end{equation*}
$$

Assume that $B$ generates $N$. One has the fan sequence

$$
\begin{equation*}
0 \rightarrow L \rightarrow \mathbb{Z}^{m} \xrightarrow{b} N \rightarrow 0 \tag{129}
\end{equation*}
$$

and the divisor sequence

$$
\begin{equation*}
0 \rightarrow \boldsymbol{M} \xrightarrow{\boldsymbol{b}^{\vee}} \mathbb{Z}^{m} \rightarrow \check{\boldsymbol{L}} \rightarrow 0 \tag{130}
\end{equation*}
$$

where $\boldsymbol{b}$ sends the $i$ th coordinate vector $e_{i} \in \mathbb{Z}^{m}$ to $\boldsymbol{b}_{i}$. Recall that

$$
\begin{equation*}
\check{\boldsymbol{L}} \cong \operatorname{Pic}(X) \cong H^{2}(X ; \mathbb{Z}), \quad \operatorname{Eff}(X) \subset \boldsymbol{L} \subset \mathbb{Z}^{m} \tag{131}
\end{equation*}
$$

where $\operatorname{Eff}(X)$ denotes the semigroup of the effective curves (see [3, Sect. 3]). We write the group ring of $\boldsymbol{M}$ as $\mathbb{C}[\boldsymbol{M}]$ and define $\mathbb{T}:=\boldsymbol{N}_{\mathbb{G}_{m}}:=\operatorname{Spec} \mathbb{C}[\boldsymbol{M}]$. We also set $\check{\mathbb{T}}:=\operatorname{Spec} \mathbb{C}[\boldsymbol{N}], \check{\mathbb{L}}:=\operatorname{Spec} \mathbb{C}[\boldsymbol{L}]$, and $\mathbb{L}:=\operatorname{Spec} \mathbb{C}[\check{\boldsymbol{L}}]$. The fan sequence induces the exact sequences

$$
\begin{equation*}
1 \rightarrow \mathbb{L} \xrightarrow{\chi}\left(\mathbb{G}_{m}\right)^{m} \rightarrow \mathbb{T} \rightarrow 1 \tag{132}
\end{equation*}
$$

and

$$
\begin{equation*}
1 \rightarrow \check{\mathbb{T}} \rightarrow\left(\mathbb{G}_{m}\right)^{m} \rightarrow \check{\mathbb{L}} \rightarrow 1 \tag{133}
\end{equation*}
$$

of algebraic tori. We write the $i$ th components of the map $\chi: \mathbb{L} \rightarrow\left(\mathbb{G}_{m}\right)^{m}$ in (132) as $\chi_{i}$, and the affine line $\mathbb{A}^{1}$ equipped with the action of $\mathbb{L}$ through $\chi_{i}$ as $\mathbb{A}_{i}$. Then one has

$$
\begin{equation*}
X \cong\left(\prod_{i=1}^{m} \mathbb{A}_{i}\right) / / \theta \mathbb{L} \tag{134}
\end{equation*}
$$

for a suitable choice of a character $\theta \in \check{\boldsymbol{L}} \cong \operatorname{Hom}\left(\mathbb{L}, \mathbb{G}_{m}\right)$. The right-hand side of (134) denotes the GIT quotient with respect to the linearization determined by $\theta$.

## 6.4

We define a graded ring $S_{\Delta}:=\bigoplus_{k=0}^{\infty} S_{\Delta}^{k}$ by

$$
\begin{equation*}
S_{\Delta}^{k}:=\bigoplus_{\boldsymbol{m} \in \boldsymbol{M} \cap(k \Delta)} \mathbb{C} \cdot y_{0}^{k} y^{\boldsymbol{m}} \tag{135}
\end{equation*}
$$

which is a subalgebra of the semigroup ring

$$
\begin{equation*}
\mathbb{C}[\mathbb{N} \times \boldsymbol{M}]=\mathbb{C}\left[y_{0}, \boldsymbol{y}^{ \pm 1}\right]:=\mathbb{C}\left[y_{0}, y_{1}^{ \pm 1}, \ldots, y_{n}^{ \pm 1}\right] \tag{136}
\end{equation*}
$$

of $\mathbb{N} \times \boldsymbol{M}$. It is the anti-canonical ring of $X$, so that one has $X \cong \operatorname{Proj} S_{\Delta}$ if and only if $X$ is Fano. The ring $S_{\Delta}$ is Cohen-Macaulay with the dualizing module $I_{\Delta}:=$ $\bigoplus_{k=0}^{\infty} I_{\Delta}^{k}$ given by

$$
\begin{equation*}
I_{\Delta}^{k}:=\bigoplus_{m \in M \cap \operatorname{Int}(k \Delta)} \mathbb{C} \cdot y_{0}^{k} y^{\boldsymbol{m}} \tag{137}
\end{equation*}
$$

where $\operatorname{Int}(k \Delta)$ is the interior of $k \Delta$.

## 6.5

For $\boldsymbol{\alpha}=\left(\alpha_{1}, \ldots, \alpha_{m}\right) \in\left(\mathbb{G}_{m}\right)^{m}$ (this $\left(\mathbb{G}_{m}\right)^{m}$ can be naturally considered as the dual torus of the big torus of $X_{\Sigma}$ ), we define an element of the group ring $\mathbb{C}[N]$ by

$$
\begin{equation*}
\check{W}_{\alpha}(\check{\boldsymbol{y}}):=\sum_{i=1}^{m} \alpha_{i} \check{\boldsymbol{y}}^{\boldsymbol{b}_{i}} \in \mathbb{C}[\boldsymbol{N}] \tag{138}
\end{equation*}
$$

An element $\check{f} \in \mathbb{C}[N]$ is said to be $\check{\Delta}$-regular if

$$
\begin{equation*}
\check{F}:=\left(\check{F}_{0}, \check{F}_{1}, \ldots, \check{F}_{n}\right):=\left(\check{y}_{0} \check{f}, \check{y}_{0} \check{y}_{1} \partial_{\check{y}_{1}} \check{f}, \ldots, \check{y}_{0} \check{y}_{n} \partial_{\check{y}_{n}} \check{f}\right) \tag{139}
\end{equation*}
$$

is a regular sequence in $S_{\check{\Delta}}$. We write

$$
\begin{equation*}
\left(\left(\mathbb{G}_{m}\right)^{m}\right)^{\mathrm{reg}}:=\left\{\alpha \in\left(\mathbb{G}_{m}\right)^{m} \mid \check{f}_{\alpha}:=1-\check{W}_{\alpha}(\check{\boldsymbol{y}}) \text { is } \check{\Delta} \text {-regular }\right\} . \tag{140}
\end{equation*}
$$

## 6.6

Let $\widetilde{\varphi}: \widetilde{\mathfrak{Y}} \rightarrow\left(\left(\mathbb{G}_{m}\right)^{m}\right)^{\text {reg }}$ be the second projection from

$$
\begin{equation*}
\widetilde{\tilde{\mathfrak{Y}}}=\left\{(\check{\boldsymbol{y}}, \boldsymbol{\alpha}) \in \check{\mathbb{T}} \times\left(\left(\mathbb{G}_{m}\right)^{m}\right)^{\mathrm{reg}} \mid \check{W}_{\alpha}(\check{\boldsymbol{y}})=1\right\} \tag{141}
\end{equation*}
$$

Assume that $X$ is Fano. Any fiber $\check{Y}_{\boldsymbol{\alpha}}:=\widetilde{\varphi}^{-1}(\boldsymbol{\alpha})$ is an uncompactified mirror of a general anti-canonical hypersurface $Y \subset X$. The closure of $\check{Y}_{\alpha}$ in $\check{X}$ is a smooth anticanonical Calabi-Yau hypersurface, which is the compact mirror of $Y$. The quotient of the family $\widetilde{\varphi}: \widetilde{\mathfrak{Y}} \rightarrow\left(\left(\mathbb{G}_{m}\right)^{m}\right)^{\text {reg }}$ by the free $\check{\mathbb{T}}$-action

$$
\begin{equation*}
\check{\mathbb{T}} \ni \check{\boldsymbol{y}}:\left(\check{\boldsymbol{y}}^{\prime},\left(\alpha_{1}, \ldots, \alpha_{m}\right)\right) \mapsto\left(\check{\boldsymbol{y}}^{-1} \check{\boldsymbol{y}}^{\prime},\left(\check{\boldsymbol{y}}^{\boldsymbol{b}_{1}} \alpha_{1}, \ldots, \check{\boldsymbol{y}}^{\boldsymbol{b}_{m}} \alpha_{m}\right)\right) \tag{142}
\end{equation*}
$$

will be denoted by $\check{\varphi}: \check{\mathfrak{Y}} \rightarrow \check{\mathbb{L}}^{\text {reg }}$, where $\check{\mathfrak{Y}}:=\check{\check{\mathfrak{Y}}} / \check{\mathbb{T}}$ and $\check{\mathbb{L}}^{\text {reg }}:=\left(\left(\mathbb{G}_{m}\right)^{m}\right)^{\mathrm{reg}} / \check{\mathbb{T}}$.

## 6.7

Choose an integral basis $\boldsymbol{p}_{1}, \ldots, \boldsymbol{p}_{r}$ of $\check{\boldsymbol{L}} \cong$ Pic $X$ such that each $\boldsymbol{p}_{i}$ is nef. This gives the corresponding coordinate $\boldsymbol{q}=\left(q_{1}, \ldots, q_{r}\right)$ on $\check{\mathbb{L}}$. Let $\check{U}^{\prime} \subset \check{\mathbb{L}}^{\text {reg }}$ be a neighborhood of $q_{1}=\cdots=q_{r}=0$, and $\check{U}$ be the universal cover of $\check{U}^{\prime}$.

## 6.8

We write the image of the Poincaré residue as

$$
\begin{equation*}
H_{\mathrm{res}}^{n-1}\left(\check{Y}_{\alpha}\right):=\operatorname{Im}\left(\operatorname{Res}: H^{0}\left(\check{X}, \Omega_{\check{X}}^{n}\left(* \check{Y}_{\alpha}\right)\right) \rightarrow H^{n-1}\left(\check{Y}_{\alpha}\right)\right) . \tag{143}
\end{equation*}
$$

Let $H_{\mathrm{B}}$ be the pull-back to $\check{U}$ of the local system $\mathrm{gr}_{n-1}^{W} R^{n-1} \check{\varphi}_{!} \mathbb{C}_{\check{\mathfrak{Y}}}$ on $\check{U}^{\prime}$, and $H_{\mathrm{B}}^{\text {res }}$ be the sub-system with stalks $H_{\text {res }}^{n-1}\left(\check{Y}_{\alpha}\right)$. Here $\operatorname{gr}_{n-1}^{W}$ is the weight $n-1$ piece of Deligne's mixed Hodge structure. The residual B-model VHS $\left(\mathcal{H}_{\mathrm{B}}, \nabla_{\mathrm{B}}, \mathscr{F}_{\mathrm{B}}^{\bullet}, Q_{\mathrm{B}}\right)$ on $\check{U}$ consists of the locally free sheaf $\mathcal{H}_{\mathrm{B}}:=H_{\mathrm{B}}^{\mathrm{res}} \otimes_{\mathbb{C}} \mathcal{O}_{\check{U}}$, the Gauss-Manin connection $\nabla_{B}$, the Hodge filtration $\mathscr{F}_{B}^{\bullet}$, and the polarization $Q_{B}: \mathcal{H}_{\mathrm{B}} \otimes_{\mathcal{O}_{\check{U}}} \mathcal{H}_{\mathrm{B}} \rightarrow \mathcal{O}_{\check{U}}$ given by

$$
\begin{equation*}
Q_{\mathrm{B}}\left(\omega_{1}, \omega_{2}\right):=(-1)^{(n-1)(n-2) / 2} \int_{\check{Y}_{\alpha}} \omega_{1} \cup \omega_{2} . \tag{144}
\end{equation*}
$$

## 6.9

On the A-model side, let

$$
\begin{equation*}
H_{\mathrm{amb}}^{\bullet}(Y ; \mathbb{C}):=\operatorname{Im}\left(\iota^{*}: H^{\bullet}(X ; \mathbb{C}) \rightarrow H^{\bullet}(Y ; \mathbb{C})\right) \tag{145}
\end{equation*}
$$

be the subspace of $H^{\bullet}(Y ; \mathbb{C})$ coming from the cohomology classes of the ambient toric variety, and set

$$
\begin{equation*}
U:=\left\{\boldsymbol{\tau}=\boldsymbol{\beta}+\sqrt{-1} \boldsymbol{\omega} \in H_{\mathrm{amb}}^{2}(Y ; \mathbb{C}) \mid\langle\boldsymbol{\omega}, \boldsymbol{d}\rangle \gg 0 \text { for any non-zero } \boldsymbol{d} \in \operatorname{Eff}(Y)\right\} \tag{146}
\end{equation*}
$$

This open subset $U$ is considered as a neighborhood of the large radius limit point. Let $\left(\tau_{i}\right)_{i=1}^{r}$ be the coordinate on $H_{\mathrm{amb}}^{2}(Y ; \mathbb{C})$ dual to the basis $\left\{\boldsymbol{p}_{i}\right\}_{i=1}^{r}$ so that $\boldsymbol{\tau}=$ $\sum_{i=1}^{r} \tau_{i} \boldsymbol{p}_{i}$.

### 6.10

The ambient A-model VHS $\left(\mathcal{H}_{\mathrm{A}}, \nabla_{\mathrm{A}}, \mathscr{F}_{\mathrm{A}}^{\bullet}, Q_{\mathrm{A}}\right)$ consists ([43, Definition 6.2], cf. also [26, Sect. 8.5]) of the locally free sheaf $\mathcal{H}_{\mathrm{A}}=H_{\mathrm{amb}}^{\bullet}(Y ; \mathbb{C}) \otimes_{\mathbb{C}} \mathcal{O}_{U}$, the connection

$$
\begin{equation*}
\nabla_{\mathrm{A}}=d+\sum_{i=1}^{r}\left(\boldsymbol{p}_{i} \circ_{\tau}\right) d \tau_{i}: \mathcal{H}_{\mathrm{A}} \rightarrow \mathcal{H}_{\mathrm{A}} \otimes \Omega_{U}^{1} \tag{147}
\end{equation*}
$$

the Hodge filtration

$$
\begin{equation*}
\mathscr{F}_{\mathrm{A}}^{p}:=H_{\mathrm{amb}}^{\leq 2(n-1-p)}(Y ; \mathbb{C}) \otimes_{\mathbb{C}} \mathcal{O}_{U} \tag{148}
\end{equation*}
$$

and the pairing

$$
\begin{equation*}
Q_{A}: \mathcal{H}_{\mathrm{A}} \otimes_{\mathcal{O}_{U}} \mathcal{H}_{\mathrm{A}} \rightarrow \mathcal{O}_{U}, \quad(\alpha, \beta) \mapsto(2 \pi \sqrt{-1})^{n-1} \int_{Y}(-1)^{\operatorname{deg} \alpha / 2} \alpha \cup \beta \tag{149}
\end{equation*}
$$

which is $(-1)^{n-1}$-symmetric and $\nabla_{\mathrm{A}}$-flat.

### 6.11

Let $\boldsymbol{u}_{i} \in H_{\text {amb }}^{2}(Y ; \mathbb{Z})$ be the first Chern class of the line bundle on $Y$ corresponding to the one-dimensional cone $\mathbb{R}_{\geq 0} \cdot \boldsymbol{b}_{i} \in \Sigma$ and $\boldsymbol{v}=\boldsymbol{u}_{1}+\cdots+\boldsymbol{u}_{m}$ be the restriction of the anti-canonical class of $X$. Denote $\boldsymbol{t}:=\sum_{i=1}^{r} t_{i} \boldsymbol{p}_{i}$. Givental's $I$-function is defined as the series

$$
\begin{equation*}
I_{Y}(\boldsymbol{t} ; \mathbf{z})=e^{t / \mathbf{z}} \sum_{\boldsymbol{d} \in \operatorname{Eff}(X)} e^{\boldsymbol{d} \cdot \boldsymbol{t}} \frac{\prod_{k=-\infty}^{\langle\boldsymbol{d}, \boldsymbol{v}\rangle}(\boldsymbol{v}+k \mathbf{z}) \prod_{j=1}^{m} \prod_{k=-\infty}^{0}\left(\boldsymbol{u}_{j}+k \mathbf{z}\right)}{\prod_{k=-\infty}^{0}(\boldsymbol{v}+k \mathbf{z}) \prod_{j=1}^{m} \prod_{k=-\infty}^{\left\langle\boldsymbol{d}, \boldsymbol{u}_{j}\right\rangle}\left(\boldsymbol{u}_{j}+k \mathbf{z}\right)} \tag{150}
\end{equation*}
$$

which is a multi-valued map from $\check{U}^{\prime}$ (or a single-valued map from $\check{U}$ ) to the classical cohomology group $H_{\mathrm{amb}}^{\bullet}\left(Y ; \mathbb{C}\left[\mathrm{z}^{-1}\right]\right)$. The $J$-function is defined by

$$
\begin{equation*}
J_{Y}(\boldsymbol{\tau} ; \mathbf{z})=L_{Y}(\boldsymbol{\tau}, \mathbf{z})^{-1}(1) \tag{151}
\end{equation*}
$$

where $L_{Y}(\boldsymbol{\tau}, \mathbf{z})$ is the fundamental solution of the quantum differential equation defined explicitly by using the Gromov-Witten invariants as in [43, Eq. (2.3)] with c set to 1 . If we write

$$
\begin{equation*}
I_{Y}(\boldsymbol{t} ; \mathbf{z})=F(\boldsymbol{t}) \mathbf{1}+\frac{G(\boldsymbol{t})}{\mathbf{z}}+O\left(\mathbf{z}^{-2}\right) \tag{152}
\end{equation*}
$$

then Givental's mirror theorem [32] states that

$$
\begin{equation*}
I_{Y}(\boldsymbol{t} ; \mathbf{z})=F(\boldsymbol{t}) \cdot J_{Y}(\varsigma(\boldsymbol{t}) ; \mathbf{z}) \tag{153}
\end{equation*}
$$

where the mirror map $\varsigma: \check{U} \rightarrow H_{\mathrm{amb}}^{2}(Y ; \mathbb{C})$ is defined by

$$
\begin{equation*}
\varsigma(t)=\iota^{*}\left(\frac{G(t)}{F(t)}\right) . \tag{154}
\end{equation*}
$$

The relation between $\boldsymbol{\tau}=\boldsymbol{\varsigma}(\boldsymbol{t})$ and $\boldsymbol{\sigma}=\boldsymbol{\beta}+\sqrt{-1} \boldsymbol{\omega}$ is given by $\boldsymbol{\tau}=2 \pi \sqrt{-1} \boldsymbol{\sigma}$, so that $\mathfrak{I m}(\boldsymbol{\sigma}) \gg 0$ corresponds to $\exp (\boldsymbol{\tau}) \sim 0$. The functions $F(\boldsymbol{t})$ and $G(\boldsymbol{t})$ satisfy the Picard-Fuchs equations, and give periods for the B-model VHS $\left(\mathcal{H}_{\mathrm{B}}, \nabla^{B}, \mathscr{F}_{B}^{\bullet}, Q_{B}\right)$.

### 6.12

Equation (153) implies the existence of an isomorphism

$$
\begin{equation*}
\operatorname{Mir}_{\mathcal{Y}}: \varsigma^{*}\left(\mathcal{H}_{\mathrm{A}}, \nabla_{\mathrm{A}}, \mathscr{F}_{\mathrm{A}}^{\bullet}, Q_{A}\right) \xrightarrow{\sim}\left(\mathcal{H}_{\mathrm{B}}, \nabla_{\mathrm{B}}, \mathscr{F}_{\mathrm{B}}^{\bullet}, Q_{B}\right) \tag{155}
\end{equation*}
$$

of variations of polarized Hodge structures, which sends $F(\boldsymbol{t}) \mathbf{1}$ on the left-hand side to

$$
\begin{equation*}
\Omega:=\operatorname{Res}\left(\frac{1}{\check{f}_{\alpha}} \frac{d \check{y}_{1}}{\check{y}_{1}} \wedge \cdots \wedge \frac{d \check{y}_{n}}{\check{y}_{n}}\right) \tag{156}
\end{equation*}
$$

on the right-hand side. A stronger statement, which gives an isomorphism of the $\widehat{\Gamma}$-integral structure on the A-side and the natural integral structure on the B-side, is proved in [43, Theorem 6.9].

## 7 Quasimap Correlation Functions for Anti-canonical Hypersurfaces in Toric Varieties

## 7.1

For $\boldsymbol{d} \in \operatorname{Eff}(X)$ and $i \in\{1, \ldots, m\}$, we set

$$
k_{i}:= \begin{cases}\left\langle\boldsymbol{u}_{i}, \boldsymbol{d}\right\rangle & \left\langle\boldsymbol{u}_{i}, \boldsymbol{d}\right\rangle \geq 0  \tag{157}\\ -1 & \left\langle\boldsymbol{u}_{i}, \boldsymbol{d}\right\rangle<0\end{cases}
$$

and define the quasimap space of degree $\boldsymbol{d}$ by

$$
\begin{equation*}
\mathbf{Q}(X ; \boldsymbol{d}):=\left(\prod_{i=1}^{m} \mathbb{A}_{i}^{k_{i}+1}\right) / / \theta \mathbb{L} \tag{158}
\end{equation*}
$$

with (134) in mind. An argument parallel to that in Sect. 3.1 shows that $\mathbf{Q}(X ; \boldsymbol{d})$ is a compactification of the space of holomorphic maps $\mathbb{P}^{1} \rightarrow X$ of degree $\boldsymbol{d}$. Later in Sect. 10.4, we will introduce the moduli spaces $\mathbf{Q}(W / / G ; \boldsymbol{d})$ of degree $\boldsymbol{d}$ quasimaps from $\mathbb{P}^{1}$ to more general GIT quotients $W / / G$, and $\mathbf{Q}(X ; \boldsymbol{d})$ here is the special case for $W=\prod_{i=1}^{m} \mathbb{A}_{i}^{k_{i}+1}$ and $G=\mathbb{L}$. The first Chern class of the line bundle on $\mathbf{Q}(X ; \boldsymbol{d})$ associated with the character $\chi_{i}$ of $\mathbb{L}$ will also be denoted by $\boldsymbol{u}_{i}$ by abuse of notation. The Morrison-Plesser class is defined by

$$
\begin{equation*}
\Phi_{\boldsymbol{d}}:=\left(\boldsymbol{u}_{1}+\cdots+\boldsymbol{u}_{m}\right)^{\left\langle\boldsymbol{u}_{1}+\cdots+\boldsymbol{u}_{m}, \boldsymbol{d}\right\rangle} \prod_{\left\langle\boldsymbol{u}_{i}, \boldsymbol{d}\right\rangle<0} \boldsymbol{u}_{i}^{-\left\langle\boldsymbol{u}_{i}, \boldsymbol{d}\right\rangle-1} \tag{159}
\end{equation*}
$$

Here, the latter part $\prod_{\left\langle\boldsymbol{u}_{i}, \boldsymbol{d}\right\rangle<0} \boldsymbol{u}_{i}^{-\left\langle\boldsymbol{u}_{i}, \boldsymbol{d}\right\rangle-1}$ is the Euler class of $h^{1}\left(E^{\vee}\right)$ where $E$ is the canonical obstruction theory for $\mathbf{Q}(X ; \boldsymbol{d})$ defined in Sect. 10.7. The first part $\left(\boldsymbol{u}_{1}+\cdots+\boldsymbol{u}_{m}\right)^{\left\langle\boldsymbol{u}_{1}+\cdots+\boldsymbol{u}_{m}, \boldsymbol{d}\right\rangle}$ is the Euler class of the vector bundle formed by the obstruction spaces to being a quasimap to an anti-canonical hypersurface $Y \subset X$ for each element in $\mathbf{Q}(X ; \boldsymbol{d})$. For a polynomial $P\left(\alpha_{1}, \ldots, \alpha_{m}\right) \in \mathbb{C}\left[\alpha_{1}, \ldots, \alpha_{m}\right]$, we set

$$
\begin{equation*}
\left\langle P\left(\boldsymbol{u}_{1}, \ldots, \boldsymbol{u}_{m}\right)\right\rangle_{X, Y, \boldsymbol{d}}:=\int_{\mathbf{Q}(X ; \boldsymbol{d})} P\left(\boldsymbol{u}_{1}, \ldots, \boldsymbol{u}_{m}\right) \Phi_{\boldsymbol{d}} \tag{160}
\end{equation*}
$$

and

$$
\begin{equation*}
\left\langle P\left(\boldsymbol{u}_{1}, \ldots, \boldsymbol{u}_{m}\right)\right\rangle_{X, Y}:=\sum_{\boldsymbol{d} \in \operatorname{Eff}(X)} \boldsymbol{\alpha}^{\boldsymbol{d}}\left\langle P\left(\boldsymbol{u}_{1}, \ldots, \boldsymbol{u}_{m}\right)\right\rangle_{X, Y, \boldsymbol{d}} \in \mathbb{Z} \llbracket \alpha^{\boldsymbol{d}}: \boldsymbol{d} \in \operatorname{Eff}(X) \rrbracket \tag{161}
\end{equation*}
$$

where the completion is taken with respect to the ideal generated by $\operatorname{Eff}(X) \backslash\{0\}$. Here $\boldsymbol{\alpha}^{\boldsymbol{d}}$ is defined by (131).

## 8 Toric Residue Mirror Symmetry

## 8.1

Let $\check{G}=\left(\check{G}_{0}, \ldots, \check{G}_{n}\right)$ be a regular sequence in $S_{\check{\Delta}}$. If we set $I_{\check{G}}:=I_{\check{\Delta}} /\left(\check{G}_{0}, \ldots, \check{G}_{n}\right)$ $I_{\check{\Delta}}$, then the graded piece $I_{\breve{G}}^{n+1}$ is one-dimensional and spanned by $J_{\check{G}}:=\operatorname{det}$ $\left(\check{y}_{i} \partial_{\check{y}_{i}} \check{G}_{j}\right)_{i, j=0}^{n}$. The toric residue [25] is the map $\operatorname{Res}_{\check{G}}: I_{\check{\Delta}}^{n+1} \rightarrow \mathbb{C}$ sending $\left(\check{G}_{0}, \ldots\right.$, $\left.\check{G}_{n}\right) I_{\check{\Delta}}$ to zero and $J_{\check{G}}$ to the normalized volume $\operatorname{vol}(\check{\Delta})$, i.e., $n!$ times the Euclidean volume of $\check{\Delta}$. For $\boldsymbol{\alpha} \in \check{\mathbb{L}}^{\text {reg }}$, we define $\check{F}_{\alpha}$ as in (139) and write $\operatorname{Res}_{\check{f}_{\alpha}}:=\operatorname{Res}_{\check{F}_{\alpha}}$. Theorem 3 below is introduced in [3, Conjecture 4.6] and proved in [11, 63].

Theorem 3 For any homogeneous polynomial $P\left(\alpha_{1}, \ldots, \alpha_{m}\right) \in \mathbb{C}\left[\alpha_{1}, \ldots, \alpha_{m}\right]$ of degree $n$, the generating function (161) gives the Laurent expansion of the toric residue

$$
\begin{equation*}
\left\langle P\left(\boldsymbol{u}_{1}, \ldots, \boldsymbol{u}_{m}\right)\right\rangle_{X, Y}=(-1)^{n} \operatorname{Res}_{\check{f}_{\alpha}}\left(\check{y}_{0}^{n+1} P\left(\alpha_{1} \check{\boldsymbol{y}}^{\boldsymbol{b}_{1}}, \ldots, \alpha_{m} \check{\boldsymbol{y}}^{\boldsymbol{b}_{m}}\right)\right) \tag{162}
\end{equation*}
$$

around the large radius limit point associated with the fan $\Sigma$.

Reference [3, Conjecture 4.6] is generalized to toric complete intersections in [4, Conjecture 4.6] and proved in [46, 64].

## 8.2

The family $\varphi: \check{\mathcal{Y}} \rightarrow \check{\mathbb{L}}^{\text {reg }}$ of Calabi-Yau manifolds comes with the holomorphic volume form

$$
\begin{equation*}
\Omega:=\operatorname{Res}\left(\frac{1}{\check{f}_{\alpha}} \frac{d \check{y}_{1}}{\check{y}_{1}} \wedge \cdots \wedge \frac{d \check{y}_{n}}{\check{y}_{n}}\right) \in H^{0}\left(\mathcal{H}_{\mathrm{B}}\right) . \tag{163}
\end{equation*}
$$

For a homogeneous polynomial $Q\left(\alpha_{1}, \ldots, \alpha_{m}\right) \in \mathbb{Q}\left[\alpha_{1}, \ldots, \alpha_{m}\right]$ of degree $n-1$, the $Q$-Yukawa $(n-1)$-point function is defined in [3, Definition 9.1] by

$$
\begin{equation*}
Y_{Q}(\boldsymbol{\alpha}):=(-1)^{(n-1)(n-2) / 2} \frac{1}{(2 \pi \sqrt{-1})^{n-1}} \int_{\check{Y}_{\alpha}} \Omega \wedge Q\left(\alpha_{1} \frac{\partial}{\partial \alpha_{1}}, \ldots, \alpha_{m} \frac{\partial}{\partial \alpha_{m}}\right) \Omega \tag{164}
\end{equation*}
$$

where the differential operators $\alpha_{1} \partial / \partial \alpha_{1}, \ldots, \alpha_{m} \partial / \partial \alpha_{m}$ act on $\mathcal{H}_{\mathrm{B}}$ by the GaussManin connection.

## 8.3

For $Q\left(\alpha_{1}, \ldots, \alpha_{m}\right) \in \mathbb{Q}\left[\alpha_{1}, \ldots, \alpha_{m}\right]$, we set

$$
\begin{equation*}
P\left(\alpha_{1}, \ldots, \alpha_{m}\right):=\left(\alpha_{1}+\cdots+\alpha_{m}\right) Q\left(\alpha_{1}, \ldots, \alpha_{m}\right) \in \mathbb{Q}\left[\alpha_{1}, \ldots, \alpha_{m}\right] . \tag{165}
\end{equation*}
$$

By [3, Theorem 9.7], which is attributed to [50], one has an equality

$$
\begin{equation*}
Y_{Q}(\boldsymbol{\alpha})=(-1)^{n} \operatorname{Res}_{\check{f}_{\alpha}}\left(\check{y}_{0}^{n} P\left(\alpha_{1} \check{\boldsymbol{y}}^{\boldsymbol{b}_{1}}, \ldots, \alpha_{m} \breve{\boldsymbol{y}}^{\boldsymbol{b}_{m}}\right)\right) \tag{166}
\end{equation*}
$$

of the Yukawa $(n-1)$-point function and the toric residue.

## 8.4

Assume that the unstable locus of the $\mathbb{L}$-action on $\mathbb{A}^{m}$ with respect to $\theta$ has codimension strictly greater than 1 . Then one has $H^{2}\left(X_{\Sigma}\right)=\operatorname{Pic}\left(X_{\Sigma}\right)=\operatorname{Pic}^{\mathbb{L}}\left(\mathbb{A}^{m}\right)$ so that the class $\boldsymbol{p}_{i}$ corresponds to a one-dimensional representation $\mathbb{C}_{\boldsymbol{p}_{i}}$ of $\mathbb{L}$. By abuse
of notation, we let $\boldsymbol{p}_{i}$ denote the $\mathbb{G}_{m}$-equivariant Euler class of the pull-back of the line bundle $\left[\mathbb{A}^{m} \times \mathbb{C}_{\boldsymbol{p}_{i}} / \mathbb{G}_{m}\right]$ by the evaluation map $\mathrm{ev}_{0}: X_{\boldsymbol{d}} \rightarrow\left[\mathbb{A}^{m} / \mathbb{G}_{m}\right]$ at $0 \in \mathbb{P}^{1}$. Denote $\boldsymbol{v}:=\sum_{i=1}^{m} \boldsymbol{u}_{i}$.

If we set

$$
\begin{equation*}
\Phi(\boldsymbol{t}, \boldsymbol{\tau} ; \mathbf{z}):=\sum_{\boldsymbol{d} \in \operatorname{Eff}(X)} e^{\boldsymbol{\tau} \cdot \boldsymbol{d}} \int_{\mathbf{Q}(X ; \boldsymbol{d})}^{\mathbb{G}_{m}} e^{(\boldsymbol{t}-\boldsymbol{\tau}) / \mathbf{z}} \Phi_{d} \boldsymbol{v} \tag{167}
\end{equation*}
$$

then for any polynomial $R\left(t_{1}, \ldots, t_{r}\right) \in \mathbb{Q}\left[t_{1}, \ldots, t_{r}\right]$, one has

$$
\begin{equation*}
\left.R\left(\mathrm{z} \frac{\partial}{\partial t_{1}}, \ldots, \mathrm{z} \frac{\partial}{\partial t_{r}}\right) \Phi(\boldsymbol{t}, \boldsymbol{\tau} ; \mathbf{z})\right|_{\tau=t}=\sum_{\boldsymbol{d} \in \operatorname{Eff}(X)} e^{\boldsymbol{t} \cdot \boldsymbol{d}} \int_{\mathbf{Q}(X ; \boldsymbol{d})}^{\mathbb{G}_{m}} R\left(\boldsymbol{p}_{1}, \ldots, \boldsymbol{p}_{r}\right) \Phi_{\boldsymbol{d}} \boldsymbol{v} \tag{168}
\end{equation*}
$$

In addition, one has

$$
\begin{equation*}
\Phi(\boldsymbol{t}, \boldsymbol{\tau} ; \mathbf{z})=\int_{Y} I(\boldsymbol{t} ;-\mathbf{z}) \cup I(\boldsymbol{\tau} ; \mathbf{z}) \tag{169}
\end{equation*}
$$

by [32, Proposition 6.2]. This is the toric hypersurface version of (62). Note that $I(t ; 1)$ is convergent for large enough $-\operatorname{Re} \boldsymbol{t}$ by ratio test on the series (150) without the prefactor. By specializing to $\mathbf{z}=1$ and using the definition of $Q_{\mathrm{A}}$, one obtains

$$
\begin{equation*}
\Phi(\boldsymbol{t}, \boldsymbol{\tau} ; 1)=Q_{\mathrm{A}}(I(\boldsymbol{t} ; 1), I(\boldsymbol{\tau} ; 1)) \tag{170}
\end{equation*}
$$

By combining (170) with (153), one obtains

$$
\begin{equation*}
\Phi(\boldsymbol{t}, \boldsymbol{\tau} ; 1)=Q_{\mathrm{A}}\left(L^{-1}(\boldsymbol{t} ; 1) F(\boldsymbol{t}) \mathbf{1}, L^{-1}(\boldsymbol{\tau} ; 1) F(\boldsymbol{\tau}) \mathbf{1}\right) \tag{171}
\end{equation*}
$$

Since $L$ is the fundamental solution for the flat connection $\nabla_{\mathrm{B}}$, the function $\Phi(\boldsymbol{t}, \boldsymbol{\tau} ; 1)$ is obtained by parallel-transporting $F(\boldsymbol{t}) \mathbf{1} \in\left(H_{\mathrm{B}}\right)_{t}$ and $F(\boldsymbol{\tau}) \mathbf{1} \in\left(H_{\mathrm{B}}\right)_{\tau}$ to the fiber at the same point and taking the pairing $Q_{\mathrm{B}}$ at that point (the result does not depend on the choice of the point since $Q_{\mathrm{B}}$ is $\nabla_{\mathrm{B}}$-parallel). By sending (171) by (155), one obtains

$$
\begin{equation*}
(2 \pi \sqrt{-1})^{n-1} \int_{Y} I(t ;-1) I(\tau ; 1)=(-1)^{(n-1)(n-2) / 2} \int_{\check{Y}} \Omega_{t} \wedge \Omega_{\tau} \tag{172}
\end{equation*}
$$

Assume that $P\left(\alpha_{1}, \ldots, \alpha_{m}\right)=\left(\alpha_{1}+\cdots+\alpha_{m}\right) Q\left(\alpha_{1}, \ldots, \alpha_{m}\right)$ for a polynomial $Q$ and take $\left.R\left(t_{1}, \ldots, t_{r}\right):=Q\left(\sum_{i=1}^{r} a_{i, 1} t_{i}, \ldots, \sum_{i=1}^{r} a_{i, m} t_{i}\right)\right)$ where $a_{i, j}$ are integers uniquely satisfying $\chi_{j}=\sum_{i=1}^{r} a_{i, j} \boldsymbol{p}_{i}$. By differentiating (172) by $R\left(\partial_{t_{1}}, \ldots, \partial_{t_{r}}\right)$ and setting $\tau=\boldsymbol{t}$, we obtain toric residue mirror symmetry for polynomials of the form $P\left(\alpha_{1}, \ldots, \alpha_{m}\right)=\left(\alpha_{1}+\cdots+\alpha_{m}\right) Q\left(\alpha_{1}, \ldots, \alpha_{m}\right)$.

## 9 Martin's Formula

## 9.1

We use the same notations $G, T, \mathscr{W}$, and $\Delta$ for a reductive algebraic group, a maximal torus, the Weyl group, and the set of roots as in Sect. 2. Let $W$ be an affine scheme with $G$-action, and fix a character $\theta$ of $G$. We write the line bundle on $W / / T$ associated with $\alpha \in \Delta$ as $L_{\alpha}$, and set

$$
\begin{equation*}
e:=\prod_{\alpha \in \Delta} c_{1}\left(L_{\alpha}\right) \in H^{2|\Delta|}(W / / T ; \mathbb{Z}) \tag{173}
\end{equation*}
$$

We write the natural projection and inclusion as and say that $\tilde{a} \in H^{*}(W / / T)$ is a lift of $a \in H^{*}(W / / G)$ if $\pi^{*} a=\iota^{*} \tilde{a}$.

Theorem 4 (Martin [48]) If $\tilde{a}$ is a lift of a, then one has

$$
\begin{equation*}
\int_{W / / G} a=\frac{1}{|\mathscr{W}|} \int_{W / T} \tilde{a} \cup e \tag{174}
\end{equation*}
$$

## 9.2

Let $\operatorname{Mat}(r, n) \cong \mathbb{A}^{r \times n}$ be the space of $n \times r$ matrices, which is considered as the space of linear maps from an $r$-dimensional vector space to an $n$-dimensional vector space. It has a natural action of $\mathrm{GL}_{r}$, and the GIT quotient $\operatorname{Gr}(r, n):=\operatorname{Mat}(r, n) / / \mathrm{GL}_{r}$ is the Grassmannian of $r$-spaces in an $n$-space.

## 9.3

When $W=\operatorname{Mat}(r, n)$ and $G=\mathrm{GL}_{r}$, one has

$$
\begin{align*}
W / / G & \cong \operatorname{Gr}(r, n),  \tag{175}\\
W / / T & \cong\left(\mathbb{P}^{n-1}\right)^{r} \tag{176}
\end{align*}
$$

and

$$
\begin{equation*}
H^{*}(\operatorname{Gr}(r, n)) \cong \mathbb{C}\left[\sigma_{1}, \ldots, \sigma_{r}\right] /\left(h_{n-r+1}, \ldots, h_{n}\right), \tag{177}
\end{equation*}
$$

$$
\begin{equation*}
H^{*}\left(\left(\mathbb{P}^{n-1}\right)^{r}\right) \cong \mathbb{C}\left[x_{1}, \ldots, x_{r}\right] /\left(x_{1}^{n}, \ldots, x_{r}^{n}\right) \tag{178}
\end{equation*}
$$

where $\sigma_{i}=\sigma_{i}\left(x_{1}, \ldots, x_{r}\right) \in \mathbb{C}\left[x_{1}, \ldots, x_{r}\right]^{\mathfrak{G}_{r}}$ are elementary symmetric functions and $h_{i}=h_{i}\left(x_{1}, \ldots, x_{r}\right) \in \mathbb{C}\left[x_{1}, \ldots, x_{r}\right]^{\mathfrak{C}_{r}}=\mathbb{C}\left[\sigma_{1}, \ldots, \sigma_{r}\right]$ are complete symmetric functions. Martin's formula in this case gives

$$
\begin{align*}
\int_{\operatorname{Gr}(r, n)} P\left(x_{1}, \ldots, x_{r}\right) & =\frac{1}{r!} \int_{\left(\mathbb{P}^{n-1}\right)^{r}} \prod_{i \neq j}\left(x_{i}-x_{j}\right) P\left(x_{1}, \ldots, x_{r}\right)  \tag{179}\\
& =\frac{(-1)^{r(r-1) / 2}}{r!} \int_{\left(\mathbb{P}^{n-1}\right)^{r}} \Delta^{2} \cup P\left(x_{1}, \ldots, x_{r}\right) \tag{180}
\end{align*}
$$

for any $P\left(x_{1}, \ldots, x_{r}\right) \in \mathbb{C}\left[x_{1}, \ldots, x_{r}\right]^{\mathfrak{C}_{r}}$ where $\boldsymbol{\Delta}:=\prod_{1 \leq i<j \leq r}\left(x_{i}-x_{j}\right)$.

## 9.4

The equivariant cohomology ring of $\operatorname{Gr}(r, n)$ with respect to the natural action of the diagonal maximal abelian subgroup $H \subset \mathrm{GL}_{n}$ is presented as

$$
\begin{equation*}
H_{H}^{\bullet}(\operatorname{Gr}(r, n) ; \mathbb{C}) \cong \mathbb{C}\left[\sigma_{1}, \ldots, \sigma_{r}, \lambda_{1}, \ldots, \lambda_{n}\right] /\left(h_{n-r+1}(\sigma, \lambda), \ldots, h_{n}(\sigma, \lambda)\right) \tag{181}
\end{equation*}
$$

where $h_{i}$ is the degree $2 i$ part of $c^{H}(\mathcal{S}) c^{H}(\mathcal{Q})-\prod_{i=1}^{n}\left(1+\lambda_{i}\right)$. Here $\mathcal{S}$ and $\mathcal{Q}$ are the tautological bundle and the universal quotient bundle respectively, and $c^{H}(-)$ stands for the $H$-equivariant total Chern class. Note that $\sigma_{i}:=c_{i}^{H}(\mathcal{S})$ is the elementary symmetric function of the $H$-equivariant Chern roots $x_{1}, \ldots, x_{r}$ of $\mathcal{S}$, and $c_{i}^{H}(\mathcal{Q})$ for $i=1, \ldots, n-r$ are expressed in terms of $\sigma_{1}, \ldots, \sigma_{r}$ and $\lambda_{1}, \ldots, \lambda_{n}$ by the condition $h_{1}=\cdots=h_{n-r}=0$. Martin's formula gives

$$
\begin{align*}
& \int_{\operatorname{Gr}(r, n)}^{H} P\left(\sigma_{1}, \ldots, \sigma_{r}\right)  \tag{182}\\
& =\sum_{1 \leq i_{1}<i_{2}<\cdots i_{r} \leq n} \operatorname{Res}_{x=\left(\lambda_{i_{1}}, \ldots, \lambda_{\left.i_{r}\right)}\right)} P\left(\sigma_{1}, \ldots, \sigma_{r}\right) \prod_{i \neq j}\left(x_{i}-x_{j}\right) \frac{d x_{1} \wedge \cdots \wedge d x_{r}}{\prod_{i=1}^{r} \prod_{j=1}^{n}\left(x_{i}-\lambda_{j}\right)} .
\end{align*}
$$

## 10 Quasimap Spaces for GIT Quotients

## 10.1

Let $G$ be a reductive algebraic group acting on an affine variety $W$ and fix a character $\theta$ of $G$. In this paper, we will always assume the following:
(1) Semi-stability implies stability.
(2) The semi-stable locus $W^{\text {ss }}$ is smooth and non-empty.
(3) The $G$-action on $W^{\text {ss }}$ is free (however, see [17] for allowing finite non-trivial stabilizers).
(4) The codimension of the unstable locus $W \backslash W^{\text {ss }}$ is greater than one.

The GIT quotient is defined by $W / / G:=W^{\text {ss }} / G$, which is an open substack of [ $W / G$ ].

## 10.2

A map $u: \mathbb{P}^{1} \rightarrow[W / G]$ to the quotient stack $[W / G]$ is pair $(P, \widetilde{u})$ of a principal $G$-bundle $P \rightarrow \mathbb{P}^{1}$ and a $G$-equivariant map $\tilde{u}: P \rightarrow W$. It is called a quasimap if the generic point of $\mathbb{P}^{1}$ is mapped to $W / / G \subset[W / G]$. A point in the inverse image of the unstable locus will be called a base point.

## 10.3

For a quasimap $u: \mathbb{P}^{1} \rightarrow[W / G]$ and a $G$-equivariant line bundle $L$ on $W$, the pullback $\widetilde{u}^{*} L$ is a $G$-equivariant line bundle on $P$, which descends to a line bundle $u^{*} L$ on $\mathbb{P}^{1}$. The degree of a quasimap $u: \mathbb{P}^{1} \rightarrow[W / G]$ is the map $\boldsymbol{d}: \operatorname{Pic}^{G} W \rightarrow \mathbb{Z}$ sending $L \in \operatorname{Pic}^{G} W$ to $\operatorname{deg} u^{*} L$.

## 10.4

An isomorphism of quasimaps $u=(P, \widetilde{u})$ and $u^{\prime}=\left(P^{\prime}, \tilde{u}^{\prime}\right)$ is an isomorphism $\varphi: P \rightarrow P^{\prime}$ of principal $G$-bundles such that $\widetilde{u}=\widetilde{u}^{\prime} \circ \varphi$. By [23, Theorem 7.1.6], the moduli functor for quasimaps of degree $\boldsymbol{d}$ is representable by a DeligneMumford stack, which will be denoted by $\mathbf{Q}(W / / G ; \boldsymbol{d})$. This stack is denoted by $\operatorname{Qmap}_{0,0}\left(W / / G, \boldsymbol{d} ; \mathbb{P}^{1}\right)$ in [23, Sect. 7.2] and $\mathrm{QG}_{0,0, \boldsymbol{d}}^{0+}(W / / G)$ in [22, Sect. 2.6]. Note that $\mathbf{Q}(W / / G)$ depends not only on $W / / G$ and $\boldsymbol{d}$ but also on $W, G$, and $\theta$.

## 10.5

Let $\mathbf{Q}$. $(W / / G ; \boldsymbol{d}) \subset \mathbf{Q}(W / / G ; \boldsymbol{d})$ be the substack parametrizing quasimaps such that $\left.u\right|_{\mathbb{P}^{\} \backslash\{0\}}$ is a constant map to $W / / G$. This implies that $0 \in \mathbb{P}^{1}$ is a base point of length $\boldsymbol{d}(\theta)$ by [23, Lemma 7.1.2]. This stack is denoted by $\mathbf{Q}_{0,0+\bullet}(W / / G, \boldsymbol{d})_{0}$ in [22, Sect.
4.1]. There is a natural map ev: Q. $(W / / G ; \boldsymbol{d}) \rightarrow W / / G$, called the evaluation map, which sends $u \in \mathbf{Q}$. $(W / / G ; \boldsymbol{d})$ to $u(\infty) \in W / / G$.

## 10.6

There is a natural $\mathbb{G}_{m}$-action on $\mathbf{Q}(W / / G ; \boldsymbol{d})$ coming from the standard $\mathbb{G}_{m}$-action on $\mathbb{P}^{1}$. As described in [22, Sect. 4.1], the fixed locus of this action is identified with the coproduct

$$
\begin{equation*}
\coprod_{\boldsymbol{d}_{1}+\boldsymbol{d}_{2}=\boldsymbol{d}} \text { Q. }\left(W / / G ; \boldsymbol{d}_{1}\right) \times_{W / / G} \mathbf{Q} .\left(W / / G ; \boldsymbol{d}_{2}\right) \tag{183}
\end{equation*}
$$

of fiber products with respect to the evaluation map.

## 10.7

If $W$ has at worst lci singularity, then $\mathbf{Q}(W / / G ; \boldsymbol{d})$ has a canonical perfect obstruction theory, which allows one to define the virtual fundamental cycle. The canonical perfect obstruction theory is $\left(\mathbb{R} \pi_{*} \mathrm{ev}^{*} T_{[W / G]}\right)^{\vee}$, where $T_{[W / G]}$ is the tangent complex of $[W / G]$, ev : $\mathbf{Q}(W / / G ; \boldsymbol{d}) \times \mathbb{P}^{1} \rightarrow[W / G]$ is the evaluation map, and $\pi: \mathbf{Q}(W / / G ; \boldsymbol{d}) \times \mathbb{P}^{1} \rightarrow \mathbf{Q}(W / / G ; \boldsymbol{d})$ is the first projection; see Theorem 7.2.2 of [22] or Sect. 5. The virtual fundamental cycle is an element of the Chow group of $\mathbf{Q}(W / / G ; \boldsymbol{d})$ whose degree is given by the virtual dimension

$$
\begin{equation*}
\text { virt. } \operatorname{dim} \mathbf{Q}(W / / G ; \boldsymbol{d})=\left\langle\boldsymbol{d}, \operatorname{det} T_{W}\right\rangle+\operatorname{dim} W / / G \tag{184}
\end{equation*}
$$

## 10.8

Since the stack $\mathbf{Q}$. $(W / / G ; \boldsymbol{d})$ is the union of connected components of the fixed locus of the $\mathbb{G}_{m}$-action, it has a perfect obstruction theory inherited from $\mathbf{Q}(W / / G ; \boldsymbol{d})$. The virtual push-forward

$$
\begin{equation*}
\mathrm{ev}_{*}^{\mathrm{virt}}(-):=\operatorname{PD}\left(\mathrm{ev}_{*}\left((-) \cap\left[\mathbf{Q}_{\mathbf{0}}(W / / G ; \boldsymbol{d})\right]^{\mathrm{virt}}\right)\right) \tag{185}
\end{equation*}
$$

along the evaluation map ev: $\mathbf{Q} .(W / / G ; \boldsymbol{d}) \rightarrow W / / G$ allows one to define the $I$ function

$$
\begin{equation*}
I(\boldsymbol{t} ; \mathbf{z}):=e^{p \cdot t / \mathbf{z}} \sum_{\boldsymbol{d} \in \operatorname{Eff}(W / / G)} e^{\boldsymbol{d} \cdot \boldsymbol{t}} I_{\boldsymbol{d}} \tag{186}
\end{equation*}
$$

by

$$
\begin{equation*}
I_{d}:=\operatorname{ev}_{*}^{\mathrm{virt}}\left(\frac{1}{\operatorname{Eul}^{\mathbb{G}_{m}}\left(N_{\mathbf{Q} \cdot(W / G ; \boldsymbol{d}) / \mathbf{Q}(W / G ; \boldsymbol{d})}^{\mathrm{virt}}\right)}\right) \tag{187}
\end{equation*}
$$

where the denominator is the $\mathbb{G}_{m}$-equivariant Euler class of the virtual normal bundle. Here $\boldsymbol{p}$ is a basis of $H^{2}(W / / G)$, and $\boldsymbol{t}$ is the coordinate of $H^{2}(W / / G)$ corresponding to $\boldsymbol{p}$.

An $H$-action on $W$ commuting with the $G$-action induces an $H$-action on Q. $(W / / G ; \boldsymbol{d})$, which allows one to define the $H$-equivariant $I$-function of $W / / G$.

## 10.9

There exits a $G$-space $V$ with a $G$-equivariant closed embedding $W \hookrightarrow V$. Let $\mathfrak{u}: \mathbf{Q}(W / / G ; \boldsymbol{d}) \times \mathbb{P}^{1} \rightarrow[W / G]$ be the universal quasimap. It consists of a principal $G$-bundle $\mathcal{P}$ on $\mathbf{Q}(W / / G ; \boldsymbol{d}) \times \mathbb{P}^{1}$ and a $G$-equivariant morphism $\tilde{\mathfrak{u}}: \mathcal{P} \rightarrow W$. Let $\mathcal{P}^{\prime}:=\left.\mathcal{P}\right|_{\mathbf{Q}(W / / G ; \boldsymbol{d}) \times\{\mathrm{pt}\}}$ be the restriction of $\mathcal{P}$ to a fiber of the second projection $\mathbf{Q}(W / / G ; \boldsymbol{d}) \times \mathbb{P}^{1} \rightarrow \mathbb{P}^{1}$. We write the Chern-Weil homomorphism defined by $\mathcal{P}^{\prime}$ as

$$
\begin{equation*}
\mathfrak{C W I : ~} \mathbb{C}[\mathfrak{g}]^{G} \rightarrow H^{*}(\mathbf{Q}(W / / G ; \boldsymbol{d})) . \tag{188}
\end{equation*}
$$

Note that $\mathbb{C}[\mathfrak{g}]^{G}$ is isomorphic to $\mathbb{C}[\mathfrak{t}]^{\mathscr{W}}$ by Chevalley restriction theorem. For $P \in$ $\mathbb{C}[t]^{\mathscr{W}}$, we set

$$
\begin{align*}
\langle P\rangle_{W / / G, \boldsymbol{d}} & :=\int_{[\mathbf{Q}(W / / G ; d))^{\mathrm{virt}}} \mathfrak{C W}(P),  \tag{189}\\
\langle P\rangle_{W / G} & :=\sum_{d \in \operatorname{Eff}(W / / G)} e^{\langle\boldsymbol{d}, \boldsymbol{t}\rangle}\langle P\rangle_{W / / G, \boldsymbol{d}} . \tag{190}
\end{align*}
$$

Conjecture 5 Suppose that $W \subset V$ is the zero locus of $G$ semi-invariant polynomials $f_{i}, i=1, \ldots, r$. Provided with conditions in Sects. 10.1 and 10.7, for any $P \in \mathbb{C}[\mathfrak{t}]^{\mathscr{W}}$, the generating function (190) of quasimap invariants coincides with the correlation function (5) of the A-twisted gauged linear sigma model up to an overall sign;

$$
\begin{equation*}
\langle P\rangle_{\mathrm{GLSM}}= \pm\langle P\rangle_{W / / G} . \tag{191}
\end{equation*}
$$

Here the potential of GLSM is given as a G-invariant function $\sum_{i} f_{i} p_{i}$ of $V \times \mathbb{A}^{r}$ where $p_{i}$ denotes ith coordinate of $\mathbb{A}^{r}$ with $R$-charge 2 .

### 10.10

By taking $\mathcal{P}^{\prime}$ to be the fiber over a fixed point of the natural $\mathbb{G}_{m}$-action on the domain curve $\mathbb{P}^{1}$, one can define $\mathbb{G}_{m}$-equivariant quasimap invariants $\langle P\rangle_{W / / G}^{\mathbb{G}_{m}}$. If $W$ has an action of an algebraic torus $H$ commuting with the action of $G$, then one can define $H \times \mathbb{G}_{m}$-equivariant quasimap invariants $\langle P\rangle_{W / / G}^{H \times \mathbb{G}_{m}}$.

## 11 Quasimap Spaces for Grassmannians

## 11.1

The quasimap space $\mathbf{Q}(\operatorname{Gr}(r, n) ; d)$ classifies pairs $(P, u)$ of a principal $\mathrm{GL}_{r}$-bundle $P$ and a $\mathrm{GL}_{r}$-equivariant map $u$. The choice of a principal $\mathrm{GL}_{r}$-bundle $P$ is equivalent to the choice of a vector bundle $S$ of rank $r$, and the choice of a $\mathrm{GL}_{r}$-equivariant map $u$ is equivalent to the choice of a map $\mathcal{S} \rightarrow \mathcal{O}_{\mathbb{P}^{1}}^{\oplus n}$, which is a sheaf injection since the generic point must go to the semi-stable locus (but not necessarily a morphism of vector bundles). The choice of a sheaf injection $\mathcal{S} \rightarrow \mathcal{O}_{\mathbb{P}^{1}}^{\oplus n}$ is equivalent to the choice of a surjection $\mathcal{O}_{\mathbb{P}^{1}}^{\oplus n} \rightarrow \mathcal{Q}$, where $\mathcal{Q}$ is a coherent sheaf whose Hilbert polynomial is $d+(n-r)(t+1)$. This is the same as the Hilbert polynomial of a locally free sheaf of rank $n-r$ and degree $d$, and one has an isomorphism

$$
\begin{equation*}
\mathbf{Q}(\operatorname{Gr}(r, n) ; d) \cong \operatorname{Quot}_{\mathbb{P}^{1}, d}\left(\mathcal{O}_{\mathbb{P}^{1}}^{\oplus n}, n-r\right) \tag{192}
\end{equation*}
$$

## 11.2

It is shown in [10, Lemma 1.2] that the subspace $\mathbf{Q} .(\operatorname{Gr}(r, n) ; d)$ of $\mathbf{Q}(\operatorname{Gr}(r, n) ; d)$ is decomposed into connected components as

$$
\begin{equation*}
\text { Q. }(\operatorname{Gr}(r, n) ; d)=\coprod_{|\boldsymbol{d}|=d} \mathbf{Q} .(\operatorname{Gr}(r, n) ; \boldsymbol{d}) \text {, } \tag{193}
\end{equation*}
$$

where $\boldsymbol{d}=\left(d_{1}, \ldots, d_{r}\right)$ runs over elements of $\mathbb{N}^{r}$ satisfying $|\boldsymbol{d}|:=d_{1}+\cdots+d_{r}=$ $d, d_{1} \leq d_{2} \leq \cdots \leq d_{r}$ and each connected component is isomorphic to the partial flag manifold

$$
\begin{equation*}
\mathbf{Q} .(\operatorname{Gr}(r, n) ; \boldsymbol{d}) \cong \mathrm{Fl}\left(m_{1}, \ldots, m_{k}, r, n\right), \tag{194}
\end{equation*}
$$

where $1 \leq m_{1}<m_{2}<\cdots<m_{k}=r$ denote the jumping indices;

$$
\begin{equation*}
0 \leq d_{1}=\cdots=d_{m_{1}}<d_{m_{1}+1}=\cdots=d_{m+2}<\cdots . \tag{195}
\end{equation*}
$$

Let $x_{1}, \ldots, x_{r}$ be the Chern roots of the dual of the universal subbundle on $\operatorname{Gr}(r, n)$. We also define $|\boldsymbol{x}|:=\sum_{i=1}^{r} x_{i}$ and $|\boldsymbol{d}|:=\sum_{i=1}^{r} d_{i}$ for $\boldsymbol{d}=\left(d_{1}, \ldots, d_{r}\right)$. The $I$ function can be computed by localization as

$$
\begin{equation*}
I_{\operatorname{Gr}(r, n)}(t ; \mathbf{z})=\sum_{\boldsymbol{d} \in \mathbb{N}^{r}}(-1)^{(r-1)|\boldsymbol{d}|} e^{(|\boldsymbol{d}|+|\boldsymbol{x}| / \mathbf{z}) t} I_{\boldsymbol{d}}(\mathbf{z}) \tag{196}
\end{equation*}
$$

where

$$
\begin{equation*}
I_{d}(\mathbf{z})=\frac{\prod_{1 \leq i<j \leq r}\left(x_{i}-x_{j}+\left(d_{i}-d_{j}\right) \mathbf{z}\right)}{\prod_{1 \leq i<j \leq r}\left(x_{i}-x_{j}\right) \prod_{i=1}^{r} \prod_{j=1}^{n} \prod_{l=1}^{d_{i}}\left(x_{i}+l \mathbf{z}\right)} . \tag{197}
\end{equation*}
$$

As shown in [10, p. 109], the $I$-function and the $J$-function agrees for $\operatorname{Gr}(r, n)$ just as in the case of projective spaces.

## 11.3

The Hori-Vafa conjecture [41] proved in [10] shows that the $I$-functions of $\left(\mathbb{P}^{n-1}\right)^{r}$ and $\operatorname{Gr}(r, n)$ are related by

$$
\begin{equation*}
I_{\mathrm{Gr}(r, n)}(t ; \mathbf{z})=\left.e^{-\sigma_{1}(r-1) \pi \sqrt{-1} / \mathbf{z}} \frac{\mathcal{D} I_{\left(\mathbb{P}^{n-1}\right)^{r}}(\boldsymbol{t} ; \mathbf{z})}{\boldsymbol{\Delta}}\right|_{t_{i}=t+(r-1) \pi \sqrt{-1}} \tag{198}
\end{equation*}
$$

where

$$
\begin{equation*}
\mathcal{D}:=\prod_{1 \leq i<j \leq r}\left(z \frac{\partial}{\partial t_{i}}-z \frac{\partial}{\partial t_{j}}\right) . \tag{199}
\end{equation*}
$$

## 11.4

As shown in [10], the equivariant $I$-function with respect to the natural action of $H=\left(\mathbb{G}_{m}\right)^{n}$ on $\operatorname{Mat}(r, n)$ is given by

$$
\begin{equation*}
I_{\mathrm{Gr}(r, n)}^{H}(t ; \mathrm{z})=e^{t \sigma \sigma_{1} / z} \sum_{d \in \mathbb{N}^{r}}(-1)^{(r-1)|d|} e^{|d| t \mid} \frac{\prod_{1 \leq i<j \leq r}\left(x_{i}-x_{j}+\left(d_{i}-d_{j}\right) \mathrm{z}\right)}{\prod_{1 \leq i<j \leq r}\left(x_{i}-x_{j}\right) \prod_{i=1}^{r} \prod_{j=1}^{n} \prod_{l=1}^{d_{i}}\left(x_{i}-\lambda_{j}+l \mathbf{z}\right)}, \tag{200}
\end{equation*}
$$

and the factorization gives

$$
\begin{equation*}
\sum_{d=0}^{\infty} e^{d \tau}\left\langle\left. e^{(t-\tau) \sigma_{1} / z}\right|_{\operatorname{Gr}(r, n), d} ^{H \times \mathbb{G}_{m}}=\int_{\operatorname{Gr}(r, n)}^{H} I_{\mathrm{Gr}(r, n)}^{H}(t ; \mathrm{z}) \cup I_{\mathrm{Gr}(r, n)}^{H}(\tau ;-\mathrm{z}) .\right. \tag{201}
\end{equation*}
$$

Here $\sigma_{1}=\sum_{i=1}^{r} x_{i}$ is the $H$-equivariant first Chern class of the vector bundle

$$
\begin{equation*}
\mathcal{S}^{\vee}=\left(\operatorname{Mat}(r, n) \times \mathbb{C}^{r}\right) / / G \tag{202}
\end{equation*}
$$

on $\operatorname{Gr}(r, n)$, where the $G$-action on $\mathbb{C}^{r}$ is the defining representation.

## 11.5

Let $\mathcal{V}$ be an equivariant vector bundle on $\operatorname{Gr}(r, n)$ associated with a representation $V$ of $\mathrm{GL}_{r}$. If $\mathcal{V}$ is globally generated and $\operatorname{det} \mathcal{V} \cong \omega_{\mathrm{Gr}(r, n)}^{\mathcal{V}}$, then the zero $Y:=s^{-1}(0)$ of a general section $s \in H^{0}(\mathcal{V})$ is a smooth Calabi-Yau manifold by a generalization of the theorem of Bertini [53, Theorem 1.10].

## 11.6

Let $\left[\operatorname{Mat}(r, n) / \mathrm{GL}_{r}\right]$ be the quotient stack containing $\operatorname{Gr}(r, n)$ as an open substack. The complete intersection $Y \subset \operatorname{Gr}(r, n)$ is an open substack of $\mathcal{Y}:=\left[Z / \mathrm{GL}_{r}\right]$, where $Z \subset \operatorname{Mat}(r, n)$ is the zero of the map $\widetilde{s}: \operatorname{Mat}(r, n) \rightarrow V$ underlying $s$. Indeed, $Y$ has a GIT quotient description $Y=Z / / \mathrm{GL}_{r}$, which allows us to define $\mathbf{Q}(Y ; d)$ and its virtual fundamental cycle as in Sect. 10. Let $\mathcal{S}_{\mathcal{Y}}^{\searrow}$ be the vector bundle on $\mathcal{Y}$ associated with the defining representation of $\mathrm{GL}_{r}$. Any point $p \in \mathbb{P}^{1}$ determines a map ev ${ }_{p}: \mathbf{Q}(Y ; d) \rightarrow \mathcal{Y}$ sending $f: \mathbb{P}^{1} \rightarrow \mathcal{Y}$ to $f(p) \in \mathcal{Y}$, and the Chern classes

$$
\begin{equation*}
\sigma_{i}:=c_{i}\left(\operatorname{ev}_{p}^{*} \mathcal{S}_{\mathcal{Y}}^{\vee}\right), \quad i=1, \ldots, r \tag{203}
\end{equation*}
$$

does not depend on the choice of $p \in \mathbb{P}^{1}$. For $P\left(\sigma_{1}, \ldots, \sigma_{r}\right) \in \mathbb{C}\left[\sigma_{1}, \ldots, \sigma_{r}\right]$, we set

$$
\begin{equation*}
\left\langle P\left(\sigma_{1}, \ldots, \sigma_{r}\right)\right\rangle_{Y, d}:=\int_{[\mathbf{Q}(Y ; d)]_{\text {vit }}} P\left(\sigma_{1}, \ldots, \sigma_{r}\right) \tag{204}
\end{equation*}
$$

and

$$
\begin{equation*}
\left\langle P\left(\sigma_{1}, \ldots, \sigma_{r}\right)\right\rangle_{Y}:=\sum_{d=0}^{\infty} e^{d t}\left\langle P\left(\sigma_{1}, \ldots, \sigma_{r}\right)\right\rangle_{Y, d} \tag{205}
\end{equation*}
$$

## 11.7

The equivariant $I$-function of $Y$ is given by

$$
\begin{equation*}
I_{Y}^{H}(t ; \mathbf{z})=\left.\sum_{\boldsymbol{d} \in \mathbb{N}^{r}} e^{(\boldsymbol{d}+\boldsymbol{x} / \mathbf{z}) \cdot t} I_{\boldsymbol{d}}(\boldsymbol{t} ; \mathbf{z})\right|_{t_{i}=t+(r-1) \pi \sqrt{-1}}, \tag{206}
\end{equation*}
$$

where

$$
\begin{equation*}
I_{d}(\boldsymbol{t} ; \mathbf{z}):=\frac{\prod_{\delta \in \Delta(V)} \prod_{l=1}^{\langle\delta, d\rangle}(\langle\boldsymbol{\delta}, \boldsymbol{x}\rangle+l \mathbf{z}) \prod_{1 \leq i<j \leq r}\left(x_{i}-x_{j}+\left(d_{i}-d_{j}\right) \mathbf{z}\right)}{\prod_{1 \leq i<j \leq r}\left(x_{i}-x_{j}\right) \prod_{i=1}^{r} \prod_{j=1}^{n} \prod_{l=1}^{d_{i}}\left(x_{i}-\lambda_{j}+l \mathbf{z}\right)}, \tag{207}
\end{equation*}
$$

where $\Delta(V)$ denotes the set of weights of $V$ and $\langle\boldsymbol{\delta}, \boldsymbol{x}\rangle$ denotes the first Chern class associated to the weight $\boldsymbol{\delta}$ (expressed in terms of the fundamental weights $x_{1}, \ldots, x_{r}$ of the maximal diagonal torus of $G$ ). Localization with respect to the natural $\mathbb{G}_{m^{-}}$action on $\mathbf{Q}(\operatorname{Gr}(r, n) ; d)$ shows

$$
\begin{equation*}
\left\langle\left. e^{(t-\tau) \sigma_{1} / \mathrm{z}}\right|_{Y} ^{H \times \mathbb{G}_{m}}=\int_{Y}^{H} I^{H}(t ; \mathbf{z}) \cup I^{H}(\tau ;-\mathbf{z})\right. \tag{208}
\end{equation*}
$$

just as in (169).

## 12 Residue Mirror Symmetry for Grassmannians

12.1

We define the abelianized quasimap space for $\operatorname{Gr}(r, n)$ by

$$
\begin{align*}
& \mathbf{Q}^{\mathrm{ab}}(\operatorname{Gr}(r, n) ; d):=\coprod_{|\boldsymbol{d}|=d} \mathbf{Q}^{\mathrm{ab}}(\operatorname{Gr}(r, n) ; \boldsymbol{d}),  \tag{209}\\
& \mathbf{Q}^{\mathrm{ab}}(\operatorname{Gr}(r, n) ; \boldsymbol{d}):=\mathbf{Q}\left(\mathbb{P}^{n-1} ; d_{1}\right) \times \cdots \times \mathbf{Q}\left(\mathbb{P}^{n-1} ; d_{r}\right), \tag{210}
\end{align*}
$$

where $\boldsymbol{d}$ runs over $\boldsymbol{d}=\left(d_{1}, \ldots, d_{r}\right) \in \mathbb{N}^{r}$ such that $|\boldsymbol{d}|:=d_{1}+\cdots+d_{r}=d$. An abelianized quasimap

$$
\begin{equation*}
\varphi\left(z_{1}, z_{2}\right)=\left(\left(\varphi_{i 1}\left(z_{1}, z_{2}\right), \ldots, \varphi_{i n}\left(z_{1}, z_{2}\right)\right) \in \mathbf{Q}\left(\mathbb{P}^{n-1} ; d_{i}\right)\right)_{i=1}^{r} \tag{211}
\end{equation*}
$$

defines a genuine map of degree $d$ if the matrix $\left(\varphi_{i j}\left(z_{1}, z_{2}\right)\right)_{i, j}$ has rank $r$ for any $\left(z_{1}, z_{2}\right) \neq 0$. For $P\left(\sigma_{1}, \ldots, \sigma_{r}\right) \in \mathbb{C}\left[\sigma_{1}, \ldots, \sigma_{r}\right]$, we set

$$
\begin{align*}
\left\langle P\left(\sigma_{1}, \ldots, \sigma_{r}\right)\right\rangle_{\operatorname{Gr}(r, n), \boldsymbol{d}}^{\mathrm{ab}} & =\frac{1}{r!} \int_{\mathbf{Q}^{\mathrm{ab}}(\operatorname{Gr}(r, n) ; \boldsymbol{d})} \prod_{i \neq j}\left(x_{i}-x_{j}\right) P\left(\sigma_{1}\left(x_{1}, \ldots, x_{r}\right), \ldots, \sigma_{r}\left(x_{1}, \ldots, x_{r}\right)\right),  \tag{212}\\
\left\langle P\left(\sigma_{1}, \ldots, \sigma_{r}\right)\right\rangle_{\operatorname{Gr}(r, n), d}^{\mathrm{ab}} & =\sum_{|\boldsymbol{d}|=d}\left\langle P\left(\sigma_{1}, \ldots, \sigma_{r}\right)\right\rangle_{\operatorname{Gr}(r, n), \boldsymbol{d}}^{\mathrm{ab}},  \tag{213}\\
\left\langle P\left(\sigma_{1}, \ldots, \sigma_{r}\right)\right\rangle_{\mathrm{Gr}(r, n)}^{\mathrm{ab}} & :=\sum_{d=0}^{\infty}(-1)^{(r-1) d} q^{d}\left\langle P\left(\sigma_{1}, \ldots, \sigma_{r}\right)\right\rangle_{\operatorname{Gr}(r, n), d}^{\mathrm{ab}} . \tag{214}
\end{align*}
$$

## 12.2

If we set $G:=\mathrm{GL}_{r}$ and $V:=\operatorname{Mat}(r, n)$, where $G$ acts naturally on $V$ and $\mathbb{G}_{m}$ acts trivially on $V$, then we have $Z_{\boldsymbol{d}}^{\text {vec }}(x)=\prod_{i \neq j}\left(x_{i}-x_{j}\right)$ and $Z_{d}^{\text {mat }}(x)=$ $\prod_{i=1}^{r}\left(x_{i}^{-d_{i}-1}\right)^{n}$, so that (5) gives the same result as (212);

$$
\begin{equation*}
\left\langle P\left(\sigma_{1}, \ldots, \sigma_{r}\right)\right\rangle_{\mathrm{GLSM}}=\left\langle P\left(\sigma_{1}, \ldots, \sigma_{r}\right)\right\rangle_{\mathrm{Gr}(r, n)}^{\mathrm{ab}} \tag{215}
\end{equation*}
$$

## 12.3

We write the ring homomorphism $\mathbb{C}\left[\sigma_{1}, \ldots, \sigma_{r}\right] \rightarrow \mathrm{QH}(\operatorname{Gr}(r, n))$ sending $\sigma_{i} \in$ $\mathbb{C}\left[\sigma_{1}, \ldots, \sigma_{r}\right]$ to $\sigma_{i} \in H^{*}(\operatorname{Gr}(r, n) ; \mathbb{C}) \cong \mathbb{C}\left[\sigma_{1}, \ldots, \sigma_{r}\right] /\left(h_{n-r+1}, \ldots, h_{n}\right)$ as $P\left(\sigma_{1}, \ldots, \sigma_{r}\right) \mapsto \stackrel{\circ}{P}\left(\sigma_{1}, \ldots, \sigma_{r}\right)$ just as in the case of $\mathbb{P}^{n-1}$.

Theorem 6 For any $P\left(\sigma_{1}, \ldots, \sigma_{r}\right) \in \mathbb{C}\left[\sigma_{1}, \ldots, \sigma_{r}\right]$, one has

$$
\begin{equation*}
\left\langle P\left(\sigma_{1}, \ldots, \sigma_{r}\right)\right\rangle_{\operatorname{Gr}(r, n)}^{\mathrm{ab}}=\int_{\operatorname{Gr}(r, n)} \stackrel{\circ}{P}\left(\sigma_{1}, \ldots, \sigma_{r}\right) \tag{216}
\end{equation*}
$$

Proof It follows from (34) that

$$
\begin{aligned}
& \left\langle P\left(\sigma_{1}, \ldots, \sigma_{r}\right)\right)_{\mathrm{G}_{\mathrm{G}(r, n)}^{\mathrm{ab}}} \\
& \quad=\frac{1}{r!} \sum_{d_{1}, \ldots, d_{r}=0}^{\infty}\left((-1)^{r-1} q\right)^{d_{1}+\cdots+d_{r}} \operatorname{Res} \prod_{i \neq j}\left(x_{i}-x_{j}\right) P\left(\sigma_{1}, \ldots, \sigma_{r}\right) \frac{d x_{1}}{x_{1}^{n\left(d_{1}+1\right)}} \wedge \cdots \wedge \frac{d x_{r}}{x_{r}^{n\left(d_{r}+1\right)}} \\
& \quad=\frac{1}{r!} \operatorname{Res} \prod_{i \neq j}\left(x_{i}-x_{j}\right) P\left(\sigma_{1}, \ldots, \sigma_{r}\right) \frac{d x_{1}}{x_{1}^{n}+(-1)^{r} q} \wedge \cdots \wedge \frac{d x_{r}}{x_{r}^{n}+(-1)^{r} q} \\
& \quad=\frac{1}{r!n^{r}} \sum_{x_{1}^{n}=(-1)^{r-1} q} \ldots \sum_{x_{r}^{n}=(-1)^{r-1} q} \prod_{i \neq j}\left(x_{i}-x_{j}\right) P\left(\sigma_{1}\left(x_{1}, \ldots, x_{r}\right), \ldots, \sigma_{r}\left(x_{1}, \ldots, x_{r}\right)\right) \\
& \quad=\int_{\operatorname{Gr}(r, n)} P\left(\sigma_{1}, \ldots, \sigma_{r}\right),
\end{aligned}
$$

where the last equality is the Vafa-Intriligator formula [61, Theorem 4.6].

## 12.4

Theorem 6 is related to intersection theory on the moduli space of vector bundles on a Riemann surface through a theorem of Witten [70], which states the existence of a ring isomorphism $\mathrm{QH}(\operatorname{Gr}(r, n)) /(q-1) \xrightarrow{\sim} R(U(r))_{n-r, n}$ from the quantum cohomology of $\operatorname{Gr}(r, n)$ at $q=1$ and the Verlinde algebra of $U(r)$ at $\mathrm{SU}(r)$ level $n-r$ and $U(1)$ level $n$.

## 12.5

We define the $\mathbb{G}_{m}$-equivariant correlator of $P\left(\sigma_{1}, \ldots, \sigma_{r}\right) \in \mathbb{C}\left[\sigma_{1}, \ldots, \sigma_{r}\right]$ by

$$
\begin{equation*}
\left\langle P\left(\sigma_{1}, \ldots, \sigma_{r}\right)\right\rangle_{\operatorname{Gr}(r, n)}^{\mathrm{ab}, \mathbb{G}_{m}}:=\left.\sum_{\boldsymbol{d} \in \mathbb{N}^{r}} e^{\boldsymbol{d} \cdot \boldsymbol{t}}\left\langle P\left(\sigma_{1}, \ldots, \sigma_{r}\right)\right\rangle_{\operatorname{Gr}(r, n), \boldsymbol{d}}^{\mathrm{ab}, \mathbb{G}_{m}}\right|_{t_{i}=t+(r-1) \pi \sqrt{-1}} \tag{217}
\end{equation*}
$$

where

$$
\begin{align*}
\left\langle P\left(\sigma_{1}, \ldots, \sigma_{r}\right)\right\rangle_{\operatorname{Gr}(r, n), \boldsymbol{d}}^{\mathrm{ab}, \mathbb{G}_{m}}:=\int_{\mathbf{Q}^{\mathrm{ab}}(\operatorname{Gr}(r, n) ; \boldsymbol{d})}^{\mathbb{G}_{m}} & \prod_{1 \leq i<j \leq n}\left(x_{i}-x_{j}\right)\left(x_{j}-x_{i}+\left(d_{j}-d_{i}\right) \mathbf{z}\right)  \tag{218}\\
& P\left(\sigma_{1}\left(x_{1}, \ldots, x_{r}\right), \ldots, \sigma_{r}\left(x_{1}, \ldots, x_{r}\right)\right) .
\end{align*}
$$

Since $\quad \mathbf{Q}\left(\left(\mathbb{P}^{n-1}\right)^{r} ; \boldsymbol{d}\right)^{\mathbb{G}_{m}}=\prod_{i=1}^{r} \mathbf{Q}\left(\mathbb{P}^{n-1} ; d_{i}\right)^{\mathbb{G}_{m}} \quad$ under $\quad \mathbf{Q}\left(\left(\mathbb{P}^{n-1}\right)^{r} ; \boldsymbol{d}\right)=$ $\prod_{i=1}^{r} \mathbf{Q}\left(\mathbb{P}^{n-1} ; d_{i}\right)$, we have a straightforward generalization of (62):

$$
\begin{equation*}
\sum_{\boldsymbol{d} \in \mathbb{N}^{r}} e^{\boldsymbol{d} \cdot \boldsymbol{\tau}}\left\langle e^{(\boldsymbol{t}-\boldsymbol{\tau}) \cdot \boldsymbol{x} / z}\right\rangle_{\left(\mathbb{P}^{n-1}\right)^{r}, \boldsymbol{d}}^{\mathbb{G}_{m}}=\int_{\left(\mathbb{P}^{n-1}\right)^{r}} I_{\left(\mathbb{P}^{n-1}\right)^{r}}(\boldsymbol{t} ; \mathbf{z}) \cup I_{\left(\mathbb{P}^{n-1}\right)^{r}}(\boldsymbol{\tau} ;-\mathbf{Z}), \tag{219}
\end{equation*}
$$

By acting $\mathcal{D}_{t}:=\prod_{1 \leq i<j \leq r}\left(\mathbf{z} \partial_{t_{i}}-\mathbf{z} \partial_{t_{j}}\right)$ and $-\mathcal{D}_{\tau}:=\prod_{1 \leq i<j \leq r}\left(-\mathbf{z} \partial_{\tau_{i}}+\mathbf{z} \partial_{\tau_{j}}\right)$ on both sides of (219) one obtains

$$
\begin{align*}
\sum_{\boldsymbol{d} \in \mathbb{N}^{r}} e^{\boldsymbol{d} \cdot t} & \left\langle\prod_{1 \leq i<j \leq r}\left(x_{i}-x_{j}\right) \prod_{1 \leq i<j \leq r}\left(\left(x_{j}+d_{j} \mathbf{z}\right)-\left(x_{i}+d_{i} \mathbf{z}\right)\right) \cdot e^{(\boldsymbol{t}-\tau) \cdot \boldsymbol{x} / z}\right\rangle_{\left(\mathbb{P}^{n-1}\right)^{r}}^{\mathbb{G}_{m}}  \tag{220}\\
& =\left\langle\left. e^{(\boldsymbol{t}-\tau) \cdot \boldsymbol{x} / z}\right|_{\operatorname{Gr}(r, n)} ^{\mathrm{ab}, \mathbb{G}_{m}}\right. \tag{221}
\end{align*}
$$

on the left hand side and

$$
\begin{equation*}
\int_{\left(\mathbb{P}^{n-1}\right)^{r}} \mathcal{D}_{t} I_{\left(\mathbb{P}^{n-1}\right)^{r}}(\boldsymbol{t} ; \mathbf{z}) \cup\left(-\mathcal{D}_{\tau}\right) I_{\left(\mathbb{P}^{n-1}\right)^{r}}(\boldsymbol{\tau} ;-\mathbf{Z}) \tag{222}
\end{equation*}
$$

on the right hand side. By setting $t_{i}=t+(r-1) \pi \sqrt{-1}, \tau_{i}=\tau+(r-1) \pi \sqrt{-1}$ and using (198), one obtains

$$
\begin{align*}
\left\langle\left. e^{(t-\tau) \sigma_{1} / \mathrm{z}}\right|_{\operatorname{Gr}(r, n)} ^{\mathrm{ab}, \mathbb{G}_{m}}\right. & =\frac{1}{r!} \int_{\left(\mathbb{P}^{n-1}\right) r^{r}} \Delta \cup I_{\operatorname{Gr}(r, n)}(t ; \mathbf{z}) \cup \Delta \cup I_{\mathrm{Gr}(r, n)}(\tau ;-\mathbf{z})  \tag{223}\\
& =\int_{\operatorname{Gr}(r, n)} I_{\operatorname{Gr}(r, n)}(t ; \mathbf{z}) \cup I_{\operatorname{Gr}(r, n)}(\tau ;-\mathbf{z}),
\end{align*}
$$

where the last equality is Martin's formula (182). On the other hand, localization with respect to the natural $\mathbb{G}_{m}$-action on the domain curve gives the factorization

$$
\begin{equation*}
\left\langle\left. e^{(t-\tau) \sigma_{1} / \mathrm{z}}\right|_{\operatorname{Gr}(r, n)} ^{\mathbb{G}_{m}}=\int_{\mathrm{Gr}(r, n)} I_{\mathrm{Gr}(r, n)}(t ; \mathbf{z}) \cup I_{\mathrm{Gr}(r, n)}(\tau ;-\mathbf{z}) .\right. \tag{224}
\end{equation*}
$$

Together with (223), this gives the equality

$$
\begin{equation*}
\left\langle\left. e^{(t-\tau) \sigma_{1} / \mathrm{z}}\right|_{\operatorname{Gr}(r, n)} ^{\mathrm{ab}, \mathbb{G}_{m}}=\left\langle\left. e^{(t-\tau) \sigma_{1} / \mathrm{z}}\right|_{\operatorname{Gr}(r, n)} ^{\mathbb{G}_{m}}\right.\right. \tag{225}
\end{equation*}
$$

of the abelianized correlator and the ordinary correlator.
For any $P(\boldsymbol{x}) \in \mathbb{C}\left[x_{1}, \ldots, x_{r}\right]^{\mathfrak{S}_{r}}$, the same argument gives

$$
\begin{align*}
& \left\langle\left. P(\boldsymbol{x}) e^{(\boldsymbol{t}-\boldsymbol{\tau}) \cdot \boldsymbol{x} / z}\right|_{\mathrm{Gr}(r, n)} ^{\mathrm{ab}, \mathbb{G}_{m}}\right.  \tag{226}\\
& \quad=\int_{\mathrm{Gr}(r, n)}\left(\sum_{\boldsymbol{d} \in \mathbb{N}^{r}} P(\boldsymbol{x}+\boldsymbol{d z}) I_{\mathrm{Gr}(r, n), \boldsymbol{d}}(\boldsymbol{t} ; \mathbf{z})\right) \cup\left(\sum_{\boldsymbol{d} \in \mathbb{N}^{r}} I_{\mathrm{Gr}(r, n), \boldsymbol{d}}(\boldsymbol{\tau} ;-\mathbf{z})\right) \\
& \quad=\left\langle\left. P(\boldsymbol{x}) e^{(\boldsymbol{t}-\boldsymbol{\tau}) \cdot \boldsymbol{x} / z}\right|_{\mathrm{Gr}(r, n)} ^{\mathbb{G}_{m}}\right.
\end{align*}
$$

where $t_{i}=t+(r-1) \pi \sqrt{-1}$ and $\tau_{i}=\tau+(r-1) \pi \sqrt{-1}$. By setting $t=\tau$ in (226), one obtains

$$
\begin{equation*}
\langle P(\boldsymbol{x})\rangle_{\operatorname{Gr}(r, n)}^{\mathrm{ab}, \mathbb{G}_{m}}=\langle P(\boldsymbol{x})\rangle_{\operatorname{Gr}(r, n)}^{\mathbb{G}_{m}} . \tag{227}
\end{equation*}
$$

Together with (215), this proves Conjecture 5 for Grassmannians.

## 12.6

Let $Y \subset \operatorname{Gr}(r, n)$ be the zero locus of a general section of a globally-generated vector bundle $\mathcal{V}$ on $\operatorname{Gr}(r, n)$ associated with a representation $V$ of $\mathrm{GL}_{r}$. We define the abelianized $\mathbb{G}_{m}$-equivariant Morrison-Plesser class of $Y$ by

$$
\begin{equation*}
\Phi_{d}^{\mathrm{ab}, \mathbb{G}_{m}}(Y ; \mathbf{z}):=\prod_{\delta \in \Delta(V)} \prod_{l=1}^{\langle\boldsymbol{\delta}, \boldsymbol{d}\rangle}(\langle\boldsymbol{\delta}, \boldsymbol{x}\rangle+l \mathbf{z}) \tag{228}
\end{equation*}
$$

For $P \in \mathbb{C}\left[\sigma_{1}, \ldots, \sigma_{r}\right]$, we set

$$
\begin{equation*}
\left\langle P\left(\sigma_{1}, \ldots, \sigma_{r}\right)\right\rangle_{Y}^{\mathrm{ab}, \mathbb{G}_{m}}:=\sum_{\boldsymbol{d} \in \mathbb{N}^{r}} q^{|\boldsymbol{d}|}\left\langle\left. P\left(\sigma_{1}, \ldots, \sigma_{r}\right) \Phi_{\boldsymbol{d}}^{\mathrm{ab}, \mathbb{G}_{m}}(Y ; \mathbf{z}) v\right|_{\operatorname{Gr}(r, n), \boldsymbol{d}} ^{\mathrm{ab}, \mathbb{G}_{m}},\right. \tag{229}
\end{equation*}
$$

where $v:=\prod_{\delta \in \Delta(V)}\langle\boldsymbol{\delta}, \boldsymbol{x}\rangle$ is the Euler class of the normal bundle of $Y$ in $\operatorname{Gr}(r, n)$. By the same reasoning as in Sect. 12.5 with the insertion of the abelizanized MorrisonPlesser class, one obtains

$$
\begin{equation*}
\left\langle P\left(\sigma_{1}, \ldots, \sigma_{r}\right)\right\rangle_{Y}=(-1)^{|\Delta(V)|}\left\langle P\left(\sigma_{1}, \ldots, \sigma_{r}\right)\right\rangle_{\mathrm{GLSM}} . \tag{230}
\end{equation*}
$$

Here, the identification between $q$ and the Fayet-Illiopoulos parameter $t^{\prime}$ is given by

$$
\begin{equation*}
q=(-1)^{\sum_{\delta \in \Delta(V)}\langle\delta, \mathbf{1}\rangle} e^{t^{\prime}} \tag{231}
\end{equation*}
$$

where $1:=(1, \cdots, 1) \in \mathbb{N}^{r}$.

## 12.7

As an example, consider the vector bundle of rank 3 on $\operatorname{Gr}(3,5)=\operatorname{Mat}(3,5) / / U(3)$ associated with the representation of $U(3)$ determined by the Young diagram

$$
\begin{equation*}
\lambda=\sharp \tag{232}
\end{equation*}
$$

This vector bundle is the tensor product $\wedge^{2} \mathcal{Q}(1)$ of the second exterior power $\wedge^{2} \mathcal{Q}$ of the universal quotient bundle $\mathcal{Q}$ on $\operatorname{Gr}(2,5) \cong \operatorname{Gr}(3,5)$ and the ample generator $\mathcal{O}(1)$ of the Picard group. One can immediately see from the Young diagram that the restriction of the representation of $U(3)$ associated with $\lambda$ to the diagonal maximal torus $T \cong\left(\mathbb{G}_{m}\right)^{3}$ is the direct sum $\rho_{1,2,2} \oplus \rho_{2,1,2} \oplus \rho_{2,2,1}$. The associated line bundle on the abelian quotient $\left(\mathbb{P}^{4}\right)^{3}$ is given by $\mathcal{O}(1,2,2) \oplus \mathcal{O}(2,1,2) \oplus \mathcal{O}(2,2,1)$.

The complete intersection in $\operatorname{Gr}(3,5)$ defined by $\wedge^{2} \mathcal{Q}(1)$ is a Calabi-Yau 3-fold of Picard number 1, which will be denoted by $Y$ henceforth. The Euler class of the normal bundle of $Y$ is

$$
\begin{equation*}
v:=\left(x_{1}+2 x_{2}+2 x_{3}\right)\left(2 x_{1}+x_{2}+2 x_{3}\right)\left(2 x_{1}+2 x_{2}+x_{3}\right) \tag{233}
\end{equation*}
$$

the abelianized Morrison-Plesser class is

$$
\begin{array}{r}
\Phi^{\mathrm{ab}}(Y ; \boldsymbol{d}):=\left(x_{1}+2 x_{2}+2 x_{3}\right)^{d_{1}+2 d_{2}+2 d_{3}}  \tag{234}\\
\left(2 x_{1}+x_{2}+2 x_{3}\right)^{2 d_{1}+d_{2}+2 d_{3}}\left(2 x_{1}+2 x_{2}+x_{3}\right)^{2 d_{1}+2 d_{2}+d_{3}}
\end{array}
$$

and the generating function for $\sigma_{1}^{3}$ is

$$
\begin{align*}
\left\langle\left.\sigma_{1}^{3}\right|_{Y} ^{\mathrm{ab}}=\right. & -\frac{1}{6} \sum_{d_{1}=0}^{\infty} \sum_{d_{2}=0}^{\infty} \sum_{d_{3}=0}^{\infty} q^{d_{1}+d_{2}+d_{3}} \operatorname{Res}\left(x_{1}+x_{2}+x_{3}\right)^{3}  \tag{235}\\
& \left(x_{1}-x_{2}\right)^{2}\left(x_{1}-x_{3}\right)^{2}\left(x_{2}-x_{3}\right)^{2} \Phi^{\mathrm{ab}}(Y ; \boldsymbol{d}) v \frac{d x_{1}}{x_{1}^{n\left(d_{1}+1\right)}} \wedge \frac{d x_{2}}{x_{2}^{n\left(d_{2}+1\right)}} \wedge \frac{d x_{3}}{x_{3}^{n\left(d_{3}+1\right)}} \\
= & \frac{25(1-q)}{(1+q)\left(1-123 q+q^{2}\right)} . \tag{236}
\end{align*}
$$

This matches the Yukawa coupling of the mirror computed by Miura [51, Sect. 5.2].

## 12.8

When $\mathcal{V}$ is a direct sum of line bundles, the mirror of $Y$ is constructed by toric degenerations $[6,7]$. It is an interesting problem to compare the generating function (205) with the Yukawa coupling of this mirror.

## 13 Bethe/Gauge Correspondence

## 13.1

Let $V_{1}$ and $W_{1}$ be Hermitian vector spaces of dimensions $r$ and $n$. The unitary group $U(r)$ acts naturally on $V_{1}$ and trivially on $W_{1}$, inducing an action on $T^{*} \operatorname{Hom}\left(V_{1}, W_{1}\right) \cong \operatorname{Hom}\left(V_{1}, W_{1}\right) \oplus \operatorname{Hom}\left(W_{1}, V_{1}\right)$. The real and complex moment maps for this action are given by
$\mu_{\mathbb{R}}: \operatorname{Hom}\left(W_{1}, V_{1}\right) \oplus \operatorname{Hom}\left(V_{1}, W_{1}\right) \rightarrow \operatorname{End}\left(V_{1}\right), \quad\left(i_{1}, j_{1}\right) \mapsto \frac{\sqrt{-1}}{2}\left(i_{1} i_{1}^{*}-j_{1}^{*} j_{1}\right)$,
$\mu_{\mathbb{C}}: \operatorname{Hom}\left(W_{1}, V_{1}\right) \oplus \operatorname{Hom}\left(V_{1}, W_{1}\right) \rightarrow \operatorname{End}\left(V_{1}\right), \quad\left(i_{1}, j_{1}\right) \mapsto i_{1} j_{1}$.
If $\left(i_{1}, j_{1}\right) \in \mu_{\mathbb{R}}^{-1}\left(\zeta \sqrt{-1} \mathrm{id}_{V_{1}}\right)$ for $\zeta<0$, then $j_{1}$ is injective. If $\left(i_{1}, j_{1}\right) \in \mu_{\mathbb{C}}^{-1}(0)$, then $i_{1}$ descends to a map $W_{1} / \operatorname{Im} j_{1} \rightarrow V_{1}$. It follows that the hyperKähler quotient is isomorphic to $T^{*} \operatorname{Gr}(r, n)$;

$$
\begin{equation*}
\left(\zeta_{\mathbb{R}}^{-1}\left(\zeta \sqrt{-1} \operatorname{id}_{V_{1}}\right) \cap \mu_{\mathbb{C}}^{-1}(0)\right) / U(r) \cong T^{*} \operatorname{Gr}(r, n) \tag{239}
\end{equation*}
$$

This suggests that the gauged linear sigma model with the gauge group $U(r)$ and the representation $V:=\operatorname{Hom}\left(W_{1}, V_{1}\right) \oplus \operatorname{Hom}\left(V_{1}, W_{1}\right) \oplus \operatorname{End}\left(V_{1}\right)$ describes the quantum cohomology of $T^{*} \operatorname{Gr}(r, n)$. Here $\operatorname{End}\left(V_{1}\right)$ is the Lagrange multiplier for the complex moment map equation, and the potential is given by

$$
\begin{equation*}
V \ni\left(i_{1}, j_{1}, P\right) \mapsto \operatorname{tr}\left(P i_{1} j_{1}\right) \tag{240}
\end{equation*}
$$

Let $H:=H_{1} \times H_{2}$ be the product of

- the diagonal maximal torus $H_{1}$ of $U(n)$, acting on $\operatorname{Hom}\left(W_{1}, V_{1}\right)$ and $\operatorname{Hom}\left(V_{1}, W_{1}\right)$ through the natural action on $W_{1}$, and trivially on $\operatorname{End}\left(V_{1}\right)$, and
- the group $H_{2}=U(1)$ acting trivially on $\operatorname{Hom}\left(W_{1}, V_{1}\right)$, by scalar multiplication on $\operatorname{Hom}\left(V_{1}, W_{1}\right)$, and by inverse scalar multiplication on $\operatorname{End}\left(V_{1}\right)$.

One has

$$
\begin{align*}
Z_{d}^{\mathrm{vec}}(x)= & \prod_{1 \leq i \neq j \leq r}\left(x_{i}-x_{j}\right)  \tag{241}\\
Z_{d}^{\mathrm{mat}}(x)= & \prod_{j=1}^{n} \prod_{i=1}^{r}\left(x_{i}-\lambda_{j}\right)^{-d_{i}-1}  \tag{242}\\
& \times \prod_{j=1}^{n} \prod_{i=1}^{r}\left(-x_{i}+\lambda_{j}-\mu\right)^{-\left(-d_{i}\right)-1} \tag{243}
\end{align*}
$$

$$
\begin{equation*}
\times \prod_{1 \leq i \neq j \leq r}\left(x_{i}-x_{j}+\mu\right)^{2-\left(d_{i}-d_{j}\right)-1} \tag{244}
\end{equation*}
$$

so that the $H$-equivariant correlator of $P \in \mathbb{C}\left[x_{1}, \ldots, x_{r}\right]^{\mathfrak{S}_{r}}$ is given by

$$
\begin{align*}
\langle P\rangle_{\mathrm{GLSM}}^{H}=\frac{1}{r!} \sum_{d_{1}=0}^{\infty} \cdots & \sum_{d_{r}=0}^{\infty}\left((-1)^{r-1} e^{t}\right)^{d_{1}+\cdots+d_{r}}  \tag{245}\\
\operatorname{Res}[ & \frac{\prod_{1 \leq i \neq j \leq r}\left(x_{i}-x_{j}\right)}{\prod_{1 \leq i, j \leq r}\left(x_{i}-x_{j}+\mu\right)^{\left(d_{i}-d_{j}-1\right)}} \\
& \left.\frac{\prod_{j=1}^{n} \prod_{i=1}^{r}\left(-x_{i}+\lambda_{j}-\mu\right)^{d_{i}-1}}{\prod_{j=1}^{n} \prod_{i=1}^{r}\left(x_{i}-\lambda_{j}\right)^{d_{i}+1}} P d x_{1} \wedge \cdots \wedge d x_{r}\right]
\end{align*}
$$

where Res denotes the sum of residues at the points where $x_{i}$ is one of $\lambda_{j}$ for $i=$ $1, \ldots, r$ and $j=1, \ldots, n$ (there are $n^{r}$ such points). This can formally be regarded as an equivariant integration over the projective space of dimension $\sum_{i=1}^{r}\left(d_{i}+1\right)-1$, and it is an interesting problem to give a geometric interpretation.

The effective potential (10) of this gauged linear sigma model is given by

$$
\begin{align*}
W_{\mathrm{eff}}(\boldsymbol{x} ; t)= & W_{\mathrm{FI}}\left(\boldsymbol{x} ; t^{\prime}\right)+W_{\mathrm{vec}}(\boldsymbol{x})+W_{\mathrm{mat}}(\boldsymbol{x})  \tag{246}\\
W_{\mathrm{FI}}(\boldsymbol{x} ; t)= & t\left(x_{1}+\cdots+x_{r}\right)  \tag{247}\\
W_{\mathrm{vec}}(\boldsymbol{x})= & -\pi \sqrt{-1} \sum_{1 \leq i<j \leq r}\left(x_{j}-x_{i}\right)  \tag{248}\\
= & -\pi \sqrt{-1} \sum_{i=1}^{r}(2 i-r-1) x_{i} \\
W_{\mathrm{mat}}(\boldsymbol{x})= & -\sum_{i=1}^{r} \sum_{j=1}^{n}\left(x_{i}-\lambda_{j}\right)\left(\log \left(x_{i}-\lambda_{j}\right)-1\right)  \tag{249}\\
& \quad-\sum_{i=1}^{r} \sum_{j=1}^{n}\left(-x_{i}+\lambda_{j}-\mu\right)\left(\log \left(-x_{i}+\lambda_{j}-\mu\right)-1\right) \\
& \quad-\sum_{i=1}^{r} \sum_{j=1}^{r}\left(x_{i}-x_{j}+\mu\right)\left(\log \left(x_{i}-x_{j}+\mu\right)-1\right)
\end{align*}
$$

where $\lambda_{j}$ and $\mu$ are equivariant parameters for the actions of $H_{1}$ and $H_{2}$ respectively. Note that

$$
\begin{equation*}
e^{\partial W_{\mathrm{eff}} / \partial x_{i}}=e^{t} \cdot(-1)^{2 i-r-1} \cdot \prod_{j=1}^{n}\left(x_{i}-\lambda_{j}\right)^{-1} \prod_{j=1}^{n}\left(-x_{i}+\lambda_{j}-\mu\right) \prod_{j \neq i} \frac{x_{j}-x_{i}+\mu}{x_{i}-x_{j}+\mu} \tag{250}
\end{equation*}
$$

$$
\begin{equation*}
=e^{t+n \pi \sqrt{-1}} \prod_{j=1}^{n} \frac{x_{i}-\lambda_{j}+\mu}{x_{i}-\lambda_{j}} \prod_{j \neq i} \frac{x_{i}-x_{j}-\mu}{x_{i}-x_{j}+\mu} \tag{251}
\end{equation*}
$$

so that the equations $e^{\partial_{x_{i}} W_{\text {eff }}}=1, i=1, \ldots, r$ gives

$$
\begin{equation*}
\prod_{j=1}^{n} \frac{x_{i}-\lambda_{j}}{x_{i}-\lambda_{j}+\mu}=e^{t+n \pi \sqrt{-1}} \prod_{j \neq i} \frac{x_{i}-x_{j}-\mu}{x_{i}-x_{j}+\mu} \tag{252}
\end{equation*}
$$

By taking the sum over $d_{i}$ just as in the proof of Corollary 2, one obtains

$$
\begin{align*}
\langle P\rangle_{\mathrm{GLSM}}^{H}=\frac{1}{r!} \operatorname{Res}[ & \frac{1}{\prod_{i=1}^{r}\left(\left(1-e^{\partial_{x_{i}} W_{\mathrm{eff}}}\right) \prod_{j=1}^{n}\left(x_{i}-\lambda_{j}\right)\right)}  \tag{253}\\
& \left.\frac{\prod_{1 \leq i \neq j \leq r}\left(x_{i}-x_{j}\right) \prod_{1 \leq i, j \leq r}\left(x_{i}-x_{j}+\mu\right)}{\prod_{i=1}^{r} \prod_{j=1}^{n}\left(-x_{i}+\lambda_{j}-\mu\right)} P d x_{1} \wedge \cdots \wedge d x_{r}\right]
\end{align*}
$$

where Res denotes the sum of residues at the roots of the Eqs. (252).

## 13.2

The Heisenberg model, also known as the homogeneous $X X X_{\frac{1}{2}}$ model, is the $\mathrm{SU}(2)$ spin chain model with Hamiltonian

$$
\begin{equation*}
H=\sum_{i=1}^{n} \boldsymbol{S}_{i} \cdot \boldsymbol{S}_{i+1} \tag{254}
\end{equation*}
$$

where $\boldsymbol{S}_{i}=\left(S_{i}^{x}, S_{i}^{y}, S_{i}^{z}\right)=\left(\sigma_{i}^{x} / 2, \sigma_{i}^{y} / 2, \sigma_{i}^{z} / 2\right)$ are halves of Pauli matrices acting on the $i$ th factor of the Hilbert space $\mathcal{H}:=\left(\mathbb{C}^{2}\right)^{\otimes n}$ and

$$
\begin{equation*}
\boldsymbol{S}_{i} \cdot \boldsymbol{S}_{i+1}:=S_{i}^{x} S_{i+1}^{x}+S_{i}^{y} S_{i+1}^{y}+S_{i}^{z} S_{i+1}^{z} \tag{255}
\end{equation*}
$$

The total spin

$$
\begin{equation*}
S^{z}:=\sum_{i=1}^{n} S_{i}^{z} \tag{256}
\end{equation*}
$$

clearly commutes with the Hamiltonian, and we restrict to the $S^{z}$-eigenspace $\mathcal{H}_{r} \subset \mathcal{H}$ with eigenvalue $(-n+r) / 2$. We impose the quasi-periodicity condition

$$
\begin{equation*}
\boldsymbol{S}_{n+1}=e^{\sqrt{-1} \vartheta S_{1}^{z}} \boldsymbol{S}_{1} e^{-\sqrt{-1} \vartheta S_{1}^{z}} . \tag{257}
\end{equation*}
$$

Introduce variables $\boldsymbol{x}=\left(x_{1}, \ldots, x_{r}\right)$ related to quasi-momenta $\boldsymbol{p}=\left(p_{1}, \ldots, p_{r}\right)$ by

$$
\begin{equation*}
e^{\sqrt{-1} p_{i}}=\frac{x_{i}+\frac{\sqrt{-1}}{2}}{x_{i}-\frac{\sqrt{-1}}{2}} \tag{258}
\end{equation*}
$$

Then $H$-eigenspaces in $\mathcal{H}_{r}$ correspond bijectively to solutions of the Bethe equation

$$
\begin{equation*}
\left(\frac{x_{i}+\frac{\sqrt{-1}}{2}}{x_{i}-\frac{\sqrt{-1}}{2}}\right)^{n}=e^{\sqrt{-1} \vartheta} \prod_{j \neq i} \frac{x_{i}-x_{j}+\sqrt{-1}}{x_{i}-x_{j}-\sqrt{-1}} \tag{259}
\end{equation*}
$$

with eigenvalues $n-2 r+2 \sum_{i=1}^{r} \cos p_{i}$. The integrability comes from factorization of many-body $S$-matrix into the product of the 2 -body $S$-matrix given by

$$
\begin{equation*}
S\left(p_{i}, p_{j}\right)=1-2 e^{\sqrt{-1} p_{j}}+e^{\sqrt{-1}\left(p_{i}+p_{j}\right)} \tag{260}
\end{equation*}
$$

See e.g. [62] and references therein for Bethe ansatz for the quasi-periodic Heisenberg model. The Bethe equation (259) coincides with (252) under $\lambda_{j}=\frac{\sqrt{-1}}{2}$, $j=1, \ldots, n, \mu=-\sqrt{-1}$, and $\vartheta=-\sqrt{-1} t+n / 2$. This observation and its generalizations is called Bethe/gauge correspondence [56]. The relation between classical/quantum cohomology of Grassmannians and integrable systems is studied in [14, 36, 49, 57].

## 14 Quasimaps and Instantons

## 14.1

As explained in [29, Sect. 2.3], the moduli space of framed instantons on $\mathbb{C} \times$ $[\mathbb{C} /(\mathbb{Z} / n \mathbb{Z})]$ is isomorphic to the Nakajima quiver variety associated with the chainsaw quiver shown in Fig. 1.

## 14.2

Representations of the chainsaw quiver satisfying $\operatorname{dim} V_{n}=0$ are in one-to-one correspondence with representations of the handsaw quiver shown in Fig. 2. It is shown in [29, Sect. 2.3] (see also [55, Sect. 3] for an exposition) that the Nakajima quiver variety associated with the handsaw quiver is isomorphic to the parabolic Laumon space parametrizing flags


Fig. 1 The chainsaw quiver


Fig. 2 The handsaw quiver

$$
\begin{equation*}
0=E_{0} \subset E_{1} \subset \cdots \subset E_{n-1} \subset E_{n}=W \otimes_{\mathbb{C}} \mathcal{O}_{\mathbb{P}^{1}} \tag{261}
\end{equation*}
$$

of locally free sheaves on $\mathbb{P}^{1}$ such that $\operatorname{rank} E_{i}=\sum_{j \leq i} \operatorname{dim} W_{j}, \operatorname{deg} E_{i}=-\operatorname{dim} V_{i}$, and the flag at $\infty \in \mathbb{P}^{1}$ is equal to the standard flag $0 \subset W_{1} \subset W_{1} \oplus W_{2} \subset \cdots \subset$ $W_{1} \oplus W_{2} \oplus \cdots \oplus W_{n-1} \subset W$. This coincides with the space of based quasimaps to partial flag varieties, i.e., quasimaps with specified value at infinity.

## 15 Quasimaps and Monopoles

## 15.1

Let $G$ be a compact Lie group with a maximal torus $H$. A monopole on $\mathbb{R}^{3}$ is a pair $(A, \Phi)$ of a connection $A$ on a principal $G$-bundle $P$ and a section $\Phi$ of $P \times{ }_{G} \mathfrak{g}$ satisfying the Bogomolny equation

$$
\begin{equation*}
F_{A}=* d_{A} \Phi . \tag{262}
\end{equation*}
$$

In order for the curvature to have a finite $L^{2}$-norm, it is natural to demand that the restriction of $\Phi$ to a sphere with large radius tends to a map to a fixed adjoint orbit $\mathcal{O} \cong G / H \cong G_{\mathbb{C}} / P$. The homotopy class $k \in \pi_{2}\left(G_{\mathbb{C}} / P\right)$ of the resulting map is called the charge of the monopole.

## 15.2

A choice of a gauge satisfying a certain boundary condition at infinity is called a framing of the monopole. The framed moduli space is a principal $H$-bundle over the unframed moduli space. The framed moduli space has a natural hyperKähler structure coming from the dimensional reduction of the anti-self-dual equation in dimension 4.

## 15.3

Monopoles on $\mathbb{R}^{3}$ are related to
(1) spectral curves on $T \mathbb{P}^{1}$,
(2) Nahm's equation

$$
\begin{equation*}
\frac{d T_{i}}{d s}=\epsilon_{i j k}\left[T_{j}, T_{k}\right], \quad i=1,2,3 \tag{263}
\end{equation*}
$$

for $T_{i} \in C^{\infty}((0,2), \operatorname{Mat}(k, k ; \mathbb{C}))$, and
(3) based quasimaps from $\mathbb{P}^{1}$ to $G_{\mathbb{C}} / P$ of degree $\boldsymbol{k}$.
(1) comes from the twistor correspondence [38, 39], and (2) comes from Nahm transform [54]. (3) is proved for $\mathrm{SU}(2)$ in [27], and the general case can be found in [44, 45] and references therein.

## 16 Quasimaps and Vortices

## 16.1

Let $X$ be a Kähler manifold, $(E, h)$ be a Hermitian vector bundle on $X$, and $\tau$ be a positive real number. The Yang-Mills-Higgs functional sends a pair $(A, \phi)$ of a unitary connection $d_{A}$ of $(E, h)$ and a section $\phi$ of $E$ to

$$
\begin{equation*}
\mathscr{H} \mathscr{H} \mathscr{H}(A, \phi)=\left\|F_{A}\right\|_{L^{2}}^{2}+\left\|d_{A} \phi\right\|_{L^{2}}^{2}+\frac{1}{4}\left\|\phi \otimes \phi^{*}-\tau\right\|_{L^{2}}^{2} . \tag{264}
\end{equation*}
$$

By [12, Proposition 2.1], one has

$$
\begin{array}{r}
\mathscr{H} \mathscr{M} \mathscr{H}(A, \phi)=4\left\|F^{0,2}\right\|_{L^{2}}^{2}+2\left\|\bar{\partial}_{A} \phi\right\|_{L^{2}}^{2}+\left\|\sqrt{-1} \Lambda F+\frac{1}{2} \phi \otimes \phi^{*}-\frac{\tau}{2}\right\|_{L^{2}}^{2}  \tag{265}\\
+\tau \int_{X} \sqrt{-1} \operatorname{tr} F \wedge \omega^{[n-1]}+\int_{X} \operatorname{tr} F \wedge F \wedge \omega^{[n-2]} .
\end{array}
$$

where $\omega^{[k]}:=\omega^{k} /(k!)$ and $\Lambda$ is the dual Lefschetz operator.

## 16.2

Assume that $X$ is a projective curve, so that

$$
\begin{equation*}
\operatorname{deg}(E)=\frac{\sqrt{-1}}{2 \pi} \operatorname{tr} F \tag{266}
\end{equation*}
$$

Then (265) immediately implies the Bogomolny-Prasad-Sommerfield inequality

$$
\begin{equation*}
\mathscr{U} \mathscr{H}(A, \phi) \geq 2 \pi \tau \operatorname{deg}(E), \tag{267}
\end{equation*}
$$

and the equality holds if and only if the vortex equation

$$
\begin{align*}
F^{0,2} & =0  \tag{268}\\
\bar{\partial}_{A} \phi & =0  \tag{269}\\
-\sqrt{-1} \Lambda F & =\frac{1}{2}\left(\phi \otimes \phi^{*}-\tau \mathrm{id}_{E}\right) \tag{270}
\end{align*}
$$

is satisfied. Equations (268) and (269) are holomorphicities for $E$ and $\phi$, and (270) is a generalization of the constant central curvature equation.

## 16.3

By taking the trace of (270) and integrating over $X$, one obtains

$$
\begin{equation*}
-2 \pi \operatorname{deg}(E)=\frac{1}{2}\|\phi\|_{L^{2}}^{2}-\frac{1}{2} \tau \operatorname{rank}(E) \operatorname{vol}(X) \tag{271}
\end{equation*}
$$

so that the condition

$$
\begin{equation*}
\tau \geq \frac{4 \pi \operatorname{deg}(E)}{\operatorname{rank}(E) \operatorname{vol}(X)} \tag{272}
\end{equation*}
$$

is necessary for (270) to have a solution.

## 16.4

The slope of a holomorphic vector bundle $E$ is defined by

$$
\begin{equation*}
\mu(E)=\frac{\operatorname{deg}(E)}{\operatorname{rank}(E)} \tag{273}
\end{equation*}
$$

For a holomorphic section $\phi$ of $E$, we set

$$
\begin{aligned}
\hat{\mu}(E):= & \sup \left\{\mu\left(E^{\prime}\right) \mid E^{\prime} \text { is a reflexive subsheaf of } E \text { of rank less than } E\right\}, \\
\mu_{M}(E):= & \max \{\hat{\mu}(E), \mu(E)\}, \\
\mu_{m}(E, \phi):= & \inf \left\{\left.\frac{\operatorname{rank}(E) \mu(E)-\operatorname{rank}\left(E^{\prime}\right) \mu\left(E^{\prime}\right)}{\operatorname{rank}(E)-\operatorname{rank}\left(E^{\prime}\right)} \right\rvert\,\right. \\
& \left.E^{\prime} \text { is a reflexive subsheaf of } E \text { such that rank } E^{\prime}<\operatorname{rank} E \text { and } \phi \in \Gamma\left(E^{\prime}\right)\right\} .
\end{aligned}
$$

A pair $(E, \phi)$ of a holomorphic vector bundle $E$ and its holomorphic section $\phi$ is said to be stable if

$$
\begin{equation*}
\mu_{M}(E)<\mu_{m}(E, \phi) \tag{274}
\end{equation*}
$$

Theorem 7 ([13, Theorem 2.1.6]) Let (E, $\phi$ ) be a pair of a holomorphic vector bundle and its holomorphic section. If there exists a Hermitian metric on E satisfying the vortex equation, then one has either of the following:
(i) $(E, \phi)$ is stable and satisfies

$$
\begin{equation*}
\mu_{M}<\frac{\tau \operatorname{Vol}(X)}{4 \pi}<\mu_{m}(\phi) \tag{275}
\end{equation*}
$$

(ii) $E$ has a direct sum decomposition $E=E_{\phi} \oplus E^{\prime}, \phi$ is an element of $H^{0}\left(E_{\phi}\right) \subset$ $H^{0}(E),\left(E_{\phi}, \phi\right)$ satisfies (i) above, and $E^{\prime}$ is the direct sum of stable vector bundles of slope $\tau \operatorname{Vol}(X) / 4 \pi$.

Theorem 8 ([13, Theorem 3.1.1]) Let (E, $\phi$ ) be a stable pair of a holomorphic vector bundle and its holomorphic section. Then for any real number $\tau$ satisfying (275), there exists a Hermitian metric on E satisfying (270).

Bradlow proved these results not only for projective curves but also for compact Kähler manifolds.

## 16.5

Vortex Eq. (270) admits the following generalization, which also contains Hitchin's self-duality equation [40] as a special case. Let $Q=\left(Q_{0}, Q_{1}, s, t\right)$ be a quiver and $M=\left(M_{a}\right)_{a \in Q_{1}}$ be a collection of vector bundles on $X$ labeled by $Q_{1}$. An $M$-twisted $Q$-sheaf on $X$ is a pair $R=\left(\left(E_{v}\right)_{v \in Q_{0}},\left(\phi_{a}\right)_{a \in Q_{1}}\right)$ of a collection $\left(E_{v}\right)_{v \in Q_{0}}$ of vector bundles labeled by $Q_{0}$ and a collection

$$
\begin{equation*}
\left(\phi_{a}\right)_{a \in Q_{1}} \in \prod_{a \in Q_{1}} \operatorname{Hom}\left(E_{s(a)} \otimes M_{a}, E_{t(a)}\right) \tag{276}
\end{equation*}
$$

of morphisms labeled by $Q_{1}$.
Given a collection $\left(E_{v}\right)_{v \in Q_{0}}$ of holomorphic vector bundles on a Kähler manifold $X$, another collection $\left(M_{a}\right)_{a \in Q_{1}}$ of holomorphic vector bundles on $X$, a collection $\sigma=\left(\sigma_{v}\right)_{v \in Q_{0}}$ of positive real numbers, and a collection $\tau=\left(\tau_{v}\right)_{v \in Q_{0}}$ of real numbers, the equation

$$
\begin{equation*}
\sigma_{v} \sqrt{-1} \Lambda F_{v}+\sum_{t(a)=v} \phi_{a} \circ \phi_{a}^{*}-\sum_{s(a)=v} \phi_{a}^{*} \circ \phi_{a}=\tau_{v} \operatorname{id}_{E_{v}} \tag{277}
\end{equation*}
$$

for Hermitian metrics on $\left(E_{v}\right)_{v \in Q_{0}}$ is called the $M$-twisted quiver $(\sigma, \tau)$-vortex equation.

The $(\sigma, \tau)$-degree and the $(\sigma, \tau)$-slope of an $M$-twisted $Q$-sheaf $R$ is defined by

$$
\begin{align*}
\operatorname{deg}_{\sigma, \tau}(R) & =\sum_{v \in Q_{0}}\left(\sigma_{v} \operatorname{deg} E_{v}-\tau_{v} \operatorname{rank} E_{v}\right),  \tag{278}\\
\mu_{\sigma, \tau}(R) & =\frac{\operatorname{deg}_{\sigma, \tau}(R)}{\sum_{v \in Q_{0}} \sigma_{v} \operatorname{rank} E_{v}} . \tag{279}
\end{align*}
$$

A $Q$-sheaf is stable if one has $\mu_{\sigma, \tau}\left(R^{\prime}\right)<\mu_{\sigma, \tau}(R)$ for any proper subsheaf $R^{\prime}$. A $Q$-sheaf is polystale if it is the direct sum of stable $Q$-sheaf of the same slope.

Theorem 9 ([1, Theorem 3.1])A $Q$-sheaf $R$ with $\operatorname{deg}_{\sigma, \tau}(R)=0$ admits a Hermitian metric satisfying the quiver vortex Eq. (277) if and only if $R$ is $(\sigma, \tau)$-polystable. This Hermitian metric is unique up to a multiplication by a positive constant for each stable summand.

Quasimaps to $\operatorname{Mat}(r, n) / / \mathrm{GL}_{r}$ corresponds to the case when the quiver $Q=$ $(1 \rightarrow 2)$ consists of two vertices and one arrow between them, $M_{1}$ and $M_{2}$ are the structure sheaves, rank $E_{1}=r$, and $E_{2}$ is the trivial bundle of rank $n$.

## 16.6

Note that the map $V \rightarrow \operatorname{End}(V), \phi \mapsto \phi \otimes \phi^{*}$ appearing in (270) is the moment map for the natural action of the unitary group $U(V)$ on $V$. With this in mind, a generalization

$$
\begin{equation*}
* F_{A}+\mu(\Phi)=\tau \operatorname{id}_{E} \tag{280}
\end{equation*}
$$

of the vortex Eq. (270) to the case where one has a Hamiltonian action of a compact group $G$ on a Kähler manifold $X$ is given in [19, 58]. Here $A$ is a connection on a principal $G$-bundle on a curve $C, \Phi$ is a holomorphic section of $P \times_{G} X$, and $\mu: X \rightarrow \mathfrak{g}$ is the moment map. They are used to define invariants of a symplectic manifold with a Hamiltonian group action [19, 20, 59], which are closely related to the Gromov-Witten invariants of the symplectic quotient [30, 71-74]. Reference [18] use wall-crossing in vortex invariants to study quantum cohomology of monotone toric varieties with minimal Chern number greater than or equal to 2 .

## 16.7

Let $X$ be a Kähler manifold with a Hamiltonian action of a compact connected Lie group $G$. We assume that $X$ is either compact or equivariantly convex at infinity with a proper moment map. We fix an invariant inner product to identify $\mathfrak{g}^{\vee}$ with $\mathfrak{g}$, and write the moment map as $\mu: X \rightarrow \mathfrak{g}$.

An affine vortex is a pair $(A, u)$ of a connection $A$ on the principal bundle $P=$ $\mathbb{C} \times G$ and a holomorphic section $u: C \rightarrow P \times{ }_{G} X$ satisfying the vortex equation

$$
\begin{equation*}
* F_{A}+\mu(u)=0 \tag{281}
\end{equation*}
$$

A gauged holomorphic map to $X$ with respect to the complex Lie group $G_{\mathbb{C}}$ acting on $X$ is a map to the quotient stack $\left[X / G_{\mathbb{C}}\right]$. In other words, a gauged holomorphic map from a scheme $C$ to $X$ is a pair $(P, u)$ of a principal $G_{\mathbb{C}}$-bundle $P$ over $C$ and a $G_{\mathbb{C}}$-equivariant holomorphic map $u: P \rightarrow X$.

If the $G_{\mathbb{C}}$-action on $X^{\text {ss }}$ is free, then by [68, Theorem 1.1], there is a natural bijection between the set of affine $K$-vortices with target $X$ up to gauge equivalence and the set of pairs gauged holomorphic maps such that $u(\infty) \in X^{\text {ss }}$. This is an open substack of the set of quasimaps such that $\infty$ is not contained in the base locus.

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[^3]:    ${ }^{1}$ Some authors prefer to use at this place Taylor polynomials of degree $k-1$ instead, see for instance [3] and [4].

[^4]:    ${ }^{2}$ The standard topology language adopted, among many other sources, in [13]

[^5]:    ${ }^{3}$ However, this terminology is not yet definitely settled, as shown in a recent work [10]. The authors of the latter speak just descriptively about 'the Łojasiewicz exponent for the regular separation of closed semialgebraic sets'.

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[^7]:    ${ }^{1}$ This English translation of the profound discovery of van der Waerden in German [59] is due to Schappacher [56].

[^8]:    ${ }^{2}$ The enumerative results in Schubert [51] were mutually verifiable with the results of other geometers (e.g. Salmon, Clebsch, Chasles and Zeuthen) of the same period, hence were already known to be correct at that time.

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[^11]:    ${ }^{1}$ In general, a cocharacter $\varpi: \mathbb{C}^{*} \rightarrow T \subseteq G$ determines a reductive subgroup $G_{\varpi} \subseteq G$, the centralizer of its image; the corresponding parabolic $P_{\varpi}$ is generated by $G_{\varpi}$ together with the Borel.

[^12]:    ${ }^{2}$ When $n=1$, the spaces are 0 -dimensional; when $n=2$, they coincide with type $A$ spaces. For $n=3$, there are coincidences $\mathcal{Q}^{4}=\operatorname{Gr}(2,4)$ and $O G^{+}(3,6)=O G^{-}(3,6)=\mathbb{P}^{3}$. For $n=4$, there are also coincidences $\mathcal{Q}^{6}=O G^{+}(4,8)=O G^{-}(4,8)$. The reader may use these to verify our claims, but beware that the torus actions are usually written differently.

[^13]:    ${ }^{3}$ The indexing most natural to our setup is slightly nonstandard. The rows and columns of $A(x)$ are labelled $n-1, \ldots, 0$ (left to right, top to bottom); similarly, the rows of $B_{I}(x \mid t)$ are labelled $n-1, \ldots, 0$ (top to bottom) and its columns are labelled $i_{1}, \ldots, i_{k}$ (left to right).

[^14]:    ${ }^{4}$ The non-equivariant rim hook rule was discovered by Bertram, Ciocan-Fontanine, and Fulton during the 1996-7 program on quantum cohomology at Institut Mittag-Leffler [3]. Bertiger, Milićević, and Taipale gave an equivariant generalization [2].

[^15]:    ${ }^{5}$ Our indexing of $t$ variables is reversed when compared with that of [39], since we use opposite Schubert classes.

[^16]:    ${ }^{6}$ This "spin basis" is not directly connected to spin representations and the orthogonal Grassmannian, as far as we are aware.

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[^19]:    ${ }^{1}$ The result also holds for $F=F_{q}((t))$.

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[^21]:    ${ }^{1}$ Since $u^{k}=\left[\mathcal{O}_{\mathbb{P}^{n-1-k}}\right] \in K\left(\mathbb{P}^{n-1}\right)$ thus $p_{*}\left(u^{k}\right)=\chi\left(\mathbb{P}^{n-1-k} ; \mathcal{O}\right)=1$.

[^22]:    ${ }^{2}$ This idea appeared already in an early version of [26].

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[^24]:    ${ }^{1}$ For the origin of the name Hessenberg varieties, see [22].

[^25]:    ${ }^{2}$ A paving by affines means a "(complex) cellular decomposition" in algebraic geometry. See [72, Definition 2.1] for the details.

[^26]:    ${ }^{3}$ For $w \in S_{n}$, the associated permutation flag $V_{\mathbf{0}}$ is given by $V_{i}:=\operatorname{span}_{\mathbb{C}}\left\{e_{w(1)}, \ldots, e_{w(i)}\right\}$ where $\left\{e_{1}, \ldots, e_{n}\right\}$ is the standard basis of $\mathbb{C}^{n}$.

[^27]:    ${ }^{4}$ For regular semisimple matrices $S$ and $S^{\prime}$, the associated Hessenberg varieties Hess $(S, h)$ and $\operatorname{Hess}\left(S^{\prime}, h\right)$ with a same Hessenberg function $h$ are diffeomorphic.

[^28]:    ${ }^{5} \mathrm{~A}$ complete graph is a graph in which every pair of distinct vertices is connected by an edge.

[^29]:    ${ }^{6}$ A map $\kappa:[n] \rightarrow \mathbb{N}$ is called a proper coloring of $G$ if $\kappa(i) \neq \kappa(j)$ for all pair of vertices $i$ and $j$ which are connected by an edge.

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