# Schubert functors and Schubert polynomials 

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#### Abstract

We construct a family of functors assigning an $R$-module to a flag of $R$-modules, where $R$ is a commutative ring. As particular instances, we get flagged Schur functors and Schubert functors, the latter family being indexed by permutations. We identify Schubert functors for vexillary permutations with some flagged Schur functors, thus establishing a functorial analogue of a theorem from [6] and [15]. Over an infinite field, we study the trace of a Schubert module, which is a cyclic module over a Borel subgroup $B$, restricted to the maximal torus. The main result of the paper says that this trace is equal to the corresponding Schubert polynomial of Lascoux and Schützenberger [6]. We also investigate filtrations of $B$-modules associated with the Monk formula [10] and transition formula from [8].


## Introduction

In the present paper, we study a certain extension of the notion of polynomial functors defined on the category of $R$-modules, where $R$ is a commutative ring. If $R \supset \mathbb{Q}$, the field of rationals, such polynomial functors are classified by the so-called Schur functors and have been extensively studied ([5], [14], [9], [11], and [1]). Restricted to the vector spaces over a field of characteristic 0, Schur

[^0]functors provide the set of all irreducible polynomial representations of the general linear group, and in this context were investigated already by Schur (cf. [13]). We refer to [12] for a survey concerning Schur functors and some of their applications.

The main notion introduced in the present paper is the notion of a Schubert functor. Schubert functors assign to a given flag of $R$-modules a new $R$-module, where $R$ is a commutative ring ${ }^{1}$. As Schur functors are indexed by partitions, Schubert functors are indexed by permutations; this reflects their relationship with the Schubert calculus on the flag manifold. The last subject, being a part of classical algebraic geometry, goes back to Ehresmann and Bruhat. In the beginning of the seventies, the cohomology ring of the flag variety was extensively studied by Bernstein-Gelfand-Gelfand [2] and Demazure [3] with the help of some differential-like operators defined on it. In [6], Lascoux and Schützenberger have invented a more combinatorial approach; the main notion of this approach is the notion of a Schubert polynomial (see Section 2 for a precise definition). Schubert polynomials are indexed by permutations and generalize Schur polynomials. In his Oberwolfach talk (June 1983), Lascoux suggested the existence of a family of functors generalizing Schur functors and having similar properties as Schubert polynomials. The present paper may be treated as a realization of this program.

The main aim of this paper is to relate Schubert functors and Schubert polynomials. The fundamental relationship between Schur functors and Schur polynomials says that the trace of the restricted action of the group of diagonal matrices on a Schur module is equal to the corresponding Schur polynomial (see [13], or [9] for a modern treatment). The main result of the present paper says that the trace of the restricted action of the group of diagonal matrices on a Schubert module is equal to the corresponding Schubert polynomial. This fact is proved in Section 4. The main idea is based on an analogue of the "transition formula", invented and applied by Lascoux and Schützenberger in [8], and worked out here in the setting of Schubert modules.

Schubert functors are defined and discussed in Section 1. In Section 2, we recall a definition and properties of Schubert polynomials in the setting needed for the purposes of this paper. In particular, we describe a combinatorics standing behind the transition formula. Section 3 is devoted to study of the so-called flagged Schur functors which are some interesting generalizations of Schur functors, and some special cases of Schubert functors. We prove here an analogue of a theorem of Lascoux-Schützenberger [6] and Wachs [15] relating Schubert polynomials for vexillary permutations and flagged Schur polynomials. In Section 5, a filtration associated with the Monk formula [10] is

[^1]discussed.
The main results of the present paper were announced in [4], but the full text has not been published untill now. The present article, modulo some minor modifications, follows the authors' preprint distributed in 1986. We add, at the end, a section on further developments related to the main themes of this paper.

We dedicate this paper to the mathematician who

- introduced in 1977 Schur functors to commutative algebra and combinatorics in the course of his work on the syzygies of determinantal ideals [5],
- invented in 1982 (together with M.-P. Schützenberger) Schubert polynomials ([6], [7]),
- and emphasized, with the same coauthor, the role of vexillary permutations in the study of the symmetric groups ([6], [8]).


## 1 Definition of Schubert functors

Let $R$ be a commutative $\mathbb{Q}$-algebra and let

$$
E .: E_{1} \subset E_{2} \subset \cdots \subset E_{n} \subset \cdots
$$

be a flag of $R$-modules. Assume that $\mathcal{I}=\left[i_{k, l}\right], k, l=1,2, \ldots$, is a matrix consisting of zeros and units and satisfying the following conditions:
$1.1 i_{k, l}=0$ for $k \geq l$.
$1.2 \sum_{l} i_{k, l}$ is finite for every $k$.
1.3 I has a finite number of nonzero rows.

A matrix $\mathcal{I}$ with properties $1.1-1.3$ will be called a shape. We will represent a shape graphically by replacing each unit by " $\times$ " and by omitting zeros on the diagonal and under the diagonal. Moreover, we will omit all the columns which are right to the last nonzero column.

Example 1.4 The matrix


Let $i_{k}:=\sum_{l=1}^{\infty} i_{k, l}, k=1,2, \ldots, \tilde{i}_{l}:=\sum_{k=1}^{\infty} i_{k, l}, l=1,2, \ldots$. We define the module $S_{\mathcal{I}}(E$.$) as the image of the map$

$$
\Phi_{\mathcal{I}}(E .): \bigotimes_{k} S_{i_{k}}\left(E_{k}\right) \xrightarrow{\Delta_{S}} \bigotimes_{k} \bigotimes_{l} S_{i_{k, l}}\left(E_{k}\right) \xrightarrow{m_{\lambda}} \bigotimes_{l} \bigwedge^{\widetilde{\tilde{i}_{l}}} E_{l}
$$

where $\Delta_{S}$ is a product of symmetrizations and $m_{\wedge}$ is a product of multiplications in the corresponding exterior algebras. Observe that $\Phi_{\mathcal{I}}(E$.$) is well$ defined because of properties 1.1-1.3. It is clear that $S_{\mathcal{I}}(-)$ defines, in fact, a functor: if $E$. and $F$. are two flags of $R$-modules and $f: \bigcup_{n} E_{n} \rightarrow \bigcup_{n} F_{n}$ is an $R$-homomorphism such that $f\left(E_{n}\right) \subset F_{n}$, then $f$ induces in a natural way a homomorphism $S_{\mathcal{I}}(E.) \rightarrow S_{\mathcal{I}}(F$.$) .$

For a sequence $I: i_{1} \geq i_{2} \geq \cdots \geq i_{r}$ of nonnegative integers, we define the flagged Schur module $S_{I}(E$.$) as S_{\mathcal{I}}(E$.$) where the shape \mathcal{I}=\left[i_{k, l}\right]$ is defined as follows. Let $c:=\max \left\{s+i_{s}: s=1,2, \ldots, r\right\}$, then

$$
i_{k, l}:= \begin{cases}1 & \text { for } k=1,2, \ldots, r, c-i_{k}+1 \leq l \leq c \\ 0 & \text { otherwise }\end{cases}
$$

Let now $\mu=\mu(1), \mu(2), \mu(3), \ldots$ be a permutation, i.e. bijection of the set of natural numbers which is the identity on the complement of some finite set ${ }^{2}$. For every positive integer $k$, we define

$$
I_{k}(\mu):=\{l: l>k, \mu(k)>\mu(l)\}
$$

which will be called the $k$ th inversion set of $\mu$. The sequence of numbers $i_{k}=\left|I_{k}(\mu)\right|$ (the cardinality of $\left.I_{k}(\mu)\right), k=1,2, \ldots$, is called the code of $\mu$. A code determines the corresponding permutation in a unique way. Note that $\sum i_{k}$ is the length of $\mu$. The index of the permutation $\mu$ is defined to be the number $\sum(k-1) i_{k}$. The shape of $\mu$ is defined to be the matrix

$$
\mathcal{I}_{\mu}=\left[i_{k, l}\right]:=\left[\chi_{k}(l)\right], k, l=1,2, \ldots
$$

where $\chi_{k}$ is the characteristic function of the set $I_{k}(\mu)$.
Example 1.5 For $\mu=5,2,1,6,4,3,7,8, \ldots$, the code of $\mu$ is $(4,1,0,2,1,0, \ldots)$, and the shape $\mathcal{I}_{\mu}$ is

$$
\begin{array}{ccccc}
\times & \times & 0 & \times & \times \\
& \times & 0 & 0 & 0 \\
& 0 & 0 & 0 \\
& & & \times & \times \\
& & & & \times
\end{array}
$$

[^2]We define $S_{\mu}(E$.) - the Schubert module associated with a permutation $\mu$ and a flag $E$., as the module $S_{\mathcal{I}_{\mu}}(E$.); this gives rise to the corresponding Schubert functor $S_{\mu}(-)$.

Remark 1.6 Let $E$. : $E_{1} \subset E_{2} \subset \cdots$ be a flag of vector spaces over a field $K$ of characteristic $0, \operatorname{dim} E_{i}=i$. Let $B$ be the group of linear transformations $f$ of $E:=\bigcup E_{i}$ which preserve $E$., i.e. $f\left(E_{i}\right)=E_{i}$. Observe that the vector spaces used in the definition of $S_{\mathcal{I}}(E$.) are $B$-modules (modules over the group ring $K[B]$ will be called $B$-modules for brevity) and the maps $\Delta_{S}, m_{\wedge}, \Phi_{\mathcal{I}}$ are morphisms of $B$-modules. Moreover, $\otimes_{k} S_{i_{k}}\left(E_{k}\right)$ is a cyclic $B$-module. Therefore $S_{\mathcal{I}}\left(E\right.$.) is also a cyclic $B$-module. Let $\left\{e_{i}: i=1,2, \ldots\right\}$ be a basis of $E$ such that $e_{1}, e_{2}, \ldots, e_{k}$ span $E_{k}$. Then $S_{\mathcal{I}}(E)$ as a $B$-submodule in $\otimes_{l} \wedge^{\widetilde{\Lambda}_{l}} E_{l}$ is generated by the element

$$
C_{\mathcal{I}}:=\bigotimes_{l} e_{k_{1, l}} \wedge e_{k_{2, l}} \wedge \cdots \wedge e_{k_{i, l}}
$$

where $k_{1, l}<k_{2, l}<\cdots<k_{i_{l}, l}$ are precisely those indices for which $i_{k_{r, l}, l}=1$. In particular, if in the $l$ th column of $\mathcal{I}$ the lowest " $\times$ " appears in the $s$ th row, then we can replace $\bigwedge^{\widetilde{i}_{l}} E_{l}$ by $\bigwedge^{\widetilde{i}_{l}} E_{s}$ in the definition of $S_{\mathcal{I}}(E)$. For a given permutation $\mu$, we will write $C_{\mu}$ instead of $C_{\mathcal{I}_{\mu}}$.

Example 1.7 $C_{5,2,1,6,4,3,7, \ldots}=e_{1} \otimes e_{1} \wedge e_{2} \otimes e_{1} \wedge e_{4} \otimes e_{1} \wedge e_{4} \wedge e_{5}$ (an element of $\left.E_{1} \otimes \Lambda^{2} E_{2} \otimes \Lambda^{2} E_{4} \otimes \Lambda^{3} E_{5}\right)$.

Let us record the following immediate consequence of the definition of $S_{\mathcal{I}}(E$.$) .$
Lemma 1.8 Assume that for a shape $\mathcal{I}=\left[i_{k, l}\right], k, l=1,2, \ldots$, there exists $p$ such that $i_{k, l}=0$ for $k \leq p$ and $l>p$. Let $\mathcal{I}^{\prime}$ be a shape $\left[i_{k, l}^{\prime}\right]$ such that $i_{k, l}^{\prime}=i_{k, l}$ for $k \leq p$ and $i_{k, l}^{\prime}=0$ for $k>p$. Let $\mathcal{I}^{\prime \prime}$ be a shape $\left[i_{k, l}^{\prime \prime}\right]$ such that $i_{k, l}^{\prime \prime}=0$ for $k \leq p$ and $i_{k, l}^{\prime \prime}=i_{k, l}$ for $k>p$. Pictorially


Then $S_{\mathcal{I}}(E.) \simeq S_{\mathcal{I}^{\prime}}(E.) \otimes S_{\mathcal{I}^{\prime \prime}}(E).$.

## 2 Review of Schubert polynomials

In this section, we collect some facts on Schubert polynomials, as defined in [6] and [7], useful in other parts of this paper. Let $A:=\mathbb{Z}\left[x_{1}, x_{2}, \ldots\right]$ be a polynomial ring in a countable set of variables. For every $i$, we consider an operator $\partial_{i}: A \rightarrow A$ defined for $f \in A$ by

$$
\partial_{i}(f):=\frac{f\left(x_{1}, \ldots, x_{i}, x_{i+1}, \ldots\right)-f\left(x_{1}, \ldots, x_{i+1}, x_{i}, \ldots\right)}{x_{i}-x_{i+1}} .
$$

For the $i$ th simple reflection $\tau_{i}=1, \ldots, i-1, i+1, i, i+2, \ldots$, we set $\partial_{\tau_{i}}:=\partial_{i}$.
Proposition 2.1 Let $\mu=\sigma_{k} \circ \cdots \circ \sigma_{1}=\pi_{k} \circ \cdots \circ \pi_{1}$ be two reduced decompositions of a permutation $\mu$ into simple reflections. Then

$$
\partial_{\sigma_{k}} \circ \cdots \circ \partial_{\sigma_{1}}(f)=\partial_{\pi_{k}} \circ \cdots \circ \partial_{\pi_{1}}(f)
$$

for every $f \in A$.
Therefore, for a given permutation $\mu$, we can define $\partial_{\mu}$ as $\partial_{\sigma_{k}} \circ \cdots \circ \partial_{\sigma_{1}}$ independently of a reduced decomposition of $\mu$ chosen. For a given $\mu$, let $n$ be a positive integer such that $\mu(k)=k$ for $k>n$. The Schubert polynomial $X_{\mu}$ is defined by

$$
X_{\mu}:=\partial_{\mu^{-1} \circ \omega_{n}}\left(x_{1}^{n-1} x_{2}^{n-2} \cdots x_{n-1}^{1} x_{n}^{0}\right)
$$

where $\omega_{n}$ is the permutation $n, n-1, \ldots, 2,1, n+1, n+2, \ldots$. Observe that the above definition does not depend on a choice of $n$, because

$$
\partial_{n} \circ \cdots \circ \partial_{2} \circ \partial_{1}\left(x_{1}^{n} x_{2}^{n-1} \cdots x_{n}\right)=x_{1}^{n-1} x_{2}^{n-2} \cdots x_{n-1} .
$$

Schubert polynomials satisfy the following properties:
2.2 The Schubert polynomial $X_{\mu}$ is symmetric with respect to $x_{k}$ and $x_{k+1}$ iff $\mu(k)<\mu(k+1)$ (or equivalently $i_{k} \leq i_{k+1}$ ).
2.3 If $\mu(1)<\mu(2)<\cdots<\mu(k)>\mu(k+1)<\mu(k+2)<\cdots$ (or equivalently $\left.i_{1} \leq i_{2} \leq \cdots \leq i_{k}, 0=i_{k+1}=i_{k+2}=\cdots\right)$, then $X_{\mu}$ equals the Schur polynomial $s_{i_{k}, \ldots, i_{2}, i_{1}}\left(x_{1}, \ldots, x_{k}\right)$ (cf. [9] for this last notion).

Pictorially, the shape of $\mu$ can be displayed as

2.4 If $i_{1} \geq i_{2} \geq \cdots$, then $X_{\mu}=x_{1}^{i_{1}} x_{2}^{i_{2}} \cdots$ is a monomial.

If the sets $I_{k}(\mu)$ ordered by inclusion form a chain, we call $\mu$ a vexillary permutation (see [6] and [8] for more information on vexillary permutations).

Example $2.53,6,5,1,4,2,7, \ldots$ is a vexillary permutation contrary to $6,3,5,1,2,4,7, \ldots$.

Recall the notion of a flagged Schur function $s_{i_{1}, \ldots, i_{k}}\left(b_{1}, \ldots, b_{k}\right)$ where $i_{1} \geq$ $\cdots \geq i_{k}$ and $0<b_{1} \leq \cdots \leq b_{k}$ are two sequences of nonnegative integers. One sets

$$
s_{i_{1}, \ldots, i_{k}}\left(b_{1}, \ldots, b_{k}\right):=\operatorname{det}\left(s_{i_{p}-p+q}\left(x_{1}, \ldots, x_{b_{p}}\right)\right)_{1 \leq p, q \leq k}
$$

where $s_{l}$ is the complete homogeneous symmetric function of degree $l$ (for an alternative definition in terms of standard tableaux, see [15]). For a given sequence of nonnegative integers $I=\left(i_{1}, \ldots, i_{n}\right)$, we define $I^{\leq}$(resp. $I^{\geq}$) as the increasing (resp. decreasing) reordering of $I$. The following result stems from [6] and [15]:

Theorem 2.6 Let $\mu$ be a vexillary permutation with the code $\left(i_{1}, i_{2}, \ldots\right)$ and assume $i_{n} \neq 0, i_{r}=0$ for $r>n$. Then

$$
X_{\mu}=s_{\left(i_{1}, \ldots, i_{n}\right) \geq}\left(\min I_{1}(\mu)-1, \ldots, \min I_{n}(\mu)-1\right)^{\leq} .
$$

We finish this section with some facts on multiplication of Schubert polynomials.
2.7 The Monk formula for multiplication by $X_{\tau_{k}}([10],[7])$ :

$$
X_{\mu} \cdot\left(x_{1}+\cdots+x_{k}\right)=\sum X_{\mu \circ \tau_{p, q}},
$$

where $\tau_{p, q}, p<q$, denotes the permutation

$$
1, \ldots, p-1, q, p+1, \ldots, q-1, p, q+1, \ldots
$$

and the summation ranges over all $p, q$ such that $p \leq k, q>k$ and $l\left(\mu \circ \tau_{p, q}\right)=$ $l(\mu)+1$.

This last requirement can be restated as follows: for each $i \in] p, q[, \mu(i) \notin$ $[\mu(p), \mu(q)]$. Graphically, the shapes of permutations $\mu \circ \tau_{p, q}$ can be obtained from the shape of $\mu$ in the following way. Let us restrict our attention to the places occupied by the 0's in the marked rectangle with $k$ rows


Fix one such place. Then consider the hook having its head in the fixed place and the remaining boxes in the left hand side of the row, and the bottom part of the column occupied by 0 .


If in the hook the number of the 0 's is equal to the number of the $\times$ 's plus one, then this place gives a contribution to the right-hand side of the Monk formula. The shape of this summand is obtained from the shape of $\mu$ by replacing our 0 by $\times$ and by exchanging parts of appriopriate columns and rows as indicated in the following pictures:


Example 2.8 $X_{246315879 \ldots} \cdot\left(x_{1}+x_{2}\right)=X_{346215879 \ldots}+X_{264315879 \ldots}+X_{256314879 \ldots}$. Note that multiplication by $x_{1}$ becomes

$$
X_{\mu} \cdot x_{1}=\sum X_{\mu \circ \tau_{1, q}},
$$

where $q$ is such that if $i \in] 1, q[$ then $\mu(i) \notin[\mu(1), \mu(q)]$. In the above graphical description, we consider only the places occupied by the 0's in the first row of the shape of $\mu$.

From the Monk formula, we obtain:
2.9 One has $X_{\mu} \cdot x_{k}=\sum_{j} \operatorname{sign}(j-k) \cdot X_{\mu \circ \tau_{j, k}}$ where the summation ranges over all $j$ such that

$$
\begin{aligned}
& 1^{\circ} \operatorname{sign}(j-k)=\operatorname{sign}(\mu(j)-\mu(k)), \\
& \left.2^{\circ} \text { for each } i \in\right] j, k[, \mu(i) \notin[\mu(j), \mu(k)] .
\end{aligned}
$$

## Example 2.10

$$
\begin{aligned}
& X_{31425 \ldots} \cdot x_{1}=X_{41325 \ldots} \\
& X_{31425 \ldots \ldots} \cdot x_{2}=X_{34125 \ldots}+X_{32415 \ldots} \\
& X_{31425 \ldots} \cdot x_{3}=-X_{41325 \ldots}+X_{315246 \ldots} \\
& X_{31425 \ldots} \cdot x_{4}=-X_{32415 \ldots}+X_{314526 \ldots}
\end{aligned}
$$

For the purposes of Section 4, the following special case of 2.9 will be of particular importance.

Transition Formula 2.11 ([8]) Let $(j, s)$ be a pair of positive integers such that

```
\(1^{\circ} j<s\) and \(\mu(j)>\mu(s)\),
\(2^{\circ}\) for every \(i\) in \(] j, s[, \mu(i)\) is outside of \([\mu(s), \mu(j)]\),
\(3^{\circ}\) for every \(r>j\), if \(\mu(s)<\mu(r)\) then there exists \(\left.i \in\right] j\), \(r[\) such that \(\mu(i) \in\)
    \([\mu(s), \mu(r)]\),
```

then

$$
X_{\mu}=X_{\lambda} \cdot x_{j}+\sum X_{\psi_{t}}
$$

where $\lambda=\mu \circ \tau_{j, s}, \psi_{t}=\mu \circ \tau_{j, s} \circ \tau_{k t, j}$ and the sum ranges over the set of numbers $k_{t}$ satisfying

$$
\begin{aligned}
& 4^{\circ} \quad k_{t}<j \text { and } \mu\left(k_{t}\right)<\mu(s) \\
& \left.5^{\circ} \text { if } i \in\right] k_{t}, j\left[\text { then } \mu(i) \notin\left[\mu\left(k_{t}\right), \mu(s)\right] .\right.
\end{aligned}
$$

Observe that if $(j, s)$ is the maximal pair (in the lexicographical order) satisfying condition $1^{\circ}$ then conditions $2^{\circ}-3^{\circ}$ are also satisfied. A transition determined by this pair will be called maximal. In particular, for every nontrivial permutation there exists at least one transition.

Example $2.12 \mu=5,2,1,8,6,3,4,7,9, \ldots$,

$$
\begin{aligned}
& X_{\mu}=X_{521843679 \ldots} \cdot x_{5}+X_{524813679 \ldots}+X_{541823679 \ldots} \\
& X_{\mu}=X_{521763489 \ldots} \cdot x_{4}+X_{527163489 \ldots}+X_{571263489 \ldots}+X_{721563489 \ldots} \\
& X_{\mu}=X_{512864379 \ldots} \cdot x_{2},
\end{aligned}
$$

and the first equation gives the maximal transition.

A graphical interpretation of the transition formula will be given in detail in the course of the proof of Theorem 4.1. We finish this section with the following analog of Lemma 1.8 (see also [8]).

Lemma 2.13 Assume that a shape $\mathcal{I}_{\mu}=\left[i_{k, l}\right]=1,2, \ldots$ satisfies the following condition: there exists $p$ such that $i_{k, l}=0$ for $k \leq p$ and $l>p$. Let $\mathcal{I}^{\prime}=\left[i_{k, l}^{\prime}\right]$ be a shape defined as $i_{k, l}^{\prime}=i_{k, l}$ for $k \leq p$ and $i_{k, l}^{\prime}=0$ for $k>p$. Similarly let $\mathcal{I}^{\prime \prime}=\left[i_{k, l}^{\prime \prime}\right]$ be a shape defined as $i_{k, l}^{\prime \prime}=0$ for $k \leq p$ and $i_{k, l}^{\prime \prime}=i_{k, l}$ for $k>p$. Then there exist permutations $\mu^{\prime}$ and $\mu^{\prime \prime}$ such that $\mathcal{I}_{\mu^{\prime}}=\mathcal{I}^{\prime}, \mathcal{I}_{\mu^{\prime \prime}}=\mathcal{I}^{\prime \prime}$, and the following equality holds

$$
X_{\mu}=X_{\mu^{\prime}} \cdot X_{\mu^{\prime \prime}}
$$

PROOF. Consider a transition for $\mu^{\prime}$ :

$$
X_{\mu^{\prime}}=X_{\lambda^{\prime}} \cdot x_{j}+\sum X_{\psi_{t}^{\prime}}
$$

Then

$$
X_{\mu^{\prime}} \cdot X_{\mu^{\prime \prime}}=\left(X_{\lambda^{\prime}} \cdot x_{j}+\sum X_{\psi_{t}^{\prime}}\right) \cdot X_{\mu^{\prime \prime}}=X_{\lambda^{\prime}} \cdot X_{\mu^{\prime \prime}} \cdot x_{j}+\sum X_{\psi_{t}^{\prime}} \cdot X_{\mu^{\prime \prime}}
$$

Let $\lambda$ be the permutation corresponding to the shape $\left[i_{k, l}^{\lambda^{\prime}}+i_{k, l}^{\mu^{\prime \prime}}\right]$ and let $\psi_{t}$ be the permutation corresponding to the shape $\left[i_{k, l}^{\psi^{\prime}}+i_{k, l}^{\mu^{\prime \prime}}\right]$. Then by induction on the length of a permutation we get

$$
X_{\mu^{\prime}} \cdot X_{\mu^{\prime \prime}}=X_{\lambda} \cdot x_{j}+\sum X_{\psi_{t}}
$$

It is easy to check that the right-hand side of this equality is a transition for $X_{\mu}$. This proves that $X_{\mu^{\prime}} \cdot X_{\mu^{\prime \prime}}=X_{\mu}$ as needed.

## 3 Schubert functors for vexillary permutations

In this section, we will deal with vexillary permutations. We will prove an analog of Theorem 2.6 in the context of Schubert functors and flagged Schur functors.

A flag $E$. of free $R$-modules will be called a splitting flag provided $E_{1} \simeq R$ and the $i$ th inclusion in the flag is given by $E_{i} \hookrightarrow E_{i} \oplus R \simeq E_{i+1}$.

Theorem 3.1 Let E. be a splitting flag of free $R$-modules with $\operatorname{rank} E_{i}=i$. Assume that $\mu$ is a vexillary permutation and $I_{1}, I_{2}, \ldots, I_{n}$ are all nonempty sets among the sets $I_{s}(\mu), s=1,2, \ldots$, counted with multiplicities. Then $S_{\mu}(E$.$) is the flagged Schur module associated with the sequence \left(\left|I_{k}\right|\right)^{\geq}, k=$
$1,2, \ldots, n$, and with the (sub)flag of modules of ranks $\left(\min \left(I_{k}\right)-1\right)^{\leq}, k=$ $1,2, \ldots, n$.

PROOF. Each permutation $\mu=\mu(1), \mu(2), \mu(3), \ldots$ has a unique decomposition into maximal increasing sequences

$$
\mu=a_{1}<a_{2}<\cdots<a_{k}>b_{1}<b_{2}<\cdots<b_{l}>c_{1}<\cdots .
$$

With a given permutation $\mu$ we associate the permutation $\mu^{\prime}$ as follows.
$1^{\circ}$ From the set $\left\{a_{1}, a_{2}, \ldots, a_{k}\right\}$ we subtract all the numbers which are larger than $b_{1}$. Assume that their cardinality is equal to $s$.
$2^{\circ}$ Let $\sigma:\{1,2, \ldots\} \backslash\left\{a_{k-s+1}, \ldots, a_{k}\right\} \rightarrow\{s+1, s+2, \ldots\}$ be the unique bijection preserving the order.

Define $\mu^{\prime}$ as follows

$$
\mu^{\prime}:=1,2, \ldots, s, \sigma\left(a_{1}\right), \ldots, \sigma\left(a_{k-s}\right), \sigma\left(b_{1}\right), \ldots, \sigma\left(b_{l}\right), \ldots
$$

Example 3.2 If $\mu=2,6,4,5,1,3,7, \ldots$, then $\mu^{\prime}=1,3,5,6,2,4,7, \ldots$.

$$
\begin{aligned}
& \begin{array}{llllllllll}
0 & 0 & 0 & \times & 0 & 0 & 0 & 0 & 0 & 0
\end{array} \\
& \times \times \times \times \\
& \begin{array}{ccccccc}
\mathcal{I}_{\mu}= & \times & \times & \mathcal{I}_{\mu^{\prime}}= & 0 & \times & \times \\
& \times & \times & & & \times & \times \\
& & 0 & & & & \\
& & & & &
\end{array}
\end{aligned}
$$

Observe that the shape of $\mu^{\prime}$ is obtained from $\mathcal{I}_{\mu}$ by pushing down the first $k-s$ rows by $s$ rows and introducing the first $s$ rows with the 0 's only. In particular, $l\left(\mu^{\prime}\right)<l(\mu)$, and $\mu^{\prime}$ is vexillary if $\mu$ has this property. By definition, we have the following decomposition of $\Phi_{\mu}=\Phi_{\mu}(E$.$) :$


Observe that the sequence $\left(I_{s}(\mu)\right), s=1, \ldots, k$, is nondecreasing because $a_{1}<\cdots<a_{k}$. By a standard basis theorem for Schur functors (see [14], [1]),
we can replace our initial flag by an arbitrary flag between $E_{1} \subset \cdots \subset E_{k} \subset$ $E_{k+1} \subset \cdots$ and $E_{k} \subset \cdots \subset E_{k} \subset E_{k+1} \subset \cdots$ without changing the image. Let us use the flag

$$
E_{s+1} \subset \cdots \subset E_{s+(k-s)} \subset \underbrace{E_{k} \subset \cdots \subset E_{k}}_{s} \subset E_{k+1} \subset E_{k+2} \subset \cdots
$$

The image of $\Phi_{\mu}$ is now equal to the image of the following map $\Phi_{\mu}^{\prime}$ :

$$
\begin{aligned}
& \bigotimes_{p=k-s+1}^{k} S_{i_{p}}\left(E_{k}\right) \xrightarrow{\Delta_{S}^{\prime \prime \prime}} \bigotimes_{p, q} S_{i_{p, q}}\left(E_{k}\right) \xrightarrow{m_{\curlywedge}^{\prime \prime \prime}} \bigotimes_{q} \bigwedge^{\left(\sum_{p=k-s+1}^{k} i_{p, q}\right)}\left(E_{k}\right) \\
& \otimes \\
& \otimes \\
& \left(\begin{array}{c}
\begin{array}{c}
k-s \\
p=1 \\
Q_{i_{p}}
\end{array}\left(E_{s+p}\right) \\
\otimes \\
\underset{p>k}{ } S_{i_{p}}\left(E_{p}\right)
\end{array}\right) \xrightarrow{\Delta_{S}^{\mathrm{IV}}} \bigotimes_{p, q} S_{i_{p, q}}\left(E_{p}\right) \xrightarrow{m_{\wedge}^{\mathrm{IV}}} \underset{q}{ } \bigwedge^{\left(\sum_{p=1}^{k-s} i_{p, q}+\sum_{p>k} i_{p, q}\right)}\left(E_{q}\right)
\end{aligned}
$$

The composition $m_{\wedge}^{\mathrm{IV}} \circ \Delta_{S}^{\mathrm{IV}}$ is nothing but $\Phi_{\mu^{\prime}}(E$.$) . Therefore, by induction on$ $l(\mu)$, we can replace $m_{\wedge}^{\mathrm{IV}} \circ \Delta_{S}^{\mathrm{IV}}$ by the corresponding map defining the flagged Schur module associated with the sequence and the (sub)flag of modules described by nonempty inversion sets of $\mu^{\prime}$. After this exchange, we get in the above diagram the map defining the flagged Schur module associated with the desired sequence and flag because $k+1=\min I_{j}(\mu)$ for $j=k-s+1, \ldots, k$, and these sets are the only inversion sets of $\mu$ that are not inversion sets for $\mu^{\prime}$.

This completes the proof of the theorem.

## 4 A character formula relating Schubert modules and Schubert polynomials

In this section, $E .: E_{1} \subset E_{2} \subset \cdots$ will denote a flag of vector spaces over a field $K$ of characteristic zero, $\operatorname{dim} E_{i}=i$. We choose a basis for $E=\cup E_{i}$ as in Remark 1.6. Then the group $B$ of linear automorphisms of $E$ preserving $E$. may be identified with the group of upper triangular matrices. Let $T$ be the subgroup of diagonal matrices in $B$

$$
\left(\begin{array}{cccc}
x_{1} & & & \\
& x_{2} & & 0 \\
& & x_{3} & \\
& 0 & & \ddots
\end{array}\right)
$$

Consider the action of $T$ on $S_{\mu}(E$.$) induced from the canonical action of B$ by restriction. The main result of this paper is:

Theorem 4.1 The trace of the action of $T$ on $S_{\mu}(E$.$) is equal to X_{\mu}$.
The proof is divided into several steps.
Step 1. Let

$$
X_{\mu}=X_{\lambda} \cdot x_{j}+\sum_{t=1}^{k} X_{\psi_{t}}
$$

be the maximal transition for $\mu$. We claim that there exists a filtration of $B$-modules

$$
0=\mathcal{F}_{0} \subset \mathcal{F}_{1} \subset \cdots \subset \mathcal{F}_{k} \subset \mathcal{F}=S_{\mu}(E .)
$$

together with surjections $\mathcal{F} / \mathcal{F}_{k} \rightarrow S_{\lambda}(E.) \otimes E_{j} / E_{j-1}$ and $\mathcal{F}_{t} / \mathcal{F}_{t-1} \rightarrow S_{\psi_{t}}(E$. for every $t=1, \ldots, k$.

PROOF. Assume that $\lambda$ is obtained from $\mu$ by exchanging $\mu(j)$ and $\mu(s)$ as described in 2.7. Recall that $\mu(j)>\mu(s)$ and that the sequence $\mu(j+1)$, $\mu(j+2), \ldots, \mu(s), \ldots$ is increasing. Therefore the sth row of the shape of $\mu$ has the $\times$ 's in columns with numbers $j+1, j+2, \ldots, s$. Moreover, the shape $\mathcal{I}_{\lambda}$ is obtained from the shape $\mathcal{I}_{\mu}$ by omitting the last $\triangle \times$ in the $j$ th row and by exchanging the $j$ th and $s$ th columns.


Let $\Delta_{\wedge}: \bigwedge^{\widetilde{i}_{s}} E_{s} \rightarrow \bigwedge^{\widetilde{i}_{s}-1} E_{s} \otimes E_{s}$ be the diagonalization and let $p: E_{s} \rightarrow$ $E_{s} / E_{j-1}$ be the projection. Denote by $\varphi$ the following composition of $B$-module homomorphisms:

$$
\begin{gathered}
\bigwedge^{\tilde{i}_{1}} E_{1} \otimes \cdots \otimes \bigwedge^{\widetilde{c}_{j}} E_{j} \otimes \cdots \otimes \bigwedge^{\widetilde{i}_{s}} E_{s} \otimes \cdots \\
\bigwedge^{\tilde{i}_{1}} E_{1} \otimes \cdots \otimes \bigwedge^{\widetilde{彳}_{j}} E_{j} \otimes \cdots \otimes 1 \otimes \cdots \otimes \Delta \wedge \otimes \\
\left(\bigwedge^{\tilde{i}_{1}} E_{1} \otimes \cdots \otimes \bigwedge^{\widetilde{i}_{s}-1} E_{s} \otimes E_{s} \otimes \cdots\right. \\
\left.{\widetilde{\tilde{i}_{s}}-1}^{\tilde{\tau}_{s}} E_{s} \otimes \cdots \otimes \bigwedge^{\widetilde{i}_{j}} E_{j} \otimes \cdots\right) \otimes E_{s}
\end{gathered}
$$

$$
\left(\bigwedge^{\widetilde{i}_{1}} E_{1} \otimes \cdots \otimes \bigwedge^{{\widetilde{\tilde{c}_{s}}-1}^{\downarrow}(1 \otimes \cdots \otimes 1 \otimes \cdots \otimes 1 \otimes \cdots) \otimes p} E_{s} \otimes \cdots \otimes \bigwedge^{\widetilde{i}_{j}} E_{j} \otimes \cdots\right) \otimes E_{s} / E_{j-1}
$$

Observe that the image of $\varphi$ restricted to $S_{\mu}(E$.$) is equal to S_{\lambda}(E.) \otimes E_{j} / E_{j-1}$. Indeed, we have $\varphi\left(C_{\mu}\right)=C_{\lambda} \otimes e_{j}$. Assume now that $\psi_{t}$ is obtained from $\mu$ by the cyclic exchange of $\mu\left(k_{t}\right), \mu(j), \mu(s)$ as described in 2.11. Define the shape $\mathcal{I}_{t}$ by exchanging the parts of the $k_{t}$ th and $j$ th rows corresponding to the columns $j+1, j+2, \ldots$.


The shape of $\psi_{t}$ can be obtained from the shape $\mathcal{I}_{t}$ by the following exchange of column-segments.
(a) The $j$ th column goes into the place of the $s$ th column.
(b) The upper $k_{t}-1$ rows of the $j$ th column of $\mathcal{I}_{\psi_{t}}$ are formed by the corresponding rows of the $k_{t}$ th column of $\mathcal{I}_{t}$; the remaining part of the $j$ th column of the shape of $\psi_{t}$ is composed of the corresponding rows of the sth column.
(c) The upper $k_{t}-1$ rows of the $k_{t}$ th column are formed by the corresponding rows of the $s$ th column.

We define the $t$ th piece of the filtration by

$$
\mathcal{F}_{t}:=\sum_{r<t} S_{\mathcal{I}_{r}}(E .) .
$$

Let $p=\sum_{l<k_{t}} i_{l, s}$ and $q=\sum_{l>k_{t}} i_{l, s}$. We shall now define maps from $\mathcal{F}_{t}$ to $S_{\psi_{t}}(E$.). Consider the following map

$$
\begin{gathered}
\bigwedge^{i} E_{k_{t}} \otimes \bigwedge^{p+q} E_{j} \\
\bigwedge^{i} E_{k_{t}} \otimes \bigwedge^{p} E_{j} \otimes \bigwedge^{q} E_{j}
\end{gathered}
$$

$$
\begin{gathered}
\bigwedge^{p} E_{j} \otimes \bigwedge^{i} E_{k_{t}} \otimes \bigwedge^{q} E_{j} \\
\downarrow 18 m \wedge \\
\bigwedge^{p} E_{j} \otimes \bigwedge^{i+q} E_{j}
\end{gathered}
$$

where $\Delta_{\wedge}$ is the corresponding diagonalization, $m_{\wedge}$ stands for multiplication in the exterior algebra and $i:=\widetilde{i}_{k_{t}}$. This map tensored by identities corresponding to the remaining columns gives a map from the product of exterior powers associated with $\mathcal{I}_{t}$ (and also with all the shapes $\mathcal{I}_{r}, r<t$ ) to the product of exterior powers associated with $\mathcal{I}_{\psi_{t}}$ (note that we can use exterior powers of $E_{j}$ instead of $E_{s}$ because all the rows below the $j$ th one are occupied by zeros). Denote the restriction of this map to $\mathcal{F}_{t}$ by $\varphi_{t}$. It is easy to see that $\varphi_{t}\left(C_{\mathcal{I}_{t}}\right)=C_{\psi_{t}}$ and $\varphi_{t}\left(C_{\mathcal{I}_{r}}\right)=0$ for $r<t$. Therefore $\mathcal{F}_{t-1} \subset \operatorname{Ker} \varphi_{t}$. This proves that for every $t=1,2, \ldots, k, \mathcal{F}_{t} / \mathcal{F}_{t-1}$ surjects onto $S_{\psi_{t}}(E$.$) .$

Step 2. Define $d_{\mu}:=\operatorname{dim}_{K} S_{\mu}\left(E\right.$.) and $z_{\mu}$ as the value of $X_{\mu}$ for the specialization $x_{i}=1, i=1,2, \ldots$ We claim that $d_{\mu} \geq z_{\mu}$.

PROOF. Consider the maximal transition for $\mu$ :

$$
X_{\mu}=X_{\lambda} \cdot x_{j}+\sum X_{\psi_{t}},
$$

It follows from the description in 2.11 that

1) $l(\lambda)=l(\mu)-1$,
2) for every $t$ the index of $\psi_{t}$ is strictly smaller than the index of $\mu$.

Our claim now follows from the properties of the filtration of $S_{\mu}(E$.$) proved$ in Step 1 by (double) induction on the length and on the index of $\mu$.

Step 3. We claim that $\operatorname{Ker} \varphi_{t}=\mathcal{F}_{t-1}$.

PROOF. We will say that a permutation $\mu$ (or alternatively its code) is suitable if $d_{\mu}=z_{\mu}$. We will prove by induction on the length of a permutation that every permutation is suitable. Then the claim will follow by comparing the dimensions of the both sides of the above equality. It suffices to show that if a code $\left(i_{1}, \ldots, i_{p}, 0,0, \ldots\right)$ is suitable then the code $\left(i_{1}, \ldots, i_{p}+1,0,0, \ldots\right)$ is suitable, and the code

$$
(i_{1}, \ldots, i_{p}, \underbrace{0, \ldots, 0}_{l}, 1,0, \ldots)
$$

is suitable for every $l \geq 0$. First observe that the code of this type is suitable for $l \gg 0$. Indeed, let $n$ be the smallest number such that $\mu(n)=n, \mu(n+1)=$ $n+1, \ldots$. Then the permutation

$$
\mu_{r}=\mu(1), \mu(2), \ldots, \mu(n-1), n, \ldots, n+r-1, n+r+1, n+r, \ldots
$$

(where $r \geq 0$ ) has the shape


This shape satisfies the assumptions of Lemmas 1.8 and 2.13. These lemmas and the fact that $\mu$ is suitable imply that $\mu_{r}$ is suitable. To decrease $l$ we will use 2.11. Observe that if $\mu$ is suitable then $\lambda$ and $\psi_{t}$ 's appearing in the right-hand side of 2.11 are also suitable. Indeed, we have

$$
z_{\lambda}+\sum z_{\psi_{t}}=z_{\mu}=d_{\mu} \geq d_{\lambda}+\sum d_{\psi_{t}},
$$

and thus

$$
\left(d_{\lambda}-z_{\lambda}\right)+\sum\left(d_{\psi_{t}}-z_{\psi_{t}}\right) \leq 0
$$

But by Step 2, we know that all the summands are nonnegative so they must be zero. Let us start with the permutation

$$
\mu_{0}=\mu(1), \mu(2), \ldots, \mu(n-1), n+1, n, n+2, \ldots .
$$

The maximal transition for $\mu_{0}$ contains in the right-hand side

$$
\mu_{-1}=\mu(1), \mu(2), \ldots, \mu(n-2), n, \mu(n-1), n+1, \ldots
$$

The transition for $\mu_{-1}$ contains

$$
\mu_{-2}=\mu(1), \mu(2), \ldots, \mu(n-1), \mu(n-2), \ldots
$$

By repeating this consideration we finally get

$$
\mu(1), \mu(2), \ldots, \mu(p), \mu(p+2), \mu(p+1), \mu(p+3), \ldots
$$

In this way, we get permutations with the codes of the form

$$
(i_{1}, \ldots, i_{p}, \underbrace{0, \ldots, 0}_{l}, 1,0, \ldots), \quad l \geq 0 .
$$

Moreover, if $r$ is the smallest number bigger than $p$ such that $\mu(p)>\mu(r)$, then the transition for the permutation

$$
\mu(1), \ldots, \mu(p), \ldots, \mu(r+1), \mu(r), \ldots
$$

contains the permutation

$$
\mu(1), \ldots, \mu(p-1), \mu(r), \mu(p+1), \ldots, \mu(r-1), \mu(p), \mu(r+1), \ldots
$$

with the code $\left(i_{1}, \ldots, i_{p}+1,0, \ldots\right)$. This proves our claim.
Step 4. From the above considerations, we get a filtration of $B$-modules

$$
0=\mathcal{F}_{0} \subset \mathcal{F}_{1} \subset \cdots \subset \mathcal{F}_{k} \subset \mathcal{F}=S_{\mu}(E .)
$$

such that

$$
\mathcal{F} / \mathcal{F}_{k} \simeq S_{\lambda}(E .) \otimes E_{j} / E_{j-1}, \quad \mathcal{F}_{t} / \mathcal{F}_{t-1} \simeq S_{\psi_{t}}(E .)
$$

for every $t=1, \ldots, k$. In particular, we obtain an isomorphism of $T$-modules

$$
S_{\mu}(E .) \simeq S_{\lambda}(E .) \otimes E_{j} / E_{j-1} \oplus \bigoplus_{t=1}^{k} S_{\psi_{t}}(E .)
$$

By comparing this with the transition formula for Schubert polynomials (see 2.7), the assertion of Theorem 4.1 now follows by double induction on the length and the index of $\mu$, as in Step 2.

The proof of Theorem 4.1 is finished.

## 5 A filtration associated with the Monk formula

Recall (see 2.7) that

$$
X_{\mu} \cdot x_{1}=\sum_{q} X_{\mu \circ \tau_{1, q}}
$$

where $q$ is such that for $i \in] 1, q[, \mu(i)$ is outside of $[\mu(1), \mu(q)]$. Denote the set of all such $q$ 's by $\left\{q_{1}<\cdots<q_{k}\right\}$.

Proposition 5.1 For a given flag E. of vector spaces over a field $K$ of characteristic zero, $\operatorname{dim} E_{i}=i$, there exists a filtration

$$
0=\mathcal{F}_{0} \subset \mathcal{F}_{1} \subset \cdots \subset \mathcal{F}_{k}=S_{\mu}(E .) \otimes E_{1}
$$

of $B$-modules such that $\mathcal{F}_{t} / \mathcal{F}_{t-1} \simeq S_{\mu \circ \tau_{1, q_{t}}}(E$.$) for t=1, \ldots, k$.

PROOF. Fix $t=1, \ldots, k$. Set $\mu_{t}:=\mu \circ \tau_{1, q_{t}}$ and let $\mathcal{I}=\left[i_{p, q}\right]$ (resp. $\left.\mathcal{I}^{t}=\left[i_{p, q}^{t}\right]\right)$ be the shape of $\mu$ (resp. $\mu_{t}$ ). We define an element $C_{t}$ in $\otimes_{l} \wedge^{\widetilde{i}_{l}} E_{l}$ as follows. Let $\left\{e_{i}: i=1,2, \ldots\right\}$ be a basis of $E=\cup E_{i}$ such that $e_{1}, \ldots, e_{k}$ span the vector space $E_{k}$. Set $C_{t}:=\otimes_{l} e_{k_{1}} \wedge e_{k_{2}} \wedge \cdots \wedge e_{k_{i}}$, where for $l \neq q_{t}$ (resp. $l=q_{t}$ ) $k_{r}$ 's are precisely those indices (taken in ascending order) for which $i_{k_{r}, l}^{t}=1$ (resp. those increasing indices different from 1 for which $i_{k_{r}, l}^{t}=1$ ); note that $i_{1, q_{t}}^{t}=1$. Let $\mathcal{F}_{t}$ be the $B$-module generated by $C_{1} \otimes e_{1}, \ldots, C_{t} \otimes e_{1}$; of course, $\mathcal{F}_{k}=S_{\mu}(E.) \otimes E_{1}$. We have a map

$$
\varphi_{t}: \bigotimes_{l} \bigwedge^{{\tilde{\tilde{c}_{l}}}_{l} E_{l} \otimes E_{1} \rightarrow \bigotimes_{l} \bigwedge^{\tilde{\tau}_{t}^{t}} E_{l}, ~}
$$

determined by the exterior multiplication

$$
\bigwedge^{{\widetilde{q_{q}}}} E_{q_{t}} \otimes E_{1} \rightarrow \bigwedge^{{\widetilde{\iota_{q}}}_{t}^{t}} E_{q_{t}} .
$$

This map induces a surjection because $\varphi_{t}\left(C_{t} \otimes e_{1}\right)=C_{\mu_{t}}$. Moreover, it is easy to see that $\mathcal{F}_{t-1} \subset \operatorname{Ker} \varphi_{t}$. By Theorem 4.1, we know that $d_{\mu}=z_{\mu}$ for every permutation. Thus

$$
d_{\mu} \geq \sum_{t=1}^{k} d_{\mu_{t}}=\sum_{t=1}^{k} z_{\mu_{t}}=z_{\mu}
$$

(the inequality follows from the above filtration). Therefore $d_{\mu}=\sum d_{\mu_{t}}$, and this shows that $\mathcal{F}_{t} / \mathcal{F}_{t-1} \simeq S_{\mu_{t}}(E$.$) for t=1, \ldots, k$. The proof is complete.

Remark 5.2 (i) A similar filtration can be associated with the Monk formula for multiplication by $\sum_{k=1}^{n} x_{k}$ when $\mu(r)=r$ for $r>n$. We omit the details. Does there exist an analogous filtration associated with the Monk formula for every $\mu$ ?
(ii) Observe that the first part of the proof of Proposition 5.1 gives us another justification of the inequality $d_{\mu} \geq z_{\mu}$. It suffices to use descending induction on the length of a permutation and the Monk formula for Schubert polynomials.

Remark 5.3 Suppose, as we did in [4], that $R$ is any commutative ring. Replace the symmetric powers (resp. symmetrizations) by divided powers (resp. diagonalizations in the algebra of divided powers) in the definition of a module $S_{\mathcal{I}}(E$.$) associated with flag E$. and shape $\mathcal{I}$. The proofs of Theorem 3.1, Theorem 4.1, and Proposition 5.1 suitably adapted to this modified definition, show that the so obtained flagged Schur functors and Schubert functors have the following properties:

- For any vexillary permutation and any splitting flag (in the sense of Section 3), Theorem 3.1 remains true; consequently these flagged Schur functors are universally free, following the terminology of [1].
- If $K$ is an infinite field (of arbitrary characteristic), then Theorem 4.1 remains valid.
- If $K$ is an infinite field, then Proposition 5.1 also holds true.


## 6 Further developments

The ideas and results of the present paper (in particular, Theorem 4.1) were developed mainly by V. Reiner \& M. Shimozono, and P. Magyar.

A detailed study of different aspects of flagged Schur modules was made by the former authors in the following papers: Key polynomials and a flagged Littlewood-Richardson rule [J. Combin. Theory Ser. A 70 (1995), 107-143], Percentage-avoiding, northwest shapes and peelable tableaux [J. Combin. Theory Ser. A 82 (1998), 1-73], Flagged Weyl modules for two column shapes [J. Pure Appl. Algebra 141 (1999), 59-100]. On another hand, in the paper Specht modules for column-convex diagrams [J. Algebra 174 (1995), 489-522], the present article was linked with the study of Specht modules associated with column-convex diagrams and some related diagrams. In the paper Balanced labellings and Schubert polynomials [European J. Combin. 18 (1997), 373-389], Fomin-Greene-Reiner-Shimozono constructed an explicit basis of Schubert modules.

Magyar developed a more geometric approach to Theorem 4.1 in the papers: Four new formulas for Schubert polynomials [preprint, 1995], Borel-Weil theorem for configuration varieties and Schur modules [Adv. Math. 134 (1998), 328-366], Schubert polynomials and Bott-Samelson varieties [Comment. Math. Helv. 73 (1998), 603-636]. This last paper contains another, geometric proof of Theorem 4.1 obtained in collaboration with Reiner and Shimozono. In the paper Standard monomial theory for Bott-Samelson varieties of GL( $n$ ) [Publ. Res. Inst. Math. Sci. 34 (1998), 229-248], Lakshmibai-Magyar mention a link of Theorem 4.1 with Standard monomial theory.

In a recent paper by Buch-Kresch-Tamvakis-Yong Schubert polynomials and quiver formulas [to appear in Duke Math. J.], the authors discuss a relationship between a representation-theoretic interpretation of their work and Theorem 4.1.

See also a paper by Lascoux-Schützenberger: Fonctorialité des polynômes de Schubert [Contemp. Math. 88 (1989), 585-598].

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[^1]:    ${ }^{1}$ Following the classical route ([13], [5], and [9]), we work here over commutative $\mathbb{Q}$-algebras $R$. We discuss generalizations of our results to arbitrary commutative rings or, in some cases, infinite fields in Remark 5.3.

[^2]:    ${ }^{2}$ We display a permutation by the sequence of its values; the composition of permutations will be denoted by "o".

