# LECTURES ON SCHUBERT VARIETIES 

SARA BILLEY'S SPRING 2007 COURSE

## 1. Introduction(March 28, 2007)

### 1.1. History of Schubert Calculus.

(1) Enumerative geometry

- Hermann Grassmann (1809-1877)
- Hermann Cäsar Hannibal Schubert (1848-1911)
- Francesco Severi (1879-1961)
- Mario Pieri
- Giovanni Zeno Giambelli (1879-1935)
(2) Topology of homogeneous spaces
- Charles Ehresmann (1905-1979)
- Claude Chevalley (1907-1984)
(3) Representation theory (in 1950's)
- A. Borel
- R. Bott
- B. Kostant
(4) Explicit computation
- I. Berenstein
- I. Gelfand
- S. Gelfand
- M. Demazure
(5) More concrete(combinatorial) theory
- A. Lascoux
- M.P. Schützenberger


### 1.2. Enumerative Geometry.

(1) Some of typical questions in enumerative geometry

- Given two planes $P_{1}, P_{2} \subseteq \mathbb{R}^{3}$ through the origin, what is $\operatorname{dim}\left(P_{1} \cap P_{2}\right)$ ? or equivalently, given two lines in projective space $\mathbb{P}^{2}$, how many points are in the intersection?
- How many lines meet four given lines in $\mathbb{R}^{3}$ ?
(2) An example of Schubert variety
- A family of lines meeting a point and contained in a 2 -dimensional space.
(3) Schubert problem
- Count the number of points in the intersection of Schubert varieties such that the intersection is 0 -dimensional.
(4) Many flavors
- Grassmannian manifold(variety)
- Flag manifold
- Affine Grassmannian manifold
- Partial flag varieties
- GKM ...
(5) Common Theme
- Schubert varieties form a nice basis of $H^{*}(X, \mathbb{Z})$.


## 2. Flag Variety

Let $G$ be the general linear group $G L_{n}(\mathbb{C})$, and $B$ be the set of upper triangular matrices in $G$. Then for $g \in G$, the coset $g B$ is determined by the subspaces $\left\langle\mathbf{c}_{1}, \ldots, \mathbf{c}_{i}\right\rangle, i=1, \ldots, n$, where $\mathbf{c}_{i}$ is the $i$ th column vector of $g$.
Definition 2.1. A flag is a nested sequence of vector subspaces $F_{\bullet}=\left(F_{1} \subset F_{2} \subset \cdots \subset F_{k}\right)$ of $\mathbb{C}^{n}$. A flag is complete if $\operatorname{dim} F_{i}=i$ and $k=n$.

With the above definition we have a bijection between $G / B$ and the set of complete flags in $\mathbb{C}^{n}$. Observe that we inherently chose a basis $\left\{\mathbf{e}_{1}, \ldots, \mathbf{e}_{n}\right\}$ for $\mathbb{C}^{n}$ so that $\mathbf{c}_{i}=\sum_{k=1}^{i} g_{k i} \mathbf{e}_{k}$. Equivalently, we chose a base (complete) flag

$$
\left\langle\mathbf{e}_{1}\right\rangle \subset\left\langle\mathbf{e}_{1}, \mathbf{e}_{2}\right\rangle \subset \cdots \subset\left\langle\mathbf{e}_{1}, \ldots, \mathbf{e}_{n}\right\rangle .
$$

How can we canonically represent $g B$ ? Right multiplication by $B$ adds some multiple of the $i$ th column of $g$ to $j$ th column for $i \leq j$. Therefore we can always find $b \in B$ so that $g b$ is in its column echelon form; the lowest non-zero entry in each column is a 1 and the entries to the right of each leading one are all zeros. The leading 1's in each column form a permutation matrix. For example,

$$
\left[\begin{array}{llll}
3 & 1 & 0 & 0 \\
2 & 0 & 2 & 1 \\
1 & 0 & 0 & 0 \\
0 & 0 & 1 & 0
\end{array}\right] \text { is the canonical form of the matrix } g=\left[\begin{array}{llll}
9 & 5 & 9 & 7 \\
6 & 2 & 4 & 0 \\
3 & 1 & 0 & 0 \\
0 & 0 & 2 & 2
\end{array}\right] .
$$

This permutation determines the position of $g$ with respect to the base flag, denoted $\operatorname{pos}(g)=w$.
We will write $w=\left[w_{1}, w_{2}, \cdots, w_{n}\right]$ to mean $w$ is the bijection mapping $i$ to $w_{i}$. This is the one-line notation for $w \in S_{n}$. In the above example, $w=[2413]$ is the position of $g$ with respect to the base flag. In general, we read the column numbers of leading 1's in the column echelon form of $g$ from top to bottom to get the one-line notation for $w$ when $w$ determines the position of $g$ with respect to the base flag.

WARNING: there are at least 8 different names used for this permutation matrix in the literature on Schubert varieties! We have chosen this way of naming our permutation matrices to agree with the diagram of the permutation which is defined below.
Definition 2.2. For $w \in S_{n}$, the associated Schubert cell is defined as

$$
C_{w}=\{g \in G \mid \operatorname{pos}(g)=w\} .
$$

Example 2.3. The Schubert cell corresponding to $w=\left[\begin{array}{lll}2 & 4 & 1\end{array}\right]$ is $C_{2413}=\left\{\left[\begin{array}{cccc}* & 1 & 0 & 0 \\ * & 0 & * & 1 \\ 1 & 0 & 0 & 0 \\ 0 & 0 & 1 & 0\end{array}\right]\right\}$.
We now can see that invertible matrices in column echelon form determine a set of representatives of $G / B$. Therefore, $G / B$ can be written as a disjoint union $\bigcup_{w \in S_{n}} C_{w}$.
Definition 2.4. For $w \in S_{n}$, diagram of $w$ is defined as follows:

$$
D(w)=\left\{\left(i, w_{j}\right) \in[n]^{2} \mid i<j \text { and } w_{i}>w_{j}\right\} .
$$

Here, $[n]=\{1,2, \cdots, n\}$.

An inversion in $w$ is a pair $1 \leq i<j \leq n$ such that $w_{i}>w_{j}$. The length of $w$, denoted $\ell(w)$, is the number of inversions for $w$. The pairs in $D(w)$ are in bijection with the inversions for $w$. Therefore, we can observe that for $w \in S_{n}$,

$$
\operatorname{dim}_{\mathbb{C}} C_{w}=\ell(w)=\text { number of inversions of } w=|D(w)| .
$$

Let $T$ be the set of diagonal matrices in $G$. The invertible diagonal matrices form an abelian group that is isomorphic to $\left(\mathbb{C}^{*}\right)^{n}$.

How does $T$ act on $C_{w}$ under the left multiplication? A diagonal matrix acts on a matrix $g$ by scaling the rows of $g$. Therefore, the column echelon form of $g$ is unchanged, hence $T \cdot C_{w}=C_{w}$. Moreover, the fixed points of the $T$-action on $G / B$ are exactly the permutation matrices; one $T$-fixed point in each $C_{w}$.

How does the left $B$-action work on $C_{w}$ ? It is easy to see that $B C_{w} \subset C_{w}$, and since the identity matrix is in $B$, we have $B C_{w}=C_{w}$, which is same as $B w$.

Exercise 2.5. Prove that there are enough free variables to get any matrix in $C_{w}$ by $B w$.

## Theorem 2.6. (Bruhat decomposition)

$$
G=\bigcup_{w \in S_{n}} B w B
$$

and hence we have

$$
G / B=\bigcup_{w \in S_{n}} B w=\bigcup_{w \in S_{n}} C_{w}
$$

Remark 2.7. The Schubert cell $C_{w}$ can be viewed in three ways:
(1) The set of cosets $B \cdot w B / B$.
(2) The orbit of the permutation matrix $B \cdot w$.
(3) The set of complete flags in position $w$ with respect to the base flag.

All three points of view are useful.
If we understand elements in $C_{w}$ as complete flags, what is the condition for a flag to be an element of $C_{w}$ in terms of equations?

Definition 2.8. For $w \in S_{n}$, let $r k_{\nwarrow} w_{[i, j]}$ be the rank of the submatrix of $w$ with top left corner at $(1,1)$ and lower right corner at $(i, j)$.

Note that $w$ can be completely recovered from the matrix $r k \nwarrow w$.
Then we have another description for Schubert cell:

$$
C_{w}=\left\{\text { complete flag } F_{\bullet} \mid\left\langle\mathbf{e}_{1}, \ldots, \mathbf{e}_{i}\right\rangle \cap F_{j}=r k_{\nwarrow} w_{[i, j]}\right\} .
$$

Example 2.9. For $w=[2413]$, the condition for a complete flag $F_{\bullet}$ to be in $C_{w}$ is

$$
\left[\operatorname{dim}\left(\left\langle\mathbf{e}_{1}, \ldots, \mathbf{e}_{i}\right\rangle \cap F_{j}\right)\right]_{i, j}=\left[\begin{array}{llll}
0 & 1 & 1 & 1 \\
0 & 1 & 1 & 2 \\
1 & 2 & 2 & 3 \\
1 & 2 & 3 & 4
\end{array}\right]
$$

Definition 2.10. (coordinate free description of Schubert cells) For a fixed base flag $B_{\boldsymbol{\bullet}}=\left(B_{1} \subset\right.$ $B_{2} \subset \cdots \subset B_{n}$ ), and $w \in S_{n}$, the corresponding Schubert cell is defined by

$$
C_{w}\left(B_{\bullet}\right)=\left\{\text { complete flags } F_{\bullet} \mid \operatorname{dim}\left(B_{i} \cap F_{j}\right)=r k<w_{[i, j]}\right\} .
$$

## Summary

$G=G L_{n}(\mathbb{C})$, invertible matrices.
$B=$ upper triangular matrices.
$T=$ diagonal matrices, maximal torus.
$W=S_{n}$.

## 3. Schubert varieties and Bruhat order (March 30, 2007)

Definition 3.1. The Zariski topology on $\mathbb{A}^{n}=\mathbb{C}^{n}$ (or $\mathbb{P}^{n}$ ) has closed sets given by the vanishing sets of some collection of polynomials (homogeneous polynomials, respectively).

Note that $n \times n$ complex matrices form an affine space of dimension $n^{2}$. Our goal is to describe the closure of a Schubert cell in $G / B$ by equations.
Example 3.2. For $w=\left[\begin{array}{lll}2 & 4 & 13\end{array}\right]$, an element in $C_{w}=B w / B$ has the following form: $\left[\begin{array}{llll}* & 1 & 0 & 0 \\ * & 0 & * & 1 \\ 1 & 0 & 0 & 0 \\ 0 & 0 & 1 & 0\end{array}\right]$. Therefore, if $M=\left(m_{i j}\right) \in B w B$, then the following equations on $m_{i j}$ 's must be satisfied:

$$
m_{41}=m_{42}=0 \text {, } \operatorname{det}\left[\begin{array}{ll}
m_{21} & m_{22} \\
m_{31} & m_{32}
\end{array}\right]=0 \text {, det }\left[\begin{array}{lll}
m_{21} & m_{22} & m_{23} \\
m_{31} & m_{32} & m_{33} \\
m_{41} & m_{42} & m_{43}
\end{array}\right]=0 \text {. }
$$

Definition 3.3. For any $n \times n$ matrix $M$, let $r k_{\swarrow} M_{[i, j]}$ be the rank of the smallest lower left submatrix of $M$ containing $M_{i j}$. Then the rank of $M$ is the size of largest non-vanishing minors of $M$.

Exercise 3.4. Let $M$ be any $k \times l$ matrix and let $r(p, q)=r k_{\swarrow} M_{[p, q]}$. Show that there exists a $k \times l$ matrix $U$ with all 0,1 -entries and at most one 1 n each row and column such that $r(p, q)=r k_{\swarrow} U_{[p, q]}$ for all $p, q$.
Lemma 3.5. For $w \in S_{n}, M \in B w B$ if and only if $r k_{\swarrow} M_{[p, q]}=r k_{\swarrow} w_{[p, q]}$ for all $p, q$. Equivalently, $N=b_{1} M b_{2}$ for some $b_{1}, b_{2} \in B$ if and only if $r k_{\swarrow} M=r k_{\swarrow} N$.
Proof. Left (or right) multiplication by $B$ adds some multiple of row (column) $j$ to row (column) $i$ for $i<j$.

Notation Let $r_{w}(p, q)=r k_{\swarrow} w_{[p, q]}$.
Definition 3.6. Matrix Schubert variety is defined as follows:

$$
\begin{aligned}
\tilde{X}(w) & =\left\{M \in M a t_{n \times n}(\mathbb{C}) \mid r k_{\swarrow} M_{[p, q]} \leq r k_{\swarrow} w_{[p, q]} \text { for all } p, q\right\} \\
& =\left\{M \in \operatorname{Mat}_{n \times n}(\mathbb{C}) \mid \text { all } r_{w}(p, q)+1 \text { minors vanish on }\left[\begin{array}{ccc}
m_{p 1} & \cdots & m_{p q} \\
\vdots & \cdots & \vdots \\
m_{n 1} & \cdots & m_{n q}
\end{array}\right] \text { for all } p, q\right\} .
\end{aligned}
$$

We let $I_{w}$ be the ideal generated by all $r_{w}(p, q)+1$ minors of $\left[\begin{array}{ccc}x_{p 1} & \cdots & x_{p q} \\ \vdots & \cdots & \vdots \\ x_{n 1} & \cdots & x_{n q}\end{array}\right]$ for all $p, q$.
Remark 3.7. $G \subseteq M_{n \times n}$ is open in Zariski topology. Hence $\widetilde{X}(w) \cap G$ is defined by the same equations.
Since vanishing/non-vanishing minors are unchanged by right multiplication by $B, I_{w}$ also defines $\overline{C_{w}}$ in $G / B$, and we have the following definition:

Definition 3.8. For $w \in S_{n}$, the associated Schubert variety $X_{w}$ is defined as the closure of $C_{w}$ in Zariski topology:

$$
X_{w}=\overline{C_{w}} .
$$

Note that $X_{w}$ is a projective variety.
Now, our question is the following: Which Schubert cells are in $X_{w}$ ?
Remark 3.9. For $w \in S_{n}$ and $i<j$, if $w(i)>w(j)$ then $r_{w^{\prime}}(p, q)=r_{w}(p, q)-1$ for $i<p \leq j$ and $w(j) \leq q<w(i)$, where $w^{\prime}=w t_{i j}$.

Proposition 3.10. (Ehressmann Tableaux Criterion 1934) Let $v=\left[v_{1}, v_{2}, \ldots, v_{n}\right]$ and $w=\left[w_{1}, w_{2}, \ldots, w_{n}\right]$ be two permutations in $S_{n}$.
(1) $C_{v} \subseteq X_{w} \Leftrightarrow X_{v} \subseteq X_{w} \Leftrightarrow\left\{v_{i}, \ldots, v_{n}\right\}_{<} \geq\left\{w_{i}, \ldots, w_{n}\right\}_{<}$for all i.
(2) $w$ cover $v$ in containing order of Schubert variety if and only if $w t_{i j}=v$ and $l(w)=l(v)+1$.

Example 3.11. $X_{531624} \subseteq X_{651342}$, but $X_{3214} \nsubseteq X_{4132}$ since $\{4,1\}_{<} \not \leq\{3,2\}_{<}$.
Proof. (of Proposition 3.10)
(1)

$$
\begin{aligned}
C_{v} \subseteq X_{w} & \Leftrightarrow r_{v}(p, q) \leq r_{w}(p, q) \text { for all } p, q \\
& \Leftrightarrow\left|\left\{v_{p}, \ldots, v_{n}\right\} \cap\{1, \ldots, q\}\right| \leq\left|\left\{w_{p}, \ldots, w_{n}\right\} \cap\{1, \ldots, q\}\right| \text { for all } p, q \\
& \Leftrightarrow\left\{v_{p}, \ldots, v_{n}\right\}_{<} \geq\left\{w_{p}, \ldots, w_{n}\right\}_{<} \text {for all } p .
\end{aligned}
$$

(2) $(\Leftarrow)$ Compare $r_{v}$ and $r_{w}$. Since $l(w)=l(v)+1$, no other 1 is in boundary box for $t_{i j}$. Therefore, there is no $z \in S_{n}$ such that $r_{v}<r_{z}<r_{w}$.
$(\Rightarrow)$ Suppose that $X_{v} \subset X_{w}$ and $w$ covers $v$ in containing order, so no $z$ exists such that $X_{v} \subset X_{z} \subset$ $X_{w}$. Let $j$ be the largest such that $v_{j} \neq w_{j}$, then $w_{j}<v_{j}$. Let $i$ be the largest index such that $i<j$ and $v_{j}>v_{i} \geq w_{j}$.
claim: $X_{v} \subset X_{v t_{i j}} \subseteq X_{w}$, and hence $v t_{i j}=w$ and $l(w)=l(v)+1$.
proof of claim: We need to show that $X_{v t_{i j}} \subseteq X_{w}$ that is $r_{v t_{i j}} \leq r_{w}$ everywhere. Let $\mathfrak{A}$ be the boxes consisting of rows from $i$ th to $j$ th and columns from $w_{j}$ th to $v_{j}$ th. Let $\mathfrak{B}$ be the boxes consisting of rows from $(i+1)$ st to the $j$ th and columns from $v_{i}$ th to $\left.9 v_{j}-1\right)$ st. Since $r_{v}(p, q) \leq r_{w}(p, q)$ for all $p, q$ and $r_{v}(p, q)=r_{v} t_{i j}(p, q)$ outside of $\mathfrak{B}, r_{v t_{i j}}(p, q) \leq r_{w}(p, q)$ outside of $\mathfrak{B}$. Since $v_{n}=w_{n}, \ldots v_{j+1}=$ $w_{j+1}$ and $v_{j}>w_{j}$, we have $r_{v} t_{i j}(j, q)=r_{v}(j, q)+1=r_{w}(j, q)$ for all $w_{j} \leq q<v_{j}$. Furthermore, $r_{v t_{i j}}\left(p, w_{j}-1\right)=r_{v}\left(p, w_{j}-1\right) \leq r_{w}\left(p, w_{j}-1\right)$ for $i \leq p \leq j$.

By the choice of $i$, there are no other 1's in matrix for $v$ in $\mathfrak{A}$, so all jumps in rank in $\mathfrak{B}$ happen because of 1 's to the $S W$ of $\mathfrak{B}$, so accounted for along left edge of $\mathfrak{B}$ or lower edge of $\mathfrak{B}$. Therefore, $r_{v t_{i j}}(p, q) \leq r_{w}(p, q)$ for all $p, q$.

Corollary 3.12. $X_{v} \subseteq X_{w}$ if and only if there exists a sequence of transpositions $t_{a_{1} b_{1}}, \ldots, t_{a_{k} b_{k}}$ such that $w=v t_{a_{1} b_{1}} \cdots t_{a_{k} b_{k}}$ and $l\left(v t_{a_{1} b_{1}} \cdots t_{a_{j} b_{j}}\right)=l(v)+j$ for all $1 \leq j \leq k$.

Exercise 3.13. (Chevalley's criterion) Let $s_{i}=t_{i i+1}$ and write $w=s_{a_{1}} \cdots s_{a_{p}}$ with $p=l(w)$. Then $X_{v} \subseteq X_{w}$ if and only if there exists $1 \leq i_{1}<i_{2}<\cdots<i_{k} \leq p$ such that $v=s_{a_{i_{1}}} \cdots s_{a_{i_{k}}}$.
Definition 3.14. $v \leq w$ in Bruhat-Chevalley order if $X_{v} \leq X_{w}$.
Remark 3.15. (1) $v \leq w \Leftrightarrow\left\{v_{1}, \ldots, v_{i}\right\}_{<} \leq\left\{w_{1}, \ldots, w_{i}\right\}_{<}$for all $i$.
(2) (Björner-Brenti) We only need to check the condition in (1) for $i$ 's such that $v s_{i}<v$ or only need to check $i$ 's such that $w s_{i}>w$.
(3)

$$
\begin{aligned}
X(w) & =\left\{g B \in G / B \mid r_{g} \leq r_{w}\right\} \\
& =\left\{F_{\bullet}: \text { complete flag } \mid \operatorname{dim}\left(\left\langle\mathbf{e}_{\mathbf{1}}, \ldots, \mathbf{e}_{\mathbf{i}}\right\rangle \cap F_{j}\right) \geq r k<w_{i j}\right\}
\end{aligned}
$$

## 4. Fulton's Essential Set (April 4, 2007)

Let $C_{w}$ and $X_{w}$ be the Schubert cell and Schubert variety, respectively, that are associated to the permutation $w \in S_{n}$.
Definition 4.1. Let $D^{\prime}(w)=\left\{(i, j) \in[n]^{2}\right.$ such that $w(i)>j$ and $\left.w^{-1}(j)<i\right\}$. This (alternate) diagram is obtained from the entries of the permutation matrix associated to $w$ by canceling entries that lie to the North and East of a 1 entry.

The essential set of $w$, denoted $E s s^{\prime}(w)$ consists of the Northeast corners of connected components in $D^{\prime}(w)$. More precisely, $E s s^{\prime}(w)=\left\{(i, j) \in D^{\prime}(w)\right.$ such that $(i-1, j),(i, j+1)$ and $(i-1, j+1) \notin$ $D^{\prime}(w)$.
Example 4.2. We have that $D^{\prime}(2413)$ consists of the entries which are circles in the matrix:

$$
\left[\begin{array}{cccc}
\cdot & 1 & \cdot & . \\
\cdot & \bullet & \cdot & 1 \\
1 & \cdot & \cdot & \cdot \\
0 & \bullet & 1 & \cdot
\end{array}\right]
$$

and $E s s^{\prime}(2413)$ consists of the $\bullet$ entries.
Observe that the cardinality of $D^{\prime}(w)$ is the codimension of $X_{w}$, which is $\binom{n}{2}-l(w)$.
We previously found that

$$
M \in X_{w} \Longleftrightarrow r k_{\swarrow}(M) \leq r k_{\swarrow}(w) \Longleftrightarrow r k_{\nwarrow}(M) \geq r k_{\nwarrow}(w)
$$

$\Longleftrightarrow$ all $r k_{\swarrow} w_{[i, j]}+1$ minors in $X_{[i, j]}$ vanish on $M$
where $X_{[i, j]}$ is $\left[\begin{array}{ccc}x_{i 1} & \ldots & x_{i j} \\ \ldots & & \ldots \\ x_{n 1} & \ldots & x_{n j}\end{array}\right]$.
Let $I_{w}$ be the ideal generated by all $r k_{\swarrow} w_{[i, j]}+1$ minors in $X_{[i, j]}$. Here are some observations.
(1) $E s s^{\prime}(w)$ is all on one row if and only if $w$ has at most one ascent. Here, we mean that the entries in the 1 -line notation for $w=\left(w_{1}, w_{2}, \ldots, w_{n}\right)$ satisfy $w_{1}>w_{2}>\ldots w_{k}, w_{k+1}>w_{k+2}>$ $\cdots>w_{n}$.
(2) All of the entries in $E s s^{\prime}(w)$ correspond to 0 entries in the canonical matrix form for $C_{w}$.
(3) For $(i, j) \in D^{\prime}(w)$, we have $r k_{\swarrow} w_{[i, j]}=\# 1$ entries Southwest of $[i, j]$. This quantity is also the number of lines crossed as we look South from $[i, j]$. Equivalently, it is the number of lines crossed as we look West from $[i, j]$.
(4) $r k_{\swarrow} w$ is constant on connected components of $D^{\prime}(w)$.

Proposition 4.3. (Fulton) Let I be the ideal generated by the $r k_{\swarrow} w_{[i, j]}+1$ minors of $X_{[i, j]}$ for all $(i, j) \in E s s^{\prime}(w)$. Then, for any $w \in S_{n}$, the ideal $I_{w}=I$.
Proof. It suffices to show that all of the defining minors of $I_{w}$ are in $I$.
Case 1. Suppose $(p, q) \in D^{\prime}(w)$. Then, $(p, q)$ is in the same connected component as some $(i, j) \in$ $E s s^{\prime}(w)$ with $i \leq p$ and $j \geq q$. We have $r k_{\swarrow} w_{[p, q]}=r k_{\swarrow} w_{[i, j]}$ by (4) and all $r k_{\swarrow} w_{[p, q]}+1$ minors of $X_{[p, q]}$ are also minors in $X_{[i, j]}$.

Case 2. Suppose $(p, q)$ lies on a vertical edge of the matrix describing the diagram of $w$. More precisely, suppose $(p, q) \notin D^{\prime}(w)$ with $w^{-1}(q) \geq p$ and $w(p) \geq q$. Then, we find the largest $k$ such that $w^{-1}(q-i)>p$ for all $0 \leq i<k$. If $k=q$ then $r k_{\swarrow} w_{[p, q]}=q$ so there are no $r k_{\swarrow} w_{[p, q]}+1$ minors in $X_{[p, q]}$. Otherwise, $k<q$ and $(p, q-k) \in D^{\prime}(w)$. We know the defining minors for $(p, q-k)$ are in $I$ by Case 1, and we will show the defining minors for $(p, q)$ are in the ideal generated by these.

Note that the $r$-minors of a matrix $M$ generate an ideal that contains all $(r+1)$-minors of $M^{\prime}$, where $M^{\prime}$ is obtained from $M$ by adding a column. If $r=r k_{\swarrow} w_{[p, q-k]}$ then $r k_{\swarrow} w_{[p, q-i]}=r+(k-i)$ for $0 \leq i \leq k$. Hence, all $r k_{\swarrow} w_{[p, q]}+1$ minors of $X_{[p, q]}$ are in $I$ by Case 1 .

Case 3. Suppose $(p, q)$ lies on a horizontal edge of the matrix describing the diagram of $w$. More precisely, suppose $(p, q) \notin D^{\prime}(w)$ with $w^{-1}(q)<p$ and $w(p)<q$. This case is very similar to Case 2 , and can be completed by looking down instead of left.

Case 4. Suppose $(p, q)$ lies on an intersection of edges in the matrix describing the diagram of $w$. More precisely, suppose $(p, q) \notin D^{\prime}(w)$ with $w^{-1}(q)>p$ and $w(p)<q$. If $w^{-1}(q-1) \leq p$ then by Case 2, we have that the $r k_{\swarrow} w_{[p, q-1]}+1$ minors are in $I$. By the note in Case 2, this implies that all $r k_{\swarrow} w_{[p, q]}+1$ minors of $X_{[p, q]}$ are in $I$.

If $w^{-1}(q-1)>p$ then by induction we may assume these minors are in $I$. We may use the note in Case 2 to see that all $(p, q)$ minors are in $I$.

These cases exhaust all of the possibilities, so we have completed the proof.
The following proposition shows that no proper subset of the essential set will work to define $I_{w}$.
Proposition 4.4. (Fulton) Let $w \in S_{n}$ and $\left(p_{0}, q_{0}\right) \in E s s^{\prime}(w)$. Then there exists an $n \times n$ matrix $M$ such that $r k_{\swarrow} M_{[p, q]} \leq r k_{\swarrow} w_{[p, q]}$ for all $(p, q) \in E s s^{\prime}(w) \backslash\left\{\left(p_{0}, q_{0}\right)\right\}$, and $r k_{\swarrow} M_{\left[p_{0}, q_{0}\right]}=r k_{\swarrow} w_{\left[p_{0}, q_{0}\right]}+1$.
Proof. Define $M=\left(m_{i j}\right)$ by

$$
m_{i j}= \begin{cases}w_{i j} & \text { if } i \geq p_{0}, j \leq q_{0} \text { but }(i, j) \neq\left(p_{0}, q_{0}\right) \\ 1 & \text { if }(i, j)=\left(p_{0}, q_{0}\right) \\ 0 & \text { otherwise }\end{cases}
$$

Then, $r k_{\swarrow} M_{[p, q]}$ is the number of 1 entries lying Southwest of $(p, q)$, so $M$ has the desired properties.

Warning 4.5. The set of determinantal equations corresponding to $r k_{\swarrow} w_{[p, q]}+1$ minors for $(p, q) \in$ $E s s^{\prime}(w)$ is not necessarily minimal. A Gröbner basis is given by Knutson-Miller.

Open Question 4.6. What is the analogue of the essential set for $G / B$ in other types?
Example 4.7. For $w=2413$ we have $D^{\prime}(w)$ given by

$$
\left[\begin{array}{cccc}
. & 1 & \cdot & . \\
. & \bullet & \cdot & 1 \\
1 & \cdot & . & \cdot \\
0 & \bullet & 1 & \cdot
\end{array}\right]
$$

so the essential set tells us us that $I_{w}$ is generated by the $1 \times 1$ and $2 \times 2$ minors of

$$
X_{[4,2]}=x_{4,1} \quad x_{4,2} \quad X_{[2,2]}=\begin{array}{ll}
x_{2,1} & x_{2,2} \\
x_{3,1} & x_{3,2} . \\
x_{4,1} & x_{4,2}
\end{array}
$$

So

$$
I_{w}=\left\langle x_{41}, x_{42}, \operatorname{det}\left[\begin{array}{ll}
x_{21} & x_{22} \\
x_{31} & x_{32}
\end{array}\right], \operatorname{det}\left[\begin{array}{ll}
x_{21} & x_{22} \\
x_{41} & x_{42}
\end{array}\right], \operatorname{det}\left[\begin{array}{ll}
x_{31} & x_{32} \\
x_{41} & x_{42}
\end{array}\right]\right\rangle
$$

which can be reduced to

$$
I_{w}=\left\langle x_{41}, x_{42}, \operatorname{det}\left[\begin{array}{ll}
x_{21} & x_{22} \\
x_{31} & x_{32}
\end{array}\right]\right\rangle
$$

Example 4.8. For $w=3412$ we have $D^{\prime}(w)=E s s^{\prime}(w)$ given by

$$
\left[\begin{array}{cccc}
. & . & 1 & . \\
. & . & \bullet & 1 \\
1 & . & \cdot & \cdot \\
\cdot & 1 & \cdot & \cdot
\end{array}\right]
$$

so

$$
I_{w}=\left\langle x_{41}, \operatorname{det}\left[\begin{array}{lll}
x_{21} & x_{22} & x_{23} \\
x_{31} & x_{32} & x_{33} \\
x_{41} & x_{42} & x_{43}
\end{array}\right]\right\rangle
$$

Example 4.9. For $w=4321$ we have $D^{\prime}(w)$ given by

$$
\left[\begin{array}{cccc}
. & . & . & 1 \\
. & . & 1 & . \\
. & 1 & . & \cdot \\
1 & . & . & \cdot
\end{array}\right]
$$

so $E s s^{\prime}(w)$ is the empty set. In this case, $X_{w}=G / B$.
Definition 4.10. Let $w_{0}$ be the permutation $(n, n-1, n-2, \ldots, 2,1)$. Then $\tilde{X}\left(w_{0}\right)=M_{n \times n}$ and $X_{w_{0}}=G / B=\cup_{v \leq w_{0}} C_{v}$.
Definition 4.11. The coordinate ring of $X_{w}$ is $\mathbb{C}\left[x_{11}, \ldots, x_{n n}\right] / I_{w}$, consisting of polynomial functions on $X_{w}$.

The standard monomials form a basis for the coordinate ring. This set was originally described by Lakshmibai-Musili-Seshadri. See also the paper by Reiner-Shimozono in our class library.

We claim that $G / B$ is a smooth manifold. This means that every point has an affine neighborhood of dimension $\binom{n}{2}$ and at every point, the dimension of the tangent space is $\binom{n}{2}$. Note that $C_{w_{0}}$ consists of matrices of the form $\left[\begin{array}{llll}* & * & * & 1 \\ * & * & 1 & 0 \\ * & 1 & 0 & 0 \\ 1 & 0 & 0 & 0\end{array}\right]$. In particular, the permutation matrix $w_{0}=\left[\begin{array}{cccc}0 & 0 & 0 & 1 \\ 0 & 0 & 1 & 0 \\ 0 & 1 & 0 & 0 \\ 1 & 0 & 0 & 0\end{array}\right] \in C_{w_{0}}$, and $C_{w_{0}}$ is its affine neighborhood of dimension $\binom{n}{2}$.

Since $G / B$ has a transitive $G L_{n}$ action by left multiplication, every point looks locally like any other point. Therefore, every point in $G / B$ has an affine neighborhood of dimension $\binom{n}{2}$, and the tangent space to every point has the same dimension, hence $G / B$ is smooth. This proves the claim.

Are all Schubert varieties smooth? No, and as $n$ tends to $\infty$, we find that almost all Schubert varieties are singular. We outline a way to test for smoothness based on the defining equations for $X_{w}$.

The local properties of points in a Schubert variety $X_{w}$ are determined by the permutation matrices for $v \leq w$ since $X_{w}$ is the union of the $B$-orbits of these permutation matrices: $X_{w}=\cup_{v \leq w} C_{v}=\cup_{v \leq w} B v$. Smoothness is a local property so $X_{w}$ is smooth at every point in $C_{v}$ for $v \leq w$ if and only if $X_{w}$ is smooth at $v$.

Theorem 4.12. (Jacobian criterion) Let $Y$ be an affine variety in $\mathbb{A}^{n}$ defined by $I(Y)=\left\{f_{1}, \ldots, f_{r}\right\}$ where the $f_{i}$ are polynomials in variables $x_{1}, \ldots, x_{n}$. Then, we define the Jacobian matrix $J\left(x_{1}, \ldots, x_{n}\right)=$ $\left(\frac{\partial f_{i}}{\partial x_{j}}\right)$. If $p=\left(p_{1}, \ldots, p_{n}\right) \in \mathbb{A}^{n}$ then the following hold:
(1) We always have $\operatorname{rk} J\left(p_{1}, \ldots, p_{n}\right) \leq \operatorname{codim}_{\mathbb{A}^{n}} Y=n-\operatorname{dim} Y$.
(2) The point $p$ is a smooth point of $Y$ if and only if rk $J\left(p_{1}, \ldots, p_{n}\right)=\operatorname{codim}_{\mathbb{A}^{n}} Y=n-\operatorname{dim} Y$.

Note that $\{p \in Y: r k J<k\}$ is a closed set in the Zariski topology, being defined by the vanishing of minors. The intersection of two closed sets is again a closed set. We said above that if $X_{w}$ is singular at $v \leq w$, then $X_{w}$ is singular at all points in $C_{v}$. Therefore, $X_{w}$ is singular on $\overline{C_{v}}=X_{v}$. Equivalently, $X_{w}$ is singular at $u$ for all $u \leq v$. Hence, we have the following corollary.

Corollary 4.13. The Schubert variety $X_{w}$ is smooth everywhere if and only if $X_{w}$ is smooth at $v=i d$.

## 5. Smooth Schubert Varieties(April 6,2007)

Recall that $X_{w}$ is smooth everywhere if and only if $X_{w}$ is smooth at the identity matrix by Corollary 4.13. Our goal is to test smoothness with the Jacobian criterion, which requires us to identify an affine neighborhood of the identity.

Consider $G / B=X_{w_{0}}=C_{w_{0}} \cup \cup_{v<w_{0}} C_{v}$. The Schubert cell $C_{w_{0}}$ is an affine neighborhood of $w_{0}$. For example when $n=4$

$$
\left(\begin{array}{llll}
0 & 0 & 0 & 1  \tag{5.1}\\
0 & 0 & 1 & 0 \\
0 & 1 & 0 & 0 \\
1 & 0 & 0 & 0
\end{array}\right) \in\left(\begin{array}{cccc}
* & * & * & 1 \\
* & * & 1 & 0 \\
* & 1 & 0 & 0 \\
1 & 0 & 0 & 0
\end{array}\right)=C_{w_{0}}
$$

We can move this neighborhood around to contain the identity by left multiplication by the matrix $w_{0}$

$$
w_{0} C_{w_{0}}=w_{0} B w_{0}=\left(\begin{array}{cccc}
0 & 0 & 0 & 1  \tag{5.2}\\
0 & 0 & 1 & 0 \\
0 & 1 & 0 & 0 \\
1 & 0 & 0 & 0
\end{array}\right)\left(\begin{array}{cccc}
* & * & * & 1 \\
* & * & 1 & 0 \\
* & 1 & 0 & 0 \\
1 & 0 & 0 & 0
\end{array}\right)=\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
x_{21} & 1 & 0 & 0 \\
x_{31} & x_{32} & 1 & 0 \\
x_{41} & x_{42} & x_{43} & 1
\end{array}\right) .
$$

In the matrix on the right, we have filled in the stars with affine coordinates.
Definition 5.1. Set $Y(w, i d) \equiv X_{w} \cap w_{0} C_{w_{0}}$ and $Y(w, v) \equiv X_{w} \cap v w_{0} C_{w_{0}}$.
Note, $Y(w, v)$ is an affine neighborhood containing $v$ in $X_{w}$ if $v \leq w$.
Example 5.2. Is $X(2413)$ smooth?
In Example 4.7, we determined that $X_{w}$ is defined by $I_{2413}=<x_{41}, x_{42}, x_{21} x_{32}-x_{31} x_{22}>$. Therefore, $Y(2413, i d)$ is determined by the same equations evaluated on the last matrix in (4.7): < $x_{41}, x_{42}, x_{21} x_{32}-x_{31}>$ where $x_{22}$ was set equal to 1 .

To use the Jacobian criterion, calculate

$$
J=\left(\begin{array}{ccccc}
0 & 0 & 0 & 1 & 0  \tag{5.3}\\
0 & 0 & 0 & 0 & 1 \\
x_{32} & -1 & x_{21} & 0 & 0
\end{array}\right)
$$

where the columns are indexed by the ordered list of variables $x_{21}, x_{31}, x_{32}, x_{41}, x_{42}$. The matrix $J(I)$ is obtained by setting all the variables $x_{i j}$ equal to 0 . Therefore, the rank of $J(I)$ is 3 , while the codimension of $X(2413)$ is $\binom{4}{2}-3=3$. Hence, $X(2413)$ is smooth.
Example 5.3. Is $X(3412)$ smooth?
The affine variety $Y_{3412, i d}$ is defined by $I_{3412}$ restricted to the matrix in (5.2) again. This ideal is generated by $<f_{1}, f_{2}>$, where

$$
\begin{align*}
f_{1} & =x_{41} \\
f_{2} & =\left(\left.\begin{array}{ccc}
x_{21} & 1 & 0 \\
x_{31} & x_{32} & 1 \\
0 & x_{42} & x_{43}
\end{array} \right\rvert\,=x_{21}\left(x_{32} x_{43}-x_{42}\right)-x_{31} x_{43}\right. \\
J & =\left(\begin{array}{cccccc}
0 & 0 & 0 & 1 & 0 & 0 \\
x_{32} x_{43}-x_{42} & -x_{43} & 0 & 0 & x_{21} & x_{21} x_{32}-x_{31}
\end{array}\right) \tag{5.4}
\end{align*}
$$

where the variables are in order $x_{21}, x_{22}, x_{31}, x_{41}, x_{42}, x_{43}$. Setting the variables equal to 0 we see that the rank of $J(I)$ is 1 , whereas the codimension of $X(3412)$ is $\binom{4}{2}-4=2$, so $X(3412)$ is NOT smooth.

What is $\operatorname{Sing}(X(3412))=\bigcup_{v \text { singular in } X_{w}} X(v)$, the singular locus of $X_{w}$ ?
Definition 5.4. $\operatorname{maxsing}(w)=\left\{v \in S_{n}: v \leq w\right.$ and $v$ is a maximally singular point in $\left.X_{w}\right\}$.
Claim 5.5. $\operatorname{Sing}(X(3412))=X(1324)$ or equivalently $\operatorname{maxsing}(w)=\{1324\}$.
Proof. We need to verify the following statements:
(1) $v=1324$ is a singular point of $X(3412)$.
(2) No $t_{i j}$ exists such that $v<v t_{i j} \leq w$ with $l\left(v t_{i j}\right)=l(v)+1$ and $v t_{i j}$ is singular.
(3) All $t_{i j}$ such that $l\left(v t_{i j}\right)=1$, different from $t_{i j}=1324$, are smooth.

To verify (1), let $v=1324$. We have

$$
Y(3412,1324)=\left(\begin{array}{cccc}
1 & 0 & 0 & 0  \tag{5.5}\\
x_{21} & x_{22} & 1 & 0 \\
x_{31} & 1 & 0 & 0 \\
x_{41} & x_{42} & x_{43} & 1
\end{array}\right) \cap X_{w}
$$

The ideal defining $Y(3412,1324)$ is generated by $f_{1}, f_{2}$ again but now we evaluate these polynomials on the matrix in (5.5). So

$$
\begin{gather*}
f_{1}=x_{41} \\
f_{2}=\left|\begin{array}{ccc}
x_{21} & x_{22} & 1 \\
x_{31} & 1 & 0 \\
x_{41} & x_{42} & x_{43}
\end{array}\right|=-x_{31}\left(x_{22} x_{43}-x_{42}\right)+x_{21} x_{43} \\
J=\left(\begin{array}{cccccc}
0 & 0 & 0 & 1 & 0 & 0 \\
x_{43} & -x_{31} x_{43} & -x_{22} x_{32}+x_{42} & 0 & x_{31} & -x_{31} x_{22}+x_{21}
\end{array}\right) \tag{5.6}
\end{gather*}
$$

where the variables are again in lexicographic order. Setting the variables equal to 0 we see that the rank of $J(1324)$ is $1<2=\operatorname{codim}\left(X_{3412}\right)$, so 1324 is a singular point of $X(3412)$.

Exercise 5.6. Complete the above proof by verifying that all permutations with length 2 above 1324 and below 3412 in Bruhat order are smooth points of $X_{w}$, and that 2134 and 1243 are also smooth.

This sort of calculation "works", but it does not give us a good feeling for how the singularities of one Schubert variety relates to the singularities of another. Therefore, we need to study singular Schubert varieties and their tangent spaces as a family. This is the discrete analog of studying a moduli space of curves in order to see their properties in context.

Definition 5.7. For any subgroup $H$ of $G l_{n}(\mathbb{F})$, the Lie algebra of $H, \operatorname{Lie}(H)$, is the space of tangent vectors to the identity matrix.

In many ways, the Lie algebra of a group is easier to work with because $\operatorname{Lie}(H)$ is a vector space. The representation theory of the Lie algebra and its Lie group are closely related. Furthermore, the connected component of $H$ containing the identity matrix can be recovered from $\operatorname{Lie}(H)$.

Observe that $G l_{n} / B \approx S l_{n} /\left(\overline{B \cap S l_{n}}\right)$. Therefore, to understand the tangent spaces to $G / B$ we want to understand $\operatorname{Lie}\left(S l_{n}(\mathbb{C})\right)$. We will discuss Lie algebras in general first and then restrict our attention to this special case by the end of the lecture.

Consider a differentiable path $\phi(t)=\left(\phi_{1}(t), \phi_{2}(t), \ldots, \phi_{k}(t)\right) \in \mathbb{R}^{k}$ with velocity vector

$$
v=\left(\frac{\partial \phi_{1}(x)}{\partial t}, \frac{\partial \phi_{2}(x)}{\partial t}, \ldots, \frac{\partial \phi_{k}(x)}{\partial t}\right)
$$

then $v$ is tangent to $\phi$ at $x \in \mathbb{R}^{k}$.
Definition 5.8. Given any subset $S \subseteq \mathbb{R}^{k}, v$ is tangent to $S$ at $x \in S$ if there exists a differentiable path $\phi(t)$ lying in $S$ such that $\phi(0)=x$ and $\phi^{\prime}(0)=v$.

Recall the gradient of a function $f(x)=f\left(x_{1}, \cdots, x_{k}\right)$ is

$$
\nabla f(x)=\left(\begin{array}{c}
\frac{\partial f}{\partial x_{1}}  \tag{5.7}\\
\dddot{\partial} \\
\frac{\partial f}{\partial x_{k}}
\end{array}\right)
$$

Lemma 5.9. If $S=V\left(f_{1}, \ldots, f_{r}\right) \subseteq \mathbb{R}^{k}$ is a closed set in the Zariski topology, and $f$ vanishes on $S$, then the tangent vectors to $S$ at $x \in S$ are orthogonal to the gradient of $f$ evaluated at $x$.
Proof. Let $\phi(t)$ be a differentiable path in $S$ such that $\phi(0)=x$ and $\phi^{\prime}(0)=v$. Then $f(y)=0$ for all $y \in S \Rightarrow \frac{d}{d t} f(\phi(t))=0$ for all $t \in \mathbb{R}$. This implies $\frac{d}{d t} f(\phi(t))=\frac{\partial f}{\partial x_{1}}(x) \frac{\partial \phi_{1}}{\partial t}(0)+\ldots+\frac{\partial f}{\partial x_{k}}(x) \frac{\partial \phi_{k}}{\partial t}(0)=$ $(\nabla f(x) \cdot v)=0$.
Remark 5.10. Two paths $\phi(t), \xi(t)$ have the same tangent at $x=0$ if and only if the first two terms of their Taylor expansions agree. That is, $\phi(t)=a+b t+\ldots$ and $\xi(t)=a+b t+\ldots$. Therefore, we only need to consider paths with coefficients in $E=\{a+b \epsilon: a, b \in \mathbb{R}\}$, with $\epsilon^{2}=0$. Note that $E$ is a ring with multiplication and addition defined by $(a+b \epsilon)(c+d \epsilon)=a c+(a d+b c) \epsilon$, and $(a+b \epsilon)+(c+d \epsilon)=(a+b)+(c+d) \epsilon$.

Using this map on Taylor series, the set of all paths $\phi(t)$ with $\phi(0)=x$ and $\phi^{\prime}(0)=v$ map to vector $x+v \epsilon=\left(x_{1}+v_{1} \epsilon, \ldots, x_{k}+v_{k} \epsilon\right)$, with $x_{i}, v_{i} \in \mathbb{R}$.
Definition 5.11. Define $v$ to be an infinitesimal tangent of $S=V\left(f_{1}, \ldots, f_{r}\right)$ at $x$ if $f(x+v \epsilon)=0$ as a polynomial over $E$ for all $f$ vanishing on $S$.
Lemma 5.12. Every tangent to $S=V\left(f_{1}, \ldots, f_{r}\right)$ at $x \in S$ is an infinitesimal tangent.
Proof. Given a polynomial $f(x)=f\left(x_{1}, \ldots, x_{k}\right)$ over $E$, then by multivariate Taylor expansion

$$
\begin{equation*}
f(x+v \epsilon)=f(x)+\left(v_{1} \epsilon \frac{\partial f}{\partial x_{1}}(x)+\ldots+v_{k} \epsilon \frac{\partial f}{\partial x_{k}}(x)\right)+\left(v_{1} \epsilon \frac{\partial}{\partial x_{1}}+\ldots+v_{k} \epsilon \frac{\partial}{\partial x_{k}}\right)^{2} f(x)+\ldots \tag{5.8}
\end{equation*}
$$

All the degree 2 and higher terms in this expansion vanish over $E$, hence

$$
f(x+v \epsilon)=f(x)+(\nabla f(x) \cdot v) \epsilon .
$$

If $f$ vanishes on $S$, then $f(x+v \epsilon)=0$ if and only if $\nabla f(x) . v=0$. Apply Lemma 5.9 to conclude $v$ is an infinitesimal tangent.

When $S$ is sufficiently smooth then the converse to Lemma 5.12 holds, namely every infinitesimal tangent is a tangent vector. We will state this precisely for linear algebraic groups and prove it for $S L_{n}$.

Definition 5.13. A group $H$ is a linear algebraic group if it is both a subgroup of $G l_{n}$ and it is defined by the vanishing of some set of polynomials. For example, $S l_{n}=V(\operatorname{det}-1)$ is a linear algebraic group.

Recall that for any matrix $A$,

$$
e^{t A}=I+t A+t^{2} \frac{A^{2}}{2!}+\ldots
$$

If $\left\{e^{t A}: t \in \mathbb{R}\right\}$ is a one-parameter subgroup of $H$, we have a parameterized path $\phi(t)=e^{t A}$ through the identity $\phi(0)$ with velocity vector $\frac{\partial \phi}{\partial t}(0)=A$. Therefore, $A$ is a tangent vector of $H$ at the identity so $A \in \operatorname{Lie}(H)$.

Fact 5.14. Let $H$ be a linear algebraic group. For every infinitesimal tangent $A$ of $H$ at the identity, (i.e. $f(I+A t)=0$ for all $f$ vanishing on $H)$ the group $\left\{e^{t A}: t \in \mathbb{R}\right\}$ is a subgroup of $H$.

Corollary 5.15. $\operatorname{Lie}(H)=\left\{A \in M_{n \times n}: I+A \epsilon \in H\right\}$.
For example, let $H=S L_{n}=V(\operatorname{det}-1)$. The vector $I+A \epsilon \in H$ if and only if $\operatorname{det}(I+A \epsilon)=1$. Observe that

$$
I+A \epsilon=\left(\begin{array}{cccc}
1+a_{11} \epsilon & a_{12} \epsilon & \ldots & a_{1 n} \epsilon  \tag{5.9}\\
a_{21} \epsilon & 1+a_{22} \epsilon & \ldots & a_{2 n} \epsilon \\
\ldots & \ldots & \ddots & \ldots \\
a_{n 1} \epsilon & a_{n 2} \epsilon & \ldots & 1+a_{n n} \epsilon
\end{array}\right)
$$

so $\operatorname{det}(I+A \epsilon)=1+\left(a_{11}+\ldots+a_{n n}\right) \epsilon$, which is 1 if and only if $\operatorname{tr}(A)=0$. Therefore,

$$
\operatorname{Lie}\left(S l_{n}(\mathbb{F})\right)=\left\{A \in M_{n \times n}: \operatorname{tr}(A)=0\right\}=\mathfrak{s l}_{n} .
$$

Fact 5.14 is equivalent to saying $\operatorname{tr}(A)=0 \Leftrightarrow \operatorname{det}\left(e^{t A}\right)=1$ for all $t \in \mathbb{R}$. To prove the fact in this case we simply note that

$$
\operatorname{det}\left(e^{B}\right)=e^{\operatorname{tr}(B)}
$$

Definition 5.16. A Lie algebra $V$ over a field $\mathbb{F}$ is a vector space over $\mathbb{F}$ with a law of composition $V \times V \rightarrow V$ mapping $(a, b) \mapsto[a, b]$ such that
(1) $[$,$] is bilinear.$
(2) $[a, a]=0$.
(3) The Jacobi identity: $[a,[b, c]]+[b,[c, a]]+[c,[a, b]]=0$.

Exercise 5.17. For $\mathrm{sl}_{n}$, the bracket operation is defined by

$$
[A, B]=A B-B A
$$

Verify Properties $1,2,3$ hold for this law of composition on $\mathrm{sl}_{n}$.

## 6. Criteria for Smoothness of $X_{w}$ (APril 11,2007 )

Recall, $\mathfrak{g}=\operatorname{Lie}\left(S l_{n}(\mathbb{F})\right)=\left\{A \in M_{n \times n}: \operatorname{tr}(A)=0\right\}=\mathfrak{s l}_{n}$. Observe that $\mathfrak{g}$ has a vector space basis

$$
\begin{equation*}
\mathcal{B}=\left\{E_{i j}: i \neq j\right\} \cup\left\{H_{i}: 1 \leq i<n\right\}, \tag{6.1}
\end{equation*}
$$

where $E_{i j}$ is the matrix whose $i j^{\text {th }}$ entry is 1 and all other entries are 0 and $H_{i}$ is the matrix whose $i i^{\text {th }}$ entry is $1,(i+1, i+1)^{s t}$ entry is -1 and all other entries are 0 .

Exercise 6.1. Compute $\left[E_{i j}, E_{k l}\right],\left[E_{i j}, H_{k}\right]$, and $\left[H_{i}, H_{j}\right]$ to find which elements of $\mathcal{B}$ commute.

For $B$ the subgroup of upper triangular matrices in $S l_{n}$, Corollary 5.15 and Equation 6.1 implies $\mathfrak{b}=\operatorname{Lie}\left(B \cap S l_{n}\right)$ has a basis $\left\{E_{i j}: i<j\right\} \cup\left\{H_{i}: 1 \leq i<n\right\}$. Therefore the tangent space to $G / B$ at $I$ is equal to the tangent space to $S l_{n} /\left(B \cap S l_{n}\right)$ at $I$ which is equal to $\mathfrak{g} / \mathfrak{b}$. We have bijections

$$
\begin{equation*}
\mathfrak{g} / \mathfrak{b}=\operatorname{Span}\left\{E_{j i}: i<j\right\} \longleftrightarrow\left\{t_{i j} \in S_{n}\right\}=R \longleftrightarrow\left\{e_{j}-e_{i}: i<j\right\}:=\Delta_{+} . \tag{6.2}
\end{equation*}
$$

More generally, for $v \in S_{n}$,

$$
\begin{equation*}
v^{-1}(\mathfrak{g} / \mathfrak{b}) v=\operatorname{Span}\left\{v^{-1} E_{j i} v: i<j\right\}=\operatorname{Span}\left\{E_{v(j) v(i)}: j<i\right\} \tag{6.3}
\end{equation*}
$$

is the tangent space to $G / B$ at $v$. Note, $E_{j i} v=E_{j, v(i)}$ and $v^{-1} E_{j i}=E_{v(j), i}$ so $v^{-1} E_{j i} v=E_{v(j) v(i)}$.
Definition 6.2. Let $T_{w}(v)$ be the tangent space to $X_{w}$ at $v$.
Theorem 6.3. (Lakshmibai-Seshadri) We have that

$$
\begin{align*}
T_{w}(v) & =\operatorname{Span}\left\{E_{v(j) v(i)}: i<j \text { and } t_{v(j) v(i)} v \leq w\right\}  \tag{6.4}\\
& =\operatorname{Span}\left\{E_{v(j) v(i)}: i<j \text { and } v t_{i j} \leq w\right\} \tag{6.5}
\end{align*}
$$

Proof. The two sets are equivalent since $t_{v(j) v(i)} v=v t_{i j} . X_{w} \subseteq X_{w_{0}}=G / B$ implies $T_{w}(v) \subseteq T_{w_{0}}(v)$. Thus, we only need to identify which $E_{v(j) v(i)} \in T_{w}(v)$ i.e. which $I+E_{v(j) v(i)} \epsilon \in X_{w}$ by Corollary 5.15.

If $v$ is the identity, then

$$
r k\left(I+E_{j i} \epsilon\right)=r k\left(\begin{array}{ccccccc}
1 & \cdots & 0 & 0 & 0 & 0 & 0 \\
\vdots & \ddots & & & & & \\
0 & & 1_{i} & & & & \\
0 & & \vdots & \ddots & & & \\
0 & & \epsilon & \cdots & 1_{j} & & \\
0 & & & & & \ddots & \\
0 & & & & & & 1
\end{array}\right)=r k\left(\begin{array}{ccccccc}
1 & \cdots & 0 & 0 & 0 & 0 & 0 \\
\vdots & \ddots & & & & & \\
0 & & 0_{i} & \cdots & 1 & & \\
0 & & \vdots & \ddots & \vdots & & \\
0 & & 1 & \cdots & 0_{j} & & \\
0 & & & & & \ddots & \\
0 & & & & & & 1
\end{array}\right),
$$

which implies $r k\left(I+E_{j i} \epsilon\right)=r k_{\swarrow}\left(t_{j i}\right)$. Therefore $I+E_{j i} \epsilon \in X(w)$ if and only if $t_{j i} \leq w$ in Bruhat order.

Similarly, if $v \leq w$, say $v(j)=k$ and $v(i)=\ell$.

$$
r k_{\swarrow}\left(v+E_{j i} \epsilon\right)=r k_{\swarrow}\left(\begin{array}{cccc} 
& 1 & & \\
& & & \\
& & 1 \\
& \vdots & & \\
& \epsilon & \cdots & 1
\end{array}\right)=\left\{\begin{array}{ll}
r k_{\swarrow}(v) & \text { if } l(v)>l\left(t_{k l} v\right) \\
r k_{\swarrow}\left\langle\left(t_{k l} v\right)\right. & \text { if } l(v)<l\left(t_{k l} v\right)
\end{array}\right\} .
$$

Thus, $r k_{\swarrow}\left(v+E_{j i} \epsilon\right)=r k_{\swarrow} w$ if and only if both $v \leq w$ and $v t_{j i} \leq w$.
Corollary 6.4. For all $v \leq w$,
(1) $\operatorname{dim} T_{w}(v)=\left|\left\{t_{i j} \in R: v t_{i j} \leq w\right\}\right|$.
(2) $X_{w}$ is smooth at $v \Leftrightarrow \operatorname{dim} T_{w}(v)=l(w) \Leftrightarrow\left|\left\{t_{i j} \in R: v t_{i j} \leq w\right\}\right|=l(w)$.
(3) $\operatorname{dim} T_{w}(v) \geq l(w)$.
(4) $\left|\left\{t_{i j} \in R: v<v t_{i j} \leq w\right\}\right| \geq l(w)-l(v)$.

Exercise 6.5. Prove Property 3 in Corollary 6.4 using $\left|\left\{t_{i j} \in R: v t_{i j} \leq w\right\}\right|=\operatorname{dim} T_{w}(v)$.

Is $X(4231)$ smooth? No! There are 6 transpositions in the interval from the identity to 4231 in the Hasse diagram for Bruhat order, but $l(4231)=5$. We have that $v=2143$ is not a smooth point of $X_{4231}$ along with its lower order ideal in Bruhat order.
Fact 6.6. $\operatorname{Sing}(X(4231))=X(2143)$ and $\operatorname{Sing}(X(3412))=X(1324)$. (Memorize these patterns!) All other $X(w) \in S l_{4} / B$ are smooth.
Definition 6.7. The Bruhat graph of $w$ has vertices $v \leq w$ and edges between $u$ and $v$ if $u=v t_{i j}$ for some $i, j$.
Fact 6.8. The dimension of $T_{w}(v)$ is equal to the degree of $v$ in the Bruhat graph of $w$.
Geometric Interpretation: Say $v<v t_{i j}$. Then

$$
\left.\left.\begin{array}{rl}
L_{v}^{(i, j)} & =\left\{\left(\begin{array}{ccc}
1_{i} & & 1 \\
\vdots & & \\
& * & \cdots
\end{array}\right): 1_{j}\right. \\
1 & \\
&
\end{array}\right): 1_{k} \text { is on the } k^{t h} \text { row and } * \in \mathbb{C}\right\} \cup\left\{\left(\begin{array}{lll} 
& & \\
& & \\
& & 1
\end{array}\right)\right\}
$$

Example 6.9. Let $v=4231$. Then left multiplication by $T$ does the following:

$$
\left(\begin{array}{llll}
t_{1} & & & \\
& t_{2} & & \\
& & t_{3} & \\
& & & t_{4}
\end{array}\right)\left(\begin{array}{lll}
1 & & 1 \\
& * & 1
\end{array}\right)=\left(\begin{array}{llll} 
& & & t_{1} \\
& t_{2} & & \\
& * t_{3} & t_{3} & \\
t_{4} & & &
\end{array}\right) .
$$

Exercise 6.10. Show

$$
\left\{L_{v}^{(i, j)}, v \in S_{n}, i<j: v<v t_{i j}\right\}=\{1-\text { dimensional subvarieties of } X(w) \text { fixed by } T \text {-action }\} .
$$

Remark 6.11. A GKM space (Goresky-Kac-MacPherson) is any symplectic manifold with a torus $T$ action with a finite number of $T$ fixed points and a finite number of $T$ fixed 1-dimensional orbits. Then this data determines a graph just like the Bruhat graph. One asks, What can we say about the singularities and cohomology in terms of this graph?
Theorem 6.12. (Uber Smoothness Theorem) The following are equivalent for any $w \in S_{n}$ :
(1) $X_{w}$ is smooth.
(2) $X_{w}$ is smooth at id.
(3) $\left|\left\{t_{i j} \leq w\right\}\right|=l(w)$.
(4) The Bruhat graph for $w$ is regular.
(5) $P_{i d, w}(q)=1$.
(6) $P_{v, w}(q)=1$ for all $v \leq w$.
(7) $r_{w}(t)=\sum_{v \leq w} t^{l(v)}$ is palindromic.
(8) $r_{w}(t)=\prod_{j=1}^{k}\left(1+q+\ldots+q^{i_{j}}\right)$ for some $\left\{i_{1}, \ldots, i_{k}\right\} \in \mathbb{N}$.
(9) $w$ avoids 3412 and 4231.

Definition 6.13. $w$ avoids 3412 and 4231 means there is no $1 \leq i_{1} \leq i_{2} \leq i_{3} \leq i_{4} \leq n$ such that $f l\left(w_{i_{1}}, w_{i_{2}}, w_{i_{3}}, w_{i_{4}}\right) \in\{3412,4231\}$. Here for any set $\left\{d_{1}, d_{2}, \ldots, d_{k}\right\}$ of distinct real numbers, $f l\left(d_{1}, \ldots, d_{k}\right)=\left(v_{1}, \ldots, v_{k}\right)$ if $d_{i}$ is the $v_{i}^{t h}$ largest element among $\left\{d_{1}, \ldots, d_{k}\right\}$.

Open Question 6.14. Explain why $S_{4}$ is enough to determine all singularities.
If $G$ is semisimple, then smoothness( 1,2 of Uber Theorem) implies rational smoothness(4,5,6,7,8') implies combinatorial smoothness(3). In Type A, all three types of smoothness are equivalent. In Types D and E, smoothness is equivalent to rational smoothness (Petersen, Carell-Kuttler). All three flavors of smoothness can be characterized by pattern avoidance using root systems (Billey-Postnikov).
6.1. Kazhdan-Lusztig Polynomials in a Nutshell. The Kazhdan-Lusztig polynomial $P_{v, w}(q)$ is a polynomial in 1 variable with the following properties:
(1) $P_{v, w}(q)$ has constant term 1 if $v \leq w$.
(2) The degree of $P_{v, w}(q)$ is less than or equal to $\frac{l(w)-l(v)-1}{2}$.
(3) $\left.P_{w, w}(q)\right)=1$.
(4) $P_{v, w}(q) \neq 0 \Leftrightarrow v \leq w$.
(5) If $s=t_{i, i+1}, w s<w$, and $v s<v$, then

$$
\left.P_{v, w}(q)\right)=q P_{v s, w s}(q)+P_{v, w s}(q)-\sum_{z s<z<w s} \mu(z, w s) q^{\frac{l(w)-l(z)}{2}} P_{v, z}(q)
$$

where $\mu\left(x, y\right.$ is the coefficient of $q^{\frac{l(y)-l(x)-1}{2}}$ in $P_{x, y}(q)$.
(6) If $w s<w$ and $v s>v$ then $\left.P_{v, w}(q)\right)=P_{v s, w}(q)$.

Theorem 6.15. Let $I H(w)$ be the intersection cohomology sheaf of $X_{w}$ with respect to middle perversity. Then
(1) $P_{v, w}(q)=\sum \operatorname{dim}\left(I H^{2 j}\left(X_{w}\right)_{v}\right) q^{j}$.
(2) $\sum_{v \leq w} P_{v, w}(q) q^{l}(v)=\sum \operatorname{dim}\left(I H^{2 j}\left(X_{w}\right)\right) q^{j}$.
(3) $P_{x, w}(q)=1$ for every $x \leq w$ if and only if $X_{w}$ is rationally smooth.

In Theorem 6.15, the first statement implies the coefficients of $P_{v, w)}(q)$ are nonnegative.
Open Question 6.16. Find an all positive formula for $P_{v, w}(q)$.
The second statement implies $\sum_{v \leq w} P_{v, w}(q) q^{l(v)}$ is palindromic. We will take the third statement to be the definition for rational smoothness.

## 7. The "Uber" theorem (April 13, 2007)

Theorem 7.1 (Uber Theorem). The following are equivalent for $w \in S_{n}$ :
(1) $X_{w}$ is smooth.
(2) $X_{w}$ is smooth at $i d \in S_{n}$.
(3) $\#\left\{t_{i j}: t_{i j} \leq w\right\}=\ell(w)$.
(4) The Bruhat graph of $w$ is regular of degree $\ell(w)$.
(5) $P_{i d, w}(q)=1$.
(6) $P_{x, w}(q)=1$ for all $x \leq w$.
(7) $r_{w}(t)=\sum_{v \leq w} t^{\ell(v)}$ is palindromic.
(8) $r_{w}(t)=\prod_{j}\left(1+t+t^{2}+\cdots+t^{i_{j}}\right)$.
(9) $w$ avoids the patterns 3412 and 4231.

Today's lecture is devoted to the proof of the Uber Theorem.
Proof of the Uber Theorem: A number of the cases follow from results we have discussed in past lectures.
$(1) \Leftrightarrow(2)$ : follows since the singular locus is closed.
$\overline{(2) \Leftrightarrow(3)} \Leftrightarrow(4)$ : is a corollary of the Lakshmibai-Seshadari Theorem.
$(5) \Leftrightarrow(6)$ : follows from a theorem of Ron Irving.
Theorem 7.2 (Irving). Suppose

$$
\begin{gathered}
P_{v, w}(q)=1+a_{1} q+a_{2} q^{2}+\cdots+a_{k} q^{k}, \text { and } \\
P_{x, w}(q)=1+b_{1} q+b_{2} q^{2}+\cdots+b_{k} q^{k} .
\end{gathered}
$$

If $v \leq x \leq w$, then $a_{i} \geq b_{i}$ for all $1 \leq i \leq k$.
(4) $\Leftrightarrow(6) \Leftrightarrow(7)$ : follows from the work of Carroll-Peterson.

Example 7.3. We see that $r_{4231}=1+3 t+5 t^{2}+6 t^{3}+4 t^{4}+t^{5}$ is not palindromic.
However $r_{4213}=1+3 t+4 t^{2}+3 t^{3}+t^{4}$ is palindromic, and factors as $(1+t)^{2}\left(1+t+t^{2}\right)$.
Pop quiz: Which $X_{w}$ are smooth among the following $w$ ? Which characterization from the uber theorem seems the easiest to apply? (Class consensus: (9).)
$w=45678123 \in S_{8}$ ? Contains subsequence $4512, f l(4512)=3412$. Singular.
$w=7432651 \in S_{7}$ ? Contains subsequence 7351, $f l(7351)=4231$. Singular.
$w=7654321 \in S_{7}$ ? No ascents. Avoids 3412 and 4231. Smooth.
$w=7132645 \in S_{7}$ ? Avoids 3412 and 4231. Smooth.
$w=7263154 \in S_{7}$ ? Contains subsequence $7231, f l(7231)=4231$. Singular.
We still need to show that
(1) $w$ avoids the patterns 3412 and 4231 implies that $r_{w}(t)$ factors "nicely", that this then implies that $r_{w}(t)$ is palindromic, and that this implies that $X_{w}$ is smooth.
(2) $w$ contains one of the patterns 3412 and 4231 implies that $r_{w}(t)$ is not palindromic, and that therefore $r_{w}(t)$ does not factor "nicely".
Lemma 7.4. For $w \in S_{n}, r_{w}(t)$ factors (nicely) if one of the following hold:
Rule 1: If $w=\left[\ldots, n, w_{k+1}, \ldots, w_{n}\right]$ with $n=w_{k}>w_{k+1}>\cdots>w_{n}$, then

$$
r_{w}(t)=\left(1+t+t^{2}+\cdots+t^{n-k}\right) \cdot r_{w^{\prime}}(t)
$$

where $w^{\prime}=\left[w_{1}, \ldots, \widehat{w_{k}}, \ldots, w_{n}\right] \in S_{n-1}$ (or equivalently, $w=w^{\prime} s_{n-1} s_{n-2} \cdots s_{k}$, where $s_{j}=t_{j, j+1}$ is the adjacent transposition indexed by $j$ ).

Rule 2: Say that $w$ contains a consecutive sequence if for some $1 \leq j \leq n$ the $j+1$ integers $\{n-j, n-j+1, \ldots, n\}$ appear in the one-line notation of $w$ in decreasing order, with $w_{n}=n-j$. If $w=[\ldots, n, \ldots, n-1, \ldots, n-j]$ contains the consecutive sequence of length $j+1$, then

$$
r_{w}(t)=\left(1+t+t^{2}+\cdots+t^{j}\right) \cdot r_{w^{\prime}}(t)
$$

where $w^{\prime}=s_{n-1} s_{n-2} \cdots s_{n-j} w$.
Example 7.5 (Rule 2 example:). $w=7132645 \in S_{7}$ contains the consecutive sequence 765 .
$r_{w}(t)=\left(1+t+t^{2}\right) \cdot r_{w_{6}}(t)$, where $w_{6}(t)=613254$, which contains 654 .
$r_{w}(t)=\left(1+t+t^{2}\right)^{2} \cdot r_{w_{5}}(t)$, where $w_{5}(t)=51324$, which contains 54 .
$r_{w}(t)=\left(1+t+t^{2}\right)^{2}(1+t) \cdot r_{w_{4}}(t)$, where $w_{4}(t)=4132$, which contains 432.
$r_{w}(t)=\left(1+t+t^{2}\right)^{3}(1+t) \cdot r_{w_{3}}(t)$, where $w_{3}(t)=312$, which contains 32 .
$r_{w}(t)=\left(1+t+t^{2}\right)^{3}(1+t)^{2} \cdot r_{w_{2}}(t)$, where $w_{2}(t)=21$.
$r_{w}(t)=\left(1+t+t^{2}\right)^{3}(1+t)^{3}$.
Proof of Lemma 7.4. Note that $w_{j}=i \Leftrightarrow w^{-1}(i)=j$. Note also that $w$ satisfies the hypotheses for Rule 1 if and only if $w^{-1}$ satisfies the hypotheses for Rule 2. Why? By Chevalley's Criterion for Bruhat order and because $[i d, w]$ and $\left[i d, w^{-1}\right]$ are isomorphic posets as intervals in the Bruhat order. Thus it suffices to prove Rule 1. For the next part of the proof, you are asked to do the following exercise.
Exercise 7.6. Show that if $u, v \in S_{n}$ such that $u_{k}=v_{k}$, then $u \leq v$ if and only if

$$
f l\left(u_{1}, \ldots, \widehat{u_{k}}, \ldots, u_{n}\right) \leq f l\left(v_{1}, \ldots, \widehat{v_{k}}, \ldots, v_{n}\right) .
$$

Claim: If $n=w_{k}>w_{k+1}>\cdots>w_{n}$, then the interval $[i d, w]$ in the Bruhat order can be partitioned into $n-k+1$ blocks $\left\{B_{i}\right\}$ such that each block $B_{i}$ with the induced order is isomorphic to the interval [id, $\left.w^{\prime}\right]$, and the smallest element in each block has length $0,1,2, \ldots, n-k$ respectively.

Proof: Let $B_{i}=\left\{x \leq w: x_{i}=n\right\}$. By Exercise 7.6, the Bruhat order on $B_{i}$ induces a poset isomorphic to the interval $\left[i d, w^{\prime}\right]$, where $w^{\prime}=\left[w_{1}, \ldots, \widehat{w_{k}=n}, \ldots, w_{n}\right] \in S_{n-1}$. Since $n=w_{k}>$ $w_{k+1}>\cdots>w_{n}$, the maximum length element in $B_{n}$ is $w^{\prime}$, the maximum length element in $B_{n-1}$ is $w^{\prime} s_{n-1}$, etc., and in general, he maximum length element in $B_{k}$ is $w^{\prime} s_{n-1} s_{n-2} \cdots s_{k}$. Similarly, the minimal length element in $B_{i}$ in general is $[1,2, \ldots, i-1, n, i, i+1, \ldots, n-1]$, which is of length $n-i$. Therefore the map $\varphi:\{0,1,2, \ldots, n-k\} \times\left[i d, w^{\prime}\right] \rightarrow[i d, w]$ given by $(i, x) \mapsto$ $\left[x_{1}, x_{2}, \ldots, x_{i-1}, n, x_{i}, x_{i+1}, \ldots, x_{n-1}\right]$ is a bijection. Moreover, this bijection is length-preserving in the sense that $\ell(x)+i=\ell(\varphi(x))$, and the bijection respects Bruhat order. This establishes the claim.

From the claim it follows that $r_{w}(t)=\left(1+t+t^{2}+\cdots+t^{n-k}\right) \cdot r_{w^{\prime}}(t)$.
Exercise 7.7. Show that the interval $[i d, w]$ in the Bruhat order has a symmetric chain decomposition if $w$ is smooth.

Corollary 7.8. The interval $[i d, w]$ is rank-symmetric, rank-unimodal, and $K$-Sperner if $w$ is smooth.
Lemma 7.9. If $w \in S_{n}$ avoids the patterns 3412 and 4231, then $w$ satisfies the hypotheses of either Rule 1 or Rule 2 of Lemma 7.4. Moreover, the respective $w^{\prime} \in S_{n-1}$ also avoids the patterns 3412 and 4231.
Proof. Regarding the latter part of the statement of the lemma, if Rule 1 applies and $w$ avoids both patterns, then $w^{\prime}$ also avoids the patterns since subsequences of $w^{\prime}$ are also subsequences of $w$. If Rule 2 applies, then for any $\left\{i_{1}, i_{2}, i_{3}, i_{4}\right\} \subset\{1, \ldots, n-1\}$, we have $f l\left(w_{i_{1}}^{\prime}, w_{i_{2}}^{\prime}, w_{i_{3}}^{\prime}, w_{i_{4}}^{\prime}\right)=f l\left(w_{i_{1}}, w_{i_{2}}, w_{i_{3}}, w_{i_{4}}\right)$, so again, if $w$ avoids the patterns 3412 and 4231, then so does $w^{\prime}$. It remains to show that $w$ satisfies the hypotheses of either Rule 1 or Rule 2.

Say $w_{d}=\mathrm{n}$. We consider cases by the value of $d$.
Case $d=n$ or $n-1$. In this case, Rule 1 applies.
Case $d=n-2$. If $w_{n}=n-1$, then Rule 2 applies. If $w_{n-1}=n-1$, then Rule 1 applies. Otherwise $w_{i}=n-1$ for some $i<d$. It cannot be the case that $w_{n-1}<w_{n}$, for otherwise $f l\left(w_{i}=n-1, w_{d}=\right.$
$\left.n, w_{n-1}, w_{n}\right)=(3412)$ and so $w$ would contain the pattern 3412. Thus it must be that $w_{n-1}>w_{n}$, and Rule 1 applies.

Case $d<n-2$. If $w_{d}>w_{d+1}>\cdots>w_{n}$, then Rule 1 applies. Otherwise $w_{d}>w_{e}<w_{e+1}$ for some $d<e<n$. Without loss of generality we may assume that $e$ is the largest such index. If $w^{-1}(k)<d$ for any $w_{e+1}<k<n$, then $f l\left(k, n, w_{e}, w_{e+1}\right)=(3412)$, contradicting the assumption that $w$ avoids this pattern. In particular, $w^{-1}(n-1)>d$, say $w_{f}=n-1$. If $f=n$, then Rule 2 applies. Otherwise $f \leq n-1$.

Let $w_{g}=n-2$. If $g<d$, then either

- $w_{e+1} \neq n-1$, when $f l\left(n-2, n, w_{e}, w_{e+1}\right)=(3412)$, or
- $w_{e+1}=n-1$ and $w_{e}<w_{n}$, when $f l\left(n-2, n, w_{e}, w_{n}\right)=(3412)$, or
- $w_{e+1}=n-1$ and $w_{e}>w_{n}$, when $f l\left(n-2, w_{e}, n-1, w_{n}\right)=(4231)$.

If $d<g<f$, then $f l\left(n, n-2, n-1, w_{n}\right)=(4231)$. All these cases contradict our assumption that $w$ avoids these patterns. Thus we must have $g>f$, i.e. $w^{-1}(n-2)>w^{-1}(n-1)$.

Performing similar analysis for $w^{-1}(n-3), w^{-1}(n-4)$, etc. successively, in general comparing $w^{-1}(k)$ against the values $\left\{d=w^{-1}(n), w^{-1}(n-1), \ldots, w^{-1}(k+1)\right\}$, and knowing that $w$ avoids the patterns 3412 and 4231, shows that

$$
d<w^{-1}(n-1)<w^{-1}(n-2)<\cdots<w^{-1}(k)<w^{-1}(k-1)<\cdots
$$

until at some point when $w^{-1}(k)=n$ for some $k$. At this point it is seen that $w$ contains a consecutive sequence, and Rule 2 applies.

Lemma 7.10. If $w \in S_{n}$ such that $f l\left(w_{i_{1}}, w_{i_{2}}, w_{i_{3}}, w_{i_{4}}\right)=(3,4,1,2)$ or $(4,2,3,1)$, then $X_{w}$ is not smooth.

Proof. Let $w^{\prime}=f l\left(w_{i_{1}}, w_{i_{2}}, w_{i_{3}}, w_{i_{4}}\right) \in S_{4}$. Let $u^{\prime} \in S_{4}$ be such that Sing $X\left(w^{\prime}\right)=X\left(u^{\prime}\right)$. Then either $u^{\prime}=1324$ or $u^{\prime}=2143$. Then $\#\left\{t_{i j} \in S_{4}: u^{\prime} t_{i j} \leq w^{\prime}\right\}>\ell\left(w^{\prime}\right)$ by the Lakshmibai-Seshadari Theorem.

Let $u \in S_{n}$ such that $u_{i}=w_{i}$ for all $i \notin\left\{i_{1}, i_{2}, i_{3}, i_{4}\right\}$, and $f l\left(u_{i_{1}}, u_{i_{2}}, u_{i_{3}}, u_{i_{4}}\right)=u^{\prime}$. Define $D_{w}(x):=\left\{t_{i j} \in S_{n}: x t_{i j} \leq w\right\}$, and define $D_{w}^{k}(x):=\left\{t_{i j} \in D_{w}(x):\left|\left\{i, j, i_{1}, i_{2}, i_{3}, i_{4}\right\}\right|=k\right\}$ for $k=4,5,6$. Then $D_{w}(x)=D_{w}^{4}(x) \sqcup D_{w}^{5}(x) \sqcup D_{w}^{6}(x)$.

By Exercise 7.6, if $v, w \in S_{n}$ such that $v_{i}=w_{i}$, then $v \leq w$ if and only if

$$
\left[v_{1}, \ldots, \widehat{v_{i}}, \ldots, v_{n}\right) \leq\left[w_{1}, \ldots, \widehat{w_{i}}, \ldots, w_{n}\right]
$$

Therefore, by repeated application of this fact, to determine if $u t_{i j} \leq w$, we only need to compare on positions in $\left\{i, j, i_{1}, i_{2}, i_{3}, i_{4}\right\}$. Computer-run verifications show that for $w^{\prime \prime} \in S_{6}$ containing either of the patterns 3412 or 4231 and for $u^{\prime \prime}$ constructed from $w^{\prime \prime}$ as above,

$$
\begin{aligned}
\left|D_{w^{\prime \prime}}^{4}\left(u^{\prime \prime}\right)\right| & >\left|D_{w^{\prime \prime}}^{4}\left(w^{\prime \prime}\right)\right| \\
\left|D_{w^{\prime \prime}}^{5}\left(u^{\prime \prime}\right)\right| & \geq\left|D_{w^{\prime \prime}}^{5}\left(w^{\prime \prime}\right)\right| \\
\left|D_{w^{\prime \prime}}^{6}\left(u^{\prime \prime}\right)\right| & =\left|D_{w^{\prime \prime}}^{6}\left(w^{\prime \prime}\right)\right|
\end{aligned}
$$

and hence

$$
\left|D_{w^{\prime \prime}}\left(u^{\prime \prime}\right)\right|=\sum_{k=4,5,6}\left|D_{w^{\prime \prime}}^{k}\left(u^{\prime \prime}\right)\right|>\sum_{k=4,5,6}\left|D_{w^{\prime \prime}}^{k}\left(w^{\prime \prime}\right)\right|=\ell\left(w^{\prime \prime}\right) .
$$

Thus $u$ is a singular point of $X_{w}$.
This concludes our proof of the Uber Theorem.

One can generalize some conditions of the Uber Theorem for intervals in the Bruhat order poset:
(4): The Bruhat graph of $[v, w]$ is regular of degree $\ell(w)-\ell(v)$.
(5): $P_{v, w}(q)=1$.
(6'): $P_{x, w}(q)=1$ for all $v \leq x \leq w$.
The following theorem was proved around the year 2000-2001 by several teams of mathematicians: Billey-Warrington, Cortez, Kassel-Lascoux-Reutenauer, and Manivel.

Theorem 7.11. For $w \in S_{n}$, we have

$$
\operatorname{Sing}\left(X_{w}\right)=\bigcup_{v \in \operatorname{maxsing}(w)} X_{v} .
$$

Moreover, $v \in \operatorname{maxsing}(w)$ if and only if $X_{v}$ is an irreducible component of $\operatorname{Sing}\left(X_{w}\right)$. Also, $v \in$ $\operatorname{maxsing}(w)$ is obtained from $w$ (diagrammatically) in one of three ways: ... (diagrams to be provided).

## 8. Grassmannians (April 18, 2007)

Open Question 8.1. Characterize the maximal singular set of $X_{w} \subset G / B$ where $G=S O(2 n+1)$, $S O(2 n), S p(2 n)$ or any other semisimple Lie group.

Let $V$ be a vector space over some field $\mathbb{F}$ of dimension $n$, so that $V \simeq \mathbb{F}^{n}$.
Definition 8.2. $G(k, n)=\left\{k\right.$-dimensional subspaces of $\left.\mathbb{F}^{n}\right\}$.
Example 8.3. $G(1, n)=\{$ lines through the origin $\}=\mathbb{P}(V)$.
Example 8.4. $G(2,4)=\left\{\right.$ planes through the origin in $\left.\mathbb{F}^{4}\right\}=\left\{\right.$ affine lines in $\left.\mathbb{P}^{3}\right\}$.
Fix a basis for $V: e_{1}, \ldots, e_{n}$. Then a $k$-dimensional subspace can be represented by a matrix:
$k$-dimensional subspace $\longleftrightarrow \operatorname{span}\left\{\sum a_{i 1} e_{i}, \ldots, \sum a_{i k} e_{i}\right\} \longleftrightarrow\left(\begin{array}{cccc}a_{11} & a_{12} & \ldots & a_{1 k} \\ \vdots & \vdots & \vdots & \vdots \\ a_{n 1} & a_{n 2} & \ldots & a_{n k}\end{array}\right)$.
Let $\mathcal{M}_{G(k, n)}$ be the set of $n \times k$ matrices with rank $k$. Note that $G(k, n) \simeq G(n-k . n)$, since $A \in \mathcal{M}_{G(k, n)}$ corresponds to a linear map $A: \mathbb{F}^{n} \rightarrow \mathbb{F}^{n}$ given by $\left(v_{1}, \ldots, v_{n}\right) \mapsto\left(v_{1}, \ldots, v_{n}\right) A$ and the kernel of this map is an $n-k$-dimensional subspace.

Exercise 8.5. Show that ker $A=\operatorname{ker} B$ if and only if columns of $A$ span the same subspace as columns of $B$.

Note that rescaling columns and adding to any column a linear combination of other columns does not change the span. Therefore, we can choose a column echelon form as the canonical one:

$$
\begin{gathered}
\left(\begin{array}{lll}
3 & 4 & 7 \\
9 & 2 & 9 \\
6 & 0 & 7 \\
3 & 0 & 3
\end{array}\right) \xrightarrow[\text { bottom 1's }]{\text { rescale to get }}\left(\begin{array}{ccc}
1 & 2 & 7 \\
3 & 1 & 9 \\
2 & 0 & 7 \\
1 & 0 & 3
\end{array}\right) \xrightarrow[\text { right from 1's }]{\text { clear to the }}\left(\begin{array}{lll}
1 & 2 & 4 \\
3 & 1 & 0 \\
2 & 0 & 1 \\
1 & 0 & 0
\end{array}\right) \xrightarrow[\text { left from 1's }]{\text { clear to the }}\left(\begin{array}{ccc}
-5 & 2 & 4 \\
0 & 1 & 0 \\
2 & 0 & 1 \\
1 & 0 & 0
\end{array}\right) \\
\underset{\text { left from 1's }}{\text { clear to the }}\left(\begin{array}{ccc}
-13 & 2 & 4 \\
0 & 1 & 0 \\
0 & 0 & 1 \\
1 & 0 & 0
\end{array}\right) \xrightarrow[\text { columns }]{\text { rearange }}\left(\begin{array}{ccc}
2 & 4 & -13 \\
1 & 0 & 0 \\
0 & 1 & 0 \\
0 & 0 & 1
\end{array}\right)
\end{gathered}
$$

Note that all these operations can be realized as multiplication by some invertible matrices.

Claim 8.6. If $A \in \mathcal{M}_{G(k, n)}$, then there exists $g \in G l_{K}(\mathbb{F})$ such that $A g=B$ and $B$ is in the canonical form, i.e. it has the form

$$
\left(\begin{array}{ccc}
* & * & * \\
1 & 0 & 0 \\
0 & * & * \\
0 & 1 & 0 \\
0 & 0 & * \\
0 & 0 & 1
\end{array}\right)
$$

with zeros below each of the ones and zeros to the left and right of each one.
Let $i_{1}, \ldots, i_{k}$ be the indices of rows containing ones.
Claim 8.7. There is a bijection $G(k, n) \longleftrightarrow \mathcal{M}_{G(k, n)} /\left\langle A=B g, g \in G l_{k}\right\rangle$.
For $U \in G(k, n)$ define $\mathcal{M}(U)$ to be the corresponding matrix in the canonical form. Then columns of $\mathcal{M}(U)$ span $U$. For any $k$-subset $\left\{i_{1}<\ldots, i_{k}\right\} \subset\{1, \ldots, n\}=[n]$ define the Schubert cell

$$
\left.C_{\left\{i_{1}, i_{2}, \ldots, i_{k}\right\}}=\left\{U \in G(k, n) \text { such that } \mathcal{M}(U) \text { has the form }\left(\begin{array}{cccc}
1 & & & \\
& 1 & & \\
& & 1 & \\
& & & 1
\end{array}\right) \leftarrow i_{1}\right) \leftarrow i_{k}\right\}
$$

Then $G(k, n)=\bigcup_{\left\{i_{1}, i_{2}, \ldots, i_{k}\right\} \subset[n]} C_{\left\{i_{1}, i_{2}, \ldots, i_{k}\right\}}$. Schubert varieties in $G(k, n)$ are $X_{\left\{i_{1}, i_{2}, \ldots, i_{k}\right\}}=\overline{C_{\left\{i_{1}, i_{2}, \ldots, i_{k}\right\}}}$.
Claim 8.8. There exists a map $\phi: G l_{n} / B_{n} \rightarrow G(k, n)$ given by $A \mapsto \operatorname{span}\left\{\operatorname{col}_{1}(A), \ldots, \operatorname{col}_{k}(A)\right\}$.
Proof. Since $B_{n}$ adds previous columns to later ones, it does not change spans of this kind. Thus $\phi$ is well defined.

Define parabolic subgroups of $G l_{n}$ by

$$
P_{k, n-k}=\left\{\left(\begin{array}{c|c}
* \in G l_{k} & * \\
\hline 0 & * \in G l_{n-k}
\end{array}\right)\right\} .
$$

Fact 8.9. $G(k, n) \simeq G l_{n} / P_{k, n-k}$.
Follows from the fact that the map $\phi$ forgets everything but the span informations of the first $k$ columns.

Fact 8.10. Since we have $G l_{n} / B_{n} \rightarrow G l_{n} / P_{k, n-k} \simeq G(k, n)$, we get $H^{*}(G(k, n)) \hookrightarrow H^{*}\left(G l_{n} / B_{n}\right)$ on cohomology.

Note that the map $\phi$ is given by $\phi\left(X_{w}\right)=X_{\text {sort }\left\{w_{1}^{-1}, \ldots, w_{k}^{-1}\right\}}$. Starting with a matrix in $G(k, n)$ we can extend it in the following way without adding new $*$ 's and thus obtaining a bijection:

$$
\left(\begin{array}{cccc}
* & * & * & * \\
* & * & * & * \\
* & * & * & * \\
1 & 0 & 0 & 0 \\
0 & * & * & * \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & 1
\end{array}\right) \longleftrightarrow\left(\begin{array}{cccc|cccc}
* & * & * & * & 1 & 0 & 0 & 0 \\
* & * & * & * & 0 & 1 & 0 & 0 \\
* & * & * & * & 0 & 0 & 1 & 0 \\
1 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & * & * & * & 0 & 0 & 0 & 1 \\
\hline 0 & 1 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 1 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 1 & 0 & 0 & 0 & 0
\end{array}\right)
$$

note the essential set.
Definition 8.11. Let $w\left\{i_{1}, \ldots, i_{k}\right\}=\left\{i_{1}, \ldots, i_{k}, 1 \ldots, \hat{i_{1}}, \ldots, \hat{i_{2}}, \ldots, \hat{i_{k}}, \ldots, n\right\}$.
Warning 8.12. $X_{w\left\{i_{1}, \ldots, i_{k}\right\}}$ is not homeomorphic to $X_{\left\{i_{1}, \ldots, i_{k}\right\}}$.
In order to determine which Schubert cells $C_{\left\{j_{1}, \ldots, j_{k}\right\}}$ are in $X_{\left\{i_{1}, \ldots, i_{k}\right\}}$, it is convenient to look at partitions.
Lemma 8.13. There is a bijection between $k$-subsets of [ $n$ ] and partitions ( $\lambda_{1} \leq \lambda_{2} \leq \cdots \leq \lambda_{k}$ ) such that $\lambda_{k} \leq n-k$, given by $\left\{i_{1}, \ldots, i_{k}\right\} \longleftrightarrow\left(i_{1}-1, \ldots, i_{k}-k\right)$.

Example 8.14. $\{4,6,7,8\} \longleftrightarrow(3,4,4,4) \longleftrightarrow$


Let $|\lambda|=\sum \lambda_{i}$. Then $\operatorname{dim} C_{\lambda}=|\lambda|$. (We can index Schubert cells and varieties by $k$-subsets or partitions.)
Theorem 8.15. $X_{\mu} \subset X_{\nu}$ if and only if $\mu \subset \nu$ as Ferrer's diagrams, i.e. $\mu_{i} \leq \nu_{i}$ for $i=1, \ldots, k$.
This partial order is isomorphic to $[\Phi,(\underbrace{k, k, \ldots, k}_{n-k \text { times }})]$ in Young's lattice.
Example 8.16. For $G(3,5)$ we get


This is a rank symmetric self dual lattice (for any two elements there exists a unique minimal one covering both of them).

Now let's put 1's in decreasing order:

$$
\left(\begin{array}{llll|llll}
* & * & * & * & * & * & * & 1 \\
* & * & * & * & * & * & 1 & 0 \\
* & * & * & * & * & 1 & 0 & 0 \\
* & * & * & 1 & 0 & 0 & 0 & 0 \\
* & * & * & 0 & 1 & 0 & 0 & 0 \\
* & * & 1 & 0 & 0 & 0 & 0 & 0 \\
* & 1 & 0 & 0 & 0 & 0 & 0 & 0 \\
1 & 0 & 0 & 0 & 0 & 0 & 0 & 0
\end{array}\right) .
$$

Notice that the essential set is always in the column to the left from the line.
Definition 8.17. $W\left\{i_{1}, \ldots, i_{k}\right\}=$ "largest permutation in Bruhat order mapping to $\left\{i_{1}, \ldots, i_{k}\right\} "=$ $\left(n, n-1, \ldots, \hat{i_{k}}, \ldots, \hat{i_{k-1}}, \ldots, \hat{i_{1}}, \ldots, 1, i_{k}, i_{k-1}, \ldots, i_{1}\right)$.

We will show that equations for $X_{\left\{i_{1}, \ldots, i_{k}\right\}} \subset G(k, n)$ are the same as for $X_{W\left\{i_{1}, \ldots, i_{k}\right\}} \subset G / B$.
Fact 8.18. $X_{W\left\{i_{1}, \ldots, i_{k}\right\}} \rightarrow X_{\left\{i_{1}, \ldots, i_{k}\right\}}$ and this is a smooth morphism.
When is $X_{W\left\{i_{1}, \ldots, i_{k}\right\}}$ smooth? Note that 3412 cannot happen since $W\left\{i_{1}, \ldots, i_{k}\right\}$ is a concatenation of two decreasing sequences. But our example has 4231:

$$
\left(\begin{array}{cccc|cccc}
* & * & * & * & * & * & * & 1 \\
* & * & * & * & * & * & 1 & 0 \\
* & * & * & * & * & 1 & 0 & 0 \\
* & * & * & 1 & 0 & 0 & 0 & 0 \\
* & * & * & 0 & 0 & 0 & 0 & 0 \\
* & * & 1 & 0 & 0 & 0 & 0 & 0 \\
* & 1 & 0 & 0 & 0 & 0 & 0 & 0 \\
1 & 0 & 0 & 0 & 0 & 0 & 0 & 0
\end{array}\right)
$$

where * correspond to inner corners of the partition.
Claim 8.19. If 4231 exists in $W\left\{i_{1}, \ldots, i_{k}\right\}$, then there are at least 2 gaps in $\left\{i_{1}, \ldots, i_{k}\right\}$ and Ferrer's diagram for the corresponding partition $\lambda$ has internal corners like $\square$
Theorem 8.20. $X_{\lambda} \subset G(k, n)$ is smooth if and only if $\lambda$ is a rectangle (i.e. does not have internal corners).

If $\lambda$ is a rectangle $\left(i^{m}\right)$, then $X_{\lambda}=G(i, m+i)$.
Exercise 8.21. What is the set $E_{\mu}$ such that $\operatorname{sing}\left(X_{\lambda}\right)=\bigcup_{\mu \in E_{\lambda}} X_{\mu}$ ?
Now let's consider equations for $X_{\lambda}$ and Plücker coordinates for $G_{\mathbb{R}}(2,4)=\left\{2\right.$-planes in $\left.\mathbb{R}^{4}\right\}$, i.e.

$$
\left(\begin{array}{ll}
a_{11} & a_{12} \\
a_{21} & a_{22} \\
a_{31} & a_{32} \\
a_{41} & a_{42}
\end{array}\right) / G l_{2} .
$$

Let

$$
P_{i j}(X)=\operatorname{det}\left|\begin{array}{ll}
X_{i 1} & X_{i 2} \\
X_{j 1} & X_{j 2}
\end{array}\right| .
$$

Let

$$
P(A)=\left(P_{12}(A), P_{13}(A), P_{14}(A), P_{23}(A), P_{24}(A), P_{34}(A)\right) .
$$

Then $P(A) \in \mathbb{P}^{5}$ for any $A \in \mathcal{M}_{G(k, n)}$. Now if $A=B g$, then $P_{i j}(A)=P_{i j}(B) \cdot \operatorname{det}(g)$, so $P(A)=$ $P(B)$ in $\mathbb{P}^{5}$ if and only if the span of columns of $A$ is equal to the span of columns of $B$.

There is the following relation between Plücker coordinates, which follows from Sylvester's Lemma: $P_{12} P_{34}-P_{13} P_{24}+P_{23} P_{14}=0$.

Theorem 8.22. $G(2,4)=V\left(P_{12} P_{34}-P_{13} P_{24}+P_{23} P_{14}\right)$ in $\mathbb{P}^{5}$.
Why they are equal? Look for the largest $(i, j)$ in the lexicographical order such that $p \in V\left(P_{12} P_{34}-\right.$ $\left.P_{13} P_{24}+P_{23} P_{14}\right)$ does not vanish and rescale the coordinates so that $P_{i j}=1$. Put 1's into $(i, 1)$ and $(j, 2)$. Put 0 's below them and into $(i, 2)$. For example, for $(i . j)=(2,4)$ we get

$$
\left(\begin{array}{cc}
* & * \\
1 & 0 \\
0 & * \\
0 & 1
\end{array}\right) .
$$

Now all $a_{r s}$ in this matrix can be recovered: $a_{r 1}=P_{i r}$ and $a_{r 2}=P_{r j}$.
In general we have

$$
\begin{gathered}
P_{i_{1} \ldots i_{k}}(X)=\operatorname{det}\left|\begin{array}{ccc}
X_{i_{1} 1} & \ldots & X_{i_{1} k} \\
\vdots & \vdots & \vdots \\
X_{i_{k} 1} & \ldots & X_{i_{k} k}
\end{array}\right|, \\
P(A)=\left(P_{i_{1} \ldots i_{k}}(A):\left\{i_{1}<\cdots<i_{k}\right\} \subset[n]\right) \in \mathbb{P}^{\binom{n}{k}-1} .
\end{gathered}
$$

Theorem 8.23. $G(k, n) \simeq V$ (Plücker equations) $\left.\subset \mathbb{P}^{n} \begin{array}{l}n \\ k\end{array}\right)-1$, where equations are constructed in the following way: for each $d \in[k]$

$$
P_{i_{1} \ldots i_{k}} P_{j_{1} \ldots j_{k}}-\sum P_{i_{1}^{\prime} \ldots i_{k}^{\prime}} P_{j_{1}^{\prime} \ldots j_{k}^{\prime}}
$$

with the sum over all $\left(i^{\prime}, j^{\prime}\right)$ obtained from $(i, j)$ by picking $1 \leq q_{1}<q_{2}<\cdots<q_{d} \leq k$ positions and exchanging $i_{q_{s}}$ with $j_{s}$ for $s \in[d]$ :

$$
\begin{aligned}
i^{\prime} & =\left(i_{1}, \ldots, i_{q_{1}-1}, j_{1}, i_{q_{1}+1}, \ldots, i_{q_{2}-1}, j_{2}, i_{q_{2}+1}, \ldots, i_{k}\right) \\
j^{\prime} & =\left(i_{q_{1}}, i_{q_{2}}, \ldots, i_{q_{d}}, j_{d+1}, \ldots, j_{k}\right) .
\end{aligned}
$$

Exercise 8.24. Compare equations for $X_{\left\{i_{1}, \ldots, i_{k}\right\}}$ to the equations for the essential set for $X_{W\left\{i_{1}, \ldots, i_{k}\right\}}$.

## 9. Partial Flags, Intersecting and Chow cohomology (April 20, 2007)

Definition 9.1. For $0<d_{1}<d_{2}<\cdots<d_{p} \leq n$, define the partial flag manifold $F l\left(n, d_{1}, d_{2}, \cdots, d_{p}\right)=$ $G l_{n} / P \approx\left\{F_{\bullet}=F_{d_{1}} \subset F_{d_{2}} \subset \cdots \subset F d_{p}=\mathbb{F}^{n}\right\}$, with $\operatorname{dimF}_{d_{i}}=d_{i}$. Where $P$ is the set of upper triangular block matrices with block widths (and heights) $d_{1}, d_{2}-d_{1}, \cdots, n-d_{p}$.

In the Grassmannian we have the bijections:

$$
\begin{aligned}
\{K-\text { sheets }\} & \Leftrightarrow\{\text { partitions in a } k \times(n-k) \text { box }\} \\
& \Leftrightarrow S_{k} \times S_{n-k} \backslash S_{n} \text { mod on left } \\
& \Leftrightarrow \text { max length coset representative in } S_{k} \times S_{n-k} \backslash S_{n} \\
& \Leftrightarrow \text { min length coset representative in } S_{k} \times S_{n-k} \backslash S_{n}
\end{aligned}
$$

Thus $\left\{i_{1}, i_{2}, \cdots, i_{k}\right\} \Leftrightarrow W\left\{i_{1}, i_{2}, \cdots, i_{k}\right\} \Leftrightarrow w\left\{i_{1}, i_{2}, \cdots, i_{k}\right\}$. So for smoothness, $X_{\left\{j_{1}, j_{2}, \cdots, j_{n}\right\}} \subset$ $X_{\left\{i_{1}, i_{2}, \cdots, i_{n}\right\}}$ if and only if $w\left\{j_{1}, \cdots, j_{n}\right\} \leq w\left\{i_{1}, \cdots, i_{n}\right\}$ and hence $X_{\left\{i_{1}, i_{2}, \cdots, i_{n}\right\}}$ is smooth if and only if $X_{W\left\{i_{1}, i_{2}, \cdots, i_{n}\right\}}$ is.

The $T$ fixed points in $G / P$ are permutations in $S_{I} \backslash S_{n}$.
Definition 9.2. For the cosets of $S_{I} \backslash S_{n}$, we define $\left(S_{n}\right)^{I}$ to be the max length coset representative and $\left(S_{n}\right)_{I}$ to be the min length coset representative.

Then the Schubert varieties in $G / P$ are indexed by the cosets of $S_{I} \backslash S_{n}$, and thus by $\left(S_{n}\right)^{I}$ or $\left(S_{n}\right)_{I}$. So looking at the map $\varphi: G / B \rightarrow G / P$, with $X_{\min \left(S_{I} w\right)} \rightarrow X_{S_{I} w}$ and $X_{\max \left(S_{I} w\right)} \rightarrow X_{S_{I} w}$. So we have $\min \left(S_{I} u\right) \leq \min \left(S_{I} w\right)$ if and only if $X_{S_{I} u} \subset X_{S_{I} w}$ and $X_{S_{I} w}$ is smooth if and only if $X_{\max \left(S_{I} w\right)}$ is smooth.

Now for intersecting Schubert varieties, done in the Grassmannian but can be generalized. Consider the question, how many lines meet four given lines in $\mathbb{R}^{3}$ ? We can view $\mathbb{R}^{3}$ as an open subset of $\mathbb{P}^{3}=$ $\mathbb{R}^{4} /(1, \cdots, 1)$. So fix a line in $\mathbb{R}^{3}$, say the line, $\overline{e_{1} e_{2}}$, not containing the origin. This determines a 2 dimensional subspace of $\mathbb{R}^{4}$, $\operatorname{span}\left\{e_{1}, e_{2}\right\}$. Then extend $\left\{e_{1}, e_{2}\right\}$ to a basis $\left\{e_{1}, e_{2}, e_{3}, e_{4}\right\}$ for $\mathbb{R}^{4}$. So lines meeting $\left.\overline{e_{1} e_{2}}=\left\{\begin{array}{cc}* & * \\ 1 & 0 \\ 0 & * \\ 0 & *\end{array}\right]\right\}=X_{\{2,4\}} \subset G(2,4)$. That is $X_{(1,2)}$ in partition notation, note that this is with respect to the ordered basis, $\left\{e_{1}, e_{2}, e_{3}, e_{4}\right\}$. So for lines intersecting four given lines we look at $X_{(1,2)}^{G}(E \bullet) \cap X_{(1,2)}^{G}(F \bullet) \cap X_{(1,2)}^{G}(G \bullet) \cap X_{(1,2)}^{G}(H \bullet)$. One method of doing this is to solve equations!

Another method is Intersection theory.
Definition 9.3. Let $X$ be any variety.
(1) Two subvarieties $U$ and $V$ in $X$ meet transversally if $U \cap V=\bigcup Z_{i}$, where the $Z_{i}$ 's are the irreducible components, then for each $i$ and each point $z$ in an open sets of $Z_{i} \cap U \bigcap V$ we have $T_{z}\left(Z_{i}\right)=T_{z}(U) \cap T_{z}(V)$.
(2) The Chow ring of $X$, denoted $A^{\bullet}(X)$, is the formal sum of $\{[V]: V$ is a subvariety of $X\} / \sim$, where $[U] \equiv[V]$ if they are rationally equivalent. With multiplication given by $[U][V]=[U \cap$ $\left.V^{\prime}\right]$, with $[V] \sim\left[V^{\prime}\right]$ and $V^{\prime}$ meets $V$ transversely,

Definition 9.4. If for each $z_{i}, \operatorname{Codim}\left(z_{i}\right)=\operatorname{Codim}(v)+\operatorname{Codim}(w)$ then $V \cup W$ is proper.
If $V \cup W$ is proper, then intersection theory says that $[U][V]=\Sigma m_{i}\left[Z_{i}\right]$ and the $m_{i}$ 's are non-negative integers.

Theorem 9.5. If $X$ has a cell decomposition, $X=\bigcup C_{i}$ then $A^{\bullet}(X)=($ linear $)$ span $\left\{\overline{\left[C_{i}\right]}: i \in E_{x}\right\}$.

Theorem 9.6. $A^{\bullet}\left(X_{w}\right)=H^{*}\left(X_{w}, \mathbb{Z}\right)=($ linear $) \operatorname{span}\left\{\left[X_{v}\right]: v \leq w\right\}$.

Definition 9.7. Define the Poincare polynomial to be $\sum_{i \in E_{x}} t^{\operatorname{codim}\left(\overline{C_{i}}\right)}=\sum_{i \in E_{x}} \operatorname{dim}\left(A^{i}(X)\right) i^{i}$. eg: $\sum_{\lambda \subset} t_{k \times(n-k) b o x} t^{|\lambda|}$ is the Poincare polynomial of $G(k, n)$.

On the $G(k, n)$, we have the map $\lambda \rightarrow \operatorname{dual}(\lambda)$. This map gives the Poincare duality on $H^{*}(G(k, n))$. On $G / B$, we have $\left[X_{w}\right] \rightarrow\left[X_{w w_{0}}\right]$. But since $X_{u} \subset X_{w}$ implies $X_{w w_{0}} \subset X_{u w_{0}}$, so we only get duality if $X_{w}$ is smooth.

Open Problem (Reiner) - Which smooth varieties in $G / B$ exhibit Poincare duality? That is, when is $\{v \leq w\}$ self dual, with a nice combinatorial map?

The ring structure on $A^{\bullet}(X)$ is given by $\left[X_{\lambda}^{G}\right]\left[X_{\mu}^{G}\right]=\left[X_{\lambda}^{G}\left(E_{\bullet}\right) \cap X_{\mu}^{G}\left(F_{\bullet}\right)\right]$ where $F_{\bullet} \in X_{w_{0}}^{F l}\left(E_{\bullet}\right)$, so we could take $F_{\bullet}=\left\{<e_{1}>\subset<e_{1}, e_{2}>\subset \cdots \subset<e_{1}, \ldots, e_{n}>\right\}$. Notationally, denote $X_{\mu}\left(F_{\bullet}\right)=\widetilde{X_{\mu}}$. Then we get the matrix form by canceling up, left and right, instead of down, left and right.

Proposition 9.8. If $\operatorname{cod}(\mu)=\operatorname{cod}(\lambda)=k(n-k)=\operatorname{dim} G(k, n)$ then $\left[X_{\lambda}\right]\left[X_{\mu}\right]$ is $1\left[X_{\emptyset}\right]$ when $\lambda=$ dual ( $\mu$ ) and 0 otherwise.

Proof. First consider $\mu=\operatorname{dual}(\lambda)$, then

$$
\begin{aligned}
& {\left[X_{\lambda}\right]\left[X_{\mu}\right] }=X_{\lambda}\left(E_{\bullet}\right) \cap \tilde{X}_{\text {dual }(\lambda)} \\
&\left.=\left\{\begin{array}{ll}
* & * \\
1 & 0 \\
0 & * \\
0 & 1
\end{array}\right]\right\} \\
&\left.\left(\begin{array}{ll}
0 & 0 \\
1 & 0 \\
* & 0 \\
0 & 1
\end{array}\right]\right\} \\
&\left.=\left\{\begin{array}{ll}
0 & 0 \\
1 & 0 \\
0 & 0 \\
0 & 1
\end{array}\right]\right\} \text { (as all the stars are canceled by 0's) } \\
&=\left[X_{\emptyset}\right] \text { for some basis. }
\end{aligned}
$$

Now if $\mu \neq \operatorname{dual}(\lambda)$ then a 1 of one of the matrices will be a 0 in the other, so there intersection will be empty. Hence $X_{\lambda}\left(E_{\bullet}\right) \cap \widetilde{X}_{\mu}=\emptyset$.

Trick - To find the $C_{\lambda \mu}^{\nu}$ 's in $\left[X_{\lambda}\right]\left[X_{\mu}\right]=\sum_{\nu} C_{\lambda \mu}^{\nu}\left[X_{\nu}\right]$. Multiple both sides by $\left[X_{\text {dual }(\nu)}\right]$ and expand, getting $C_{\lambda \mu}^{\nu}\left[X_{\emptyset}\right]=\left[X_{\text {dual }(\nu)}\right]\left[X_{\lambda}\right]\left[X_{\mu}\right]=\left[X_{\text {dual }(\nu)}\left(E_{\bullet}\right) \cap X_{\lambda}\left(F_{\bullet}\right) \cap X_{\mu}\left(H_{\bullet}\right)\right]$. So $C_{\lambda \mu}^{\nu}$ is the number of points in $X_{\text {dual }(\nu)}\left(E_{\bullet}\right) \cap X_{\lambda}\left(F_{\bullet}\right) \cap X_{\mu}\left(H_{\bullet}\right)$.

Proposition 9.9. $\left[X_{\lambda}\right]\left[X_{\hbar}\right]=\sum_{\nu \leq \lambda ;|\nu|=|\lambda-1|}\left[X_{\nu}\right]$, where $\hbar$ has the shape of the $k \times(n-k)$ box minus the bottom left corner.

Proof. To see this we count the number of points in $\widetilde{X}_{\text {dual }(\nu)} \cap X_{\lambda}\left(E_{\bullet}\right) \cap X_{\hbar}\left(H_{\bullet}\right)$. But $\widetilde{X}_{\text {dual }(\nu)} \cap X_{\lambda}=\emptyset$ if we have any gaps between $\nu$ and $\lambda$, thus $\operatorname{dual}(\lambda) \subset \operatorname{dual}(\nu)$ and they differ by only one box (or else $C_{\lambda \hbar}^{\nu}=0$ ). So assume $\operatorname{dual}(\nu)=\operatorname{dual}(\lambda)+1$ box. So just looking at their differences, we see $X_{\lambda}=$ $\binom{1}{\star}$ and $\widetilde{X}_{d u a l(\nu)}=\binom{\star}{1}$, where $\star \neq 0$. Then $\widetilde{X}_{d u a l(\nu)} \cap X_{\lambda}=\binom{\star}{\star}$. Now, taking $Y=\left\{V \in G(k, n): V\right.$ contains $<e_{i_{1}}, e_{i_{2}}, \cdots, \widehat{e_{i_{j}}}, \cdots, e_{i_{k}}>$ and intersects span $<$ $e_{i_{j}}, e_{i_{j+1}}>$ in 1 dimension $\}$. Then $Y \cap X_{\hbar}\left(H_{\bullet}\right)=\left\{V \in G(k, n): V\right.$ meets span $<h_{1}, \cdots, h_{n-k}>$ in 1 dimension $\}$. Generally, span $<e_{i_{1}}, \cdots, \widehat{e_{i_{j}}}, \cdots, e_{i_{k}}>\cap$ span $<h_{1}, \cdots, h_{n-k}>=\{0\}$. Thus $<e_{i_{j}}, \cdots, e_{i_{j+1}}>\cap<h_{1}, \cdots, h_{n-k}>$ is 1 dimensional. So $Y \cap X_{\hbar}\left(H_{\bullet}\right)=\{$ one point $\}=X_{\emptyset}\left(G_{\bullet}\right)$ and thus $C_{\lambda \hbar}^{\nu}=1$.

Wow! This is the same as the Schur functions! $S_{\lambda} S_{b o x}=\sum S_{\lambda+w x}$ and $S_{\lambda}\left(x_{1}, \cdots, x_{k}\right)=$ $\sum_{T \text { fillingof } \lambda} x^{T}$, with $S_{\lambda}\left(x_{1}, \cdots, x_{k}\right)=0$ if $\lambda$ has more then $k$ rows.

Theorem 9.10. Multiplication in $A^{\bullet}(G(k, n))$ (or in $A^{\bullet}(G / B)$ ) is determined by multiplication by the unique codimension 1 subvariety (or the $n-1$ codimension 1 subvarieties in $A^{\bullet}(G / B)$ ), along with stability.

Corollary 9.11. $\operatorname{Pic}(G(k, n))=\mathbb{Z}$.

Theorem 9.12. We have a bijection (as rings) between $A^{\bullet}(G(k, n)) \leftrightarrow \operatorname{span}\left\{S_{\lambda}\right.$ : $\lambda$ is contained in a $k \times$ $(n-k)$ box $\}$. Given by $\left[X_{\lambda}\right] \rightarrow \mathcal{S}_{\text {dual }(\lambda) \text {. }}$.

Now back to the line intersecting four lines problem. As we say the intersection of 4 lines is $\left[X_{(1,2)}\right]^{4}$. So we are looking for the coefficient of $\left[X_{\emptyset}\right]=S_{(2,2)}$ in $\left(S_{(1,0)}\right)^{4}$. Then we see

$$
\begin{aligned}
S_{(1,0)} \times S_{(1,0)} \times S_{(1,0)} \times S_{(1,0)} & =S_{(1,0)} \times S_{(1,0)} \times\left(S_{(1,1)}+S_{(2,0)}\right) \\
& =S_{(1,0)} \times\left(S_{(2,1)}+S_{(2,1)}\right) \\
& =S_{(1,0)} \times 2 S_{(2,1)} \\
& =2 S_{(2,2)} .
\end{aligned}
$$

So we find out that the answer is 2 .

## 10. Vakil's Geometric Littlewood-Richardson Rule (April 25, 2007)

### 10.1. From last time.

(1) If $\operatorname{codim} \mathrm{X}_{\lambda}+\operatorname{codim} \mathrm{X}_{\mu}=\operatorname{dimG}(\mathrm{k}, \mathrm{n})$ then

$$
\left[X_{\lambda}\right]\left[X_{\mu}\right]=\left\{\begin{array}{c}
{\left[X_{\mathrm{pt}}\right], \lambda=\operatorname{dual}(\mu)} \\
0 \text { otherwise }
\end{array}\right.
$$

(2) In general, we have

$$
\left[X_{\lambda}\right]\left[X_{\mu}\right]=\sum C_{\lambda \mu}^{\nu}\left[X_{\nu}\right]
$$

where this sum is taken over all $\nu$ with $\operatorname{codim} \mathrm{X}_{\nu}=\operatorname{codim} \mathrm{X}_{\lambda}+\operatorname{codim} \mathrm{X}_{\mu}$,

$$
\left[X_{\text {dual } \nu}\right]\left[X_{\lambda}\right]\left[X_{\mu}\right]=C_{\lambda \mu}^{\nu}\left[X_{\mathrm{pt.}}\right],
$$

and

$$
C_{\lambda \mu}^{\nu}=\#\left(X_{\text {dual } \nu}(E .) \cap X_{\lambda}(F .) \cap X_{\mu}(G .)\right.
$$

when these are all in transverse position. Hence the coefficients $C^{\nu} \lambda \mu$ are nonnegative integers.
(3) There is a ring isomorphism $A \cdot G(k, n) \rightarrow \mathbb{C}\left[x_{1}, \ldots, x_{k}\right]^{S_{k}} /\left(S_{\lambda} \mid \lambda_{k}>n-k\right)$ given by $\left[X_{\lambda}\right] \mapsto$ $S_{\text {dual }}$. The proof of this theorem uses formulae of Giambelli, Jacobi, and Trudi. Given the elementary symmetric functions $e_{i}\left(x_{1}, \ldots, x_{k}\right)=\sum_{j_{1}<\ldots<j_{i}} x_{j_{1}} \cdots x_{j_{i}}$, we have

$$
S_{\lambda}=\operatorname{det}\left[\begin{array}{cccc}
e_{\lambda_{1}} & e_{\lambda_{2}} & e_{\lambda_{3}} & \cdots \\
1 & e_{\lambda_{2}} & e_{\lambda_{3}} & \cdots \\
0 & 1 & e_{\lambda_{3}} & \cdots \\
0 & 0 & 1 & \cdots \\
\vdots & \vdots & \vdots & \ddots
\end{array}\right]
$$

Recall also the Pieri Formula,

$$
S_{\lambda} S_{1^{j}}=\sum S_{\lambda^{\prime}}
$$

where this sum is over all $\lambda^{\prime}$ containing $\lambda$ such that $\left|\lambda^{\prime}\right|=|\lambda|+j$ and $\lambda^{\prime}-\lambda$ contains no two boxes in the same row.
(4) $S_{\lambda} S_{\mu}=\sum C_{\lambda \mu}^{\nu}$ similarly, where the $C_{\lambda \mu}^{\nu}$ are the Littlewood-Richardson coefficients.
(5) These coefficients are traditionally computed via reverse lattice fillings: one places the upper leftmost corner of the tableau for $\mu$ just beneath and to the right of the lower rightmost corner of that for $\lambda$ and fills the resulting array from bottom to top in all possible ways, strictly increasing as one ascends in the columns and weakly increasing in the rows.
(6) The first four results also hold for the flag manifolds $F l_{n}$ as well as the Grassmannians, where the third becomes $A \cdot F l_{n} \simeq \mathbb{C}\left[x_{1}, \ldots, x_{n}\right] /\left(S_{\lambda} \mid \lambda \neq 0\right)$ with the isomorphism being given by $\left[X_{w}\right] \mapsto \mathcal{S}_{w w_{0}}$ where these last are the Schubert polynomials. However, there does not exist a combinatorial technique like reverse lattice fillings to compute the Littlewood-Richardson coefficients in this case. Coskun is working on this currently.
10.2. Vakil's Checkerboard. Vakil approaches Littlewood-Richardson rules for the Grassmannian geometrically, via a series of degenerations indexed by the moves of a checker game. Index Schubert varieties by $k$-subsets of $[n]$. Given such subsets $A, B, C$, we seek a rule to compute the integer $C_{A B}^{C}$, which is the number of irreducible components of $X_{A}\left(M_{.}\right) \cap X_{B}(F)$ rationally equivalent to $X_{C}(F$. where $F$. is some fixed frame and $M$. is a moving frame. This rule will take the form of a game involving black and white checkers on an $n \times n$ board.

The game will define a family of sets $Y_{0} \bullet \subset G(k, n)$ indexed by arrangements $\circ$ of white checkers and • of black checkers on the $n \times n$ board, which will change as the checkers move. These should have the properties that:
(1) $X_{A}\left(M_{.}\right) \cap X_{B}(F)=.\overline{Y_{\circ_{A B} \bullet_{\text {init }}}}$
(2) There is a sequence of degenerations forming a tree such that

$$
\left[\overline{Y_{O_{A B} \bullet}}\right]=\sum_{\text {leaves }}\left[\overline{Y_{O_{C}} \bullet \text { final }}\right]
$$

and at each stage of the degeneration we have only three possibilities,

$$
\left[\overline{Y_{\bullet \bullet}}\right]=\left\{\begin{array}{l}
{\left[\overline{Y_{0 \bullet}} \overline{\bullet_{\text {next }}}\right]} \\
{\left[\overline{Y_{s_{\text {step }}} \bullet_{\text {next }}}\right]} \\
{\left[\overline{Y_{o \bullet_{\text {next }}}}\right]+\left[\overline{Y_{o_{\text {step }}} \bullet_{\text {next }}}\right]}
\end{array}\right.
$$

(3) $C_{A B}^{C}$ equals the number of steps until ${ }^{\circ} C$.
10.3. Black Checkers. Black checkers encode the position of the moving frame $M$. with respect to $F$. in a table of dimensions. Initially $F .=<e_{1}, \ldots, e_{n}>$ and $M .=<e_{n}, \ldots, e_{1}>$ and the number of black checkers NW (inclusive) of $(i, j)$ is $\operatorname{dimM}_{\mathrm{i}} \cap \mathrm{F}_{\mathrm{j}}$.


This example corresponds to a flag $M$. whose marked line meets the marked line $F_{1}$ only at the point $F_{0}$ and whose marked point misses $F_{1}$ entirely.
10.4. White Checkers. Fix $M$. and $F$. and some $V \in G(k, n)$. The white checkers encode the positions of $F$. and $M$. relative to $V: \operatorname{dimV} \cap \mathrm{M}_{\mathrm{i}} \cap \mathrm{F}_{\mathrm{j}}$ equals the number of white checkers NW (inclusive) of $(i, j)$. Thus for the same $M$. and $F$. as above, an arrangement of white checkers

corresponds to a $V$ which intersects both marked lines and passes through the marked point $M_{0}$ but not $F_{0}$.
10.5. Details. On any checkerboard, let $w(i, j)$ denote the number of white checkers NW (inclusive) of $(i, j)$, and similarly $b(i, j)$ for black checkers.

$$
X_{\circ} \bullet\left\{\left(V, M_{.}, V .\right) \in G(k, n) \times F l_{n} \times F l_{n} \mid \operatorname{dimM}_{\mathrm{i}} \cap \mathrm{~F}_{\mathrm{j}}=\mathrm{b}(\mathrm{i}, \mathrm{j}), \operatorname{dimV} \cap \mathrm{M}_{\mathrm{i}} \cap \mathrm{~F}_{\mathrm{j}}=\mathrm{w}(\mathrm{i}, \mathrm{j})\right\}
$$

These sets partition $G(k, n) \times F l_{n} \times F l_{n}$. We define the sets $Y_{0} \bullet$ to be the projections of $X_{\circ} \bullet$ onto $G(k, n)$ obtained by fixing an $M$. and an $F$.

The closures $\overline{Y_{0}}$ are called closed two-step Schubert varieties, Richardson varieties, or skew Schubert varieties.
(1) $G(k, n)=\bigsqcup Y_{\circ}$. for fixed $M$. and $F$.
(2) $Y_{0} \bullet \not \emptyset \emptyset$ only if every white checker is happy: there is one black checker due north and one due west (inclusive) of every white checker, and every row and column has at most one white checker in it.
(3) $X_{A}\left(M_{.}\right) \cap X_{B}(F)=.\overline{Y_{0}}$ where $\bullet$ gives the position of $M$. with respect to $F . \circ=\left\{\left(a_{1}, b_{k}\right), \ldots,\left(a_{k}, b_{1}\right)\right\}$, $A=\left\{a_{1}<\ldots<a_{k}\right\}, B=\left\{b_{1}<\ldots<b_{k}\right\}$.
Lemma 10.1. The "variety" $Y_{\circ} \bullet$ is irreducible and smooth, and $\operatorname{dim} Y_{\bullet} \bullet \sum \mathrm{b}(\mathrm{i}, \mathrm{j})-\mathrm{w}(\mathrm{i}, \mathrm{j})$ where the sum is over all $(i, j)$ containing white checkers.

## Example 10.2.

$$
\left[\begin{array}{lllll} 
& & & & \bullet e_{5} \\
& & \bullet e_{3} & & \\
& & & \bullet e_{4} & \\
\bullet e_{1} & & \circ & & \\
& \bullet e_{2} & & & \circ
\end{array}\right]
$$

This corresponds to $F .=<e_{1}, \ldots, e_{5}>, M .=<e_{2}, e_{1} e_{4} e_{3} e_{5}>$ and a $V$ with basis $<v_{1}, v_{2}>$. $v_{1} \in \operatorname{span}\left\{\mathrm{e}_{1}, \mathrm{e}_{3}\right\}$ but is neither of these two, so we can choose $v_{1}=\star e_{1}+e_{3}$ for $\star \neq 0$ for $2-1=1$ dimension of freedom. Likewise $v_{2} \in \operatorname{span}\left\{\mathrm{e}_{1}, \ldots, \mathrm{e}_{5}\right\}$ and has nonzero projection onto $e_{5}$, so we may choose $v_{2}=\star e_{1}+\star e_{2}+\star e_{3}+\star e_{4}+e_{5}$ and can eliminate $e_{3}$ altogether using $v_{1}$.

Corollary 10.3. $Y_{\circ \bullet_{\text {init }}}=X_{A}\left(M_{.}\right) \cap X_{B}(F.) \neq \emptyset$ if and only if
(1) all white checkers are happy,
(2) the union of the partitions $A$ and $B$ covers all the squares, and
(3) $\operatorname{dim}_{{ }_{0}}{ }_{\text {init }}=\mathrm{k}(\mathrm{n}-\mathrm{k})-\operatorname{codim} \mathrm{X}_{\mathrm{A}}-\operatorname{codim} \mathrm{X}_{\mathrm{B}}$.

If these conditions hold, $\left[X_{A}\right]\left[X_{B}\right] \neq 0$ so the degenerations will lead somewhere.

### 10.6. Degeneration Rules.

(1) Black checkers start in the configuration $\bullet_{\text {init }}$ along the antidiagonal. White checkers start in the configuration $\circ_{A B}=\left\{\left(a_{1}, b_{k}\right), \ldots,\left(a_{k}, b_{1}\right)\right\}$.
(2) Black moves: sort black checkers (corresponding to basis vectors) by transposing two adjacent rows at a time, working from left to right.
(3) White moves: after each black move, white checkers move according to nine rules on $2 \times 2,3 \times 3$, and $4 \times 4$ critical submatrices, obtained by considering which white checkers can be affected by the next black move.
(4) Black moves correspond to degenerations: $\circ \mapsto \circ_{\text {next }}$ where $F$. $=<e_{1}, \ldots, e_{n}>, M_{\text {. }}=<$ $b_{1}, \ldots, b_{n}>$ corresponds to taking the limit as $t \rightarrow 0$ of $b_{i}=t e_{j}+(1-t) e_{k}$.
(5) Two key rules:
(a) If $Y_{\circ} \mapsto Y_{\circ^{o^{\prime}} \bullet_{\text {next }}}$ then $Y_{\circ^{\prime}} \bullet_{\text {next }}$ has the same dimension.
(b) All span and intersection data are preserved.

## 10.7. $2 \times 2$ Rules.




## 10.8. $\mathbf{3 \times 3} \mathbf{R}$ Rules.



## 11. Izzet Coskun's Lecture (April 27, 2007)

Let's set up the basic definitions and notation. Let $G(k, n)$ be the Grassmannian manifold of $k$-planes in $n$ space. Let

$$
\Sigma_{\lambda}=\left\{[\Lambda] \in G(k, n) \mid \operatorname{dim}\left(\Lambda \cap F_{n-k+i-\lambda_{i}} \geq i\right\} .\right.
$$

be the Schubert variety in $G(k, n)$ indexed by the partition $\lambda$ with $n-k \geq \lambda_{1} \geq \lambda_{2} \geq \cdots \geq \lambda_{k} \geq 0$. (Note, this is written in the reverse order to the lectures above.) Let $\sigma_{\lambda}$ be the corresponding Schubert cycle in the cohomology ring for $G(k, n)$. The notation follows Lecture 6 of Harris's "Algebraic Geometry" which is a good book for more details .
Example 11.1. Consider $\sigma_{(1)} \sigma_{(1)}$ in $G(2,4)$. Note $G(2,4)=\mathbb{G}(1,3)$ when we projectivize. The class $\sigma_{(1)} \cdot \sigma_{(1)}$ corresponds with the family of lines meeting two given lines in general position in 3 -space. To find the irreducible components of this family, swing one of the lines around until it crosses the other line in some plane. Now, which lines intersect both lines. There are two components:
(1) The family of lines passing through the plane exactly at the point of intersection between the two lines which corresponds with the Schubert cycle $\sigma_{(2)}$.
(2) The family of lines in the plane defined by the two lines which corresponds with the Schubert cycle $\sigma_{(1,1)}$.
Therefore, $\sigma_{(1)} \cdot \sigma_{(1)}=\sigma_{(2)}+\sigma_{(1,1)}$.
The coefficients in the expansion of $\sigma_{(1)} \cdot \sigma_{(1)}$ are 1 because of the Pieri rule of course, but we could have figured this out another way by studying tangent spaces. We could have shown that even after swinging $\ell_{2}$ around to meet $\ell_{1}$ we still have a transverse intersection.

Let $V$ be an $n$-dimension vector space and assume $G(k, n)$ is the set of $k$ planes in $V$. If $S \in G(k, n)$ the we get a short exact sequence

$$
0 \longrightarrow S \longrightarrow V \longrightarrow V / S \longrightarrow 0 .
$$

Fact 11.2. The tangent space to the Grassmannian at a point $[\Lambda]$ is given by

$$
T_{[\Lambda]}(G(k, n))=\operatorname{Hom}(\Lambda, V / \Lambda) \approx\{k \times(n-k) \text { matrices }\} .
$$

Therefore, $\operatorname{dim} T_{[\Lambda]}(G(k, n))=k \cdot(n-k)$ as expected for a smooth manifold.
Similarly, assume $[\Lambda] \in \Sigma_{\lambda}$ is a smooth point. Suppose further that $\operatorname{dim}\left(\Lambda \cap F_{n-k+i-\lambda_{i}}\right)=i$ which holds on a dense open set in $\Sigma_{\lambda}$ for all $i$. Then

$$
T_{[\Lambda]}\left(\Sigma_{\lambda}\right)=\left\{\phi \in \operatorname{Hom}(\Lambda, V / \Lambda) \mid \phi\left(\Lambda \cap F_{n-k+i-\lambda_{i}}\right) \in F_{n-k+i-\lambda_{i}} / \Lambda \quad \forall i\right\}
$$

Key Fact: To show that at $[\Lambda] \in \Sigma_{(1)}\left(\ell_{1}\right) \cap \Sigma_{(1)}\left(\ell_{2}\right)$ the two Schubert varieties meet transversally, we need to show that the corresponding tangent spaces intersect transversally, i.e.

$$
\operatorname{dim} T_{[\Lambda]}\left(\Sigma _ { ( 1 ) } ( \ell _ { 1 } ) \cap T _ { [ \Lambda ] } \left(\Sigma_{(1)}\left(\ell_{2}\right)=2\right.\right.
$$

(hmm, I am a little lost here. I think this intersection should have dimension 0 , but the example given has dimension 1 and it says above we should have dimension 2.)
11.1. Coskun's version of the Littlewood-Richardson Rule for $G(k, n)$. Fix an ordered basis say $\left(e_{1}, e_{2}, \ldots, e_{n}\right)$ for $V$. Fix $k$ vector spaces $V_{1}, \ldots, V_{k}$ which are spans of consecutive subsets of basis elements. To visualize these sets, put $e_{1}, e_{2}, \ldots, e_{n}$ along the antidiagonal of a matrix with $e_{1}$ in the lower left corner. Then if $V_{i}$ is spanned by $\left\{e_{j}, e_{j+1} \cdots, e_{k}\right\}$, we represent $V_{i}$ by the minimal square containing all the basis elements in $V_{i}$. The collection of squares is called a Mondrian tableau or tableau for short. These tableau were named after the artist Piet Mondrian 1872-1944. Check out
http://en.wikipedia.org/wiki/Piet_Mondrian
for more information.
Example 11.3. If $V_{1} \subset V_{2} \subset \cdots \subset V_{k}$ then

$$
S=\left\{[\Lambda] \in G(k, n) \mid \operatorname{dim}\left(\Lambda \cap V_{i}\right) \geq i\right\}
$$

is a Schubert variety. The set $S$ will be represented by a nested sequence of squares in the Mondrian tableau. For a more specific example, if $V_{j}$ is spanned by $\left\{e_{1}, \ldots, e_{i_{j}}\right\}$ then this set is the Schubert variety $X_{\left\{i_{1}, \ldots, i_{k}\right\}}$ in our previous class notation. This Schubert variety will correspond with a nested sequence of squares of side lengths $i_{1}, \ldots, i_{k}$ with a common lower left corner in the lower left corner of the matrix.

Example 11.4. If the $V_{i}$ 's have no inclusion relations and the basis elements in $V_{i}$ precede the basis elements for $V_{i+1}$ then the corresponding Mondrian tableau will have $k$ non-overlapping squares. This tableau will correspond with the intersection of two Schubert varieties in opposite position.

We want to generalize the idea in Example 11.1 to intersecting any two Schubert varieties in the Grassmannian. More specifically, take two Schubert varieties $\Sigma_{\lambda}$ and $\Sigma_{\mu}$ with respect to ordered bases $\left(e_{1}, \ldots, e_{n}\right)$ and $\left(e_{n}, \ldots, e_{1}\right)$ respectively. We want to find a sequence of degenerations which will keep the cohomology class of the intersection of these varieties the same but eventually leads to the union of irreducible components. The degenerations we will use are all of the form $t e_{i}+(1-t) e_{j}$. At $t=1$ is gives the basis element $e_{i}$ and then as $t$ approaches 0 this degenerates $e_{i}$ into $e_{j}$. If $U$ is the span of $\left(e_{i}, \ldots, e_{j-1}\right)$ then after the degeneration $U$ will be the span of $\left(e_{i+1}, \ldots, e_{j}\right)$, therefore the degeneration is visualized as moving the square for $U$ north east along the diagonal by one unit.

Definition 11.5. A tableau for $G(k, n)$ is a collection of $k$ distinct squares such that
(1) Each square corresponds to a subspace spanned by consecutive basis elements.
(2) No 2 squares share the same lower left corner.
(3) If $S_{1}$ and $S_{2}$ are any 2 squares of the tableau that share their upper right corner then every square whose lower left corner is southwest of the lower left corner of $S_{2}$ contains $S_{2}$.

Definition 11.6. A square $S$ is nested if all the squares containing $S$ are totally ordered by inclusion and for any other square $S^{\prime}$ either $S \subset S^{\prime}$ or $S^{\prime} \subset S$.

Definition 11.7. Let $S$ be the square containing $<e_{i}, \ldots, e_{j}>$ and $S^{\prime}$ be the square containing $<$ $e_{i+1}, \ldots, e_{j+1}>$ so $S^{\prime}$ is the result of degenerating $e_{i}$ to $e_{j+1}$. A neighbor $N$ of $S$ is a square such that
(1) $e_{j+1} \in N$.
(2) $N$ does not contain $S$.
(3) Let $\tilde{S}$ be any square whose lower left corner is between $S^{\prime}$ and $N$, then either $\tilde{S} \subset S$ or $N \subset \tilde{S}$.

Note that by definition the neighbors of $S$ are all ordered by containment.
Degeneration Algorithm: Given a tableau $M$,
Step 1. If every square of $M$ is nested, stop.
Step 2. If not, let $S$ be the square whose lower left corner is southwest most among the non-nested squares. Let $N_{1} \subset N_{2} \subset N_{r}$ be the neighbors of $S$.
Step 3. For every neighbor define $M_{1}\left(N_{i}\right)$ to be the tableau obtained from $M$ by replacing $S$ and $N_{i}$ by the square representing $S^{\prime} \cap N$ and the square representing the span of $S \cup N_{i}$. Define $M_{0}$ to be the tableau obtained from $M$ by replacing $S$ by $S^{\prime}$ and normalizing again.
Step 4. Among the collection $M_{0}, M_{1}\left(N_{1}\right), \ldots, M_{1}\left(N_{r}\right)$ retain the tableau with the same dimension as $M$ and repeat Step 1 with each of these as the given tableau.
For more information and pictures see Izzet Coskun's web page. A new preprint will be coming soon. For now see

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http://math.mit.edu/~ coskun/seattleoct17.pdf
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12. Divided difference operators and the Chow cohomology of the flag variety (MAY 2, 2007)

In previous lectures we derived a presentation for the Chow cohomology of the Grassmannian, given by:

$$
A^{*}(G(k, n)) \simeq \mathbb{Z}\left[s_{\lambda}\left(x_{1}, \ldots, x_{n}\right)\right] /<s_{\lambda}: \lambda \subsetneq k \times(n-k)>
$$

Here the isomorphism is given by $\left[X_{\lambda}\right] \mapsto s_{\text {dual }(\lambda)}$. Today we will discuss a similar result, due to Borel, for the Chow cohomology of the Flag variety:

Theorem 12.1 (Borel). The Chow cohomology of the flag variety $\left.F l_{n}\right)$ is given by $A^{*}\left(F l_{n}\right) \simeq \mathbb{Z}\left[x_{1}, \ldots, x_{n}\right] / I_{n}$, where $I_{n}$ is the ideal

$$
I_{n}:=<e_{1}\left(x_{1}, \ldots, x_{n}\right), \ldots, e_{n}\left(x_{1}, \ldots, x_{n}\right)>.
$$

Here $e_{i}\left(x_{1}, \ldots, x_{n}\right)$ is the $i^{\text {th }}$ elementary symmetric function; for instance $e_{1}\left(x_{1}, \ldots, x_{n}\right)=x_{1}+$ $\cdots+x_{n}$, and $e_{n}=x_{1} x_{2} \cdots x_{n}$.

We will be interested in determining ? ${ }_{w}$, the image of $\left[X_{w}\right]$ under the above isomorphism $\left[X_{w}\right] \mapsto ?_{w}$. We first point out some properties that ? ${ }_{w}$ that must satisfy.

- The isomorphism is a ring map: If $\left[X_{u}\right]\left[X_{v}\right]=\sum c_{u v}^{w}\left[X_{w}\right]$ then we must have $?_{u} ?_{v}=\sum c_{u v}^{w} ?_{w}$.

The map is also graded: Since $\operatorname{deg}\left(?_{u} ?_{v}\right)=\operatorname{deg}\left(?_{u}\right)+\operatorname{deg}\left(?_{v}\right)$, and $\operatorname{codim}\left(X_{u} \cap X_{v}\right)=$ $\operatorname{codim}\left(X_{u}\right)+\operatorname{codim}\left(X_{v}\right)$, we must have $\operatorname{deg}\left(?_{w}\right)=\operatorname{codim}\left(X_{w}\right)$.

- We have seen $\left[X_{w_{0}}\right]\left[X_{u}\right]=\left[X_{w_{0}}\left(F_{\bullet}\right) \cap X_{u}\left(G_{\bullet}\right)\right]=\left[X_{u}\right]$, and hence $\left[X_{w_{0}}\right] \mapsto 1$.
- If $\ell(u)+\ell(v)=\binom{n}{2}=\operatorname{dim}\left(F l_{n}\right)$, then we have the formula

$$
\left[X_{u}\right]\left[X_{v}\right]=\left\{\begin{array}{ll}
{\left[X_{i d}\right]} & \text { if } v=u w_{0} \\
0 & \text { otherwise }
\end{array}\right\}
$$

Hence in this case we must have

$$
?_{u} ?_{v}=\left\{\begin{array}{ll}
?_{w_{0}} & \text { if } v=u w_{0} \\
0 & \bmod I_{n} \\
\text { otherwise }
\end{array}\right\}
$$

Exercise 12.2. Show that $R_{n}:=\mathbb{Z}\left[x_{1}, \ldots, x_{n}\right] / I_{n}$ is isomorphic as a vector space to the (linear) span of the set $\left\{x_{1}^{i_{1}} x_{2}^{i_{2}} \cdots x_{n}^{i_{n}}: i_{k} \leq n-k\right\}$. For example, if $n=3$ then $R_{3}=\operatorname{span}\left\{1, x_{1}, x_{2}, x_{1}^{2}, x_{1} x_{2}, x_{1}^{2} x_{2}\right\}$. Garcia calls this the 'Artin basis'.

Hint: Find a Gröbner basis for $I_{n}$.
The Artin basis will help us determine the Hilbert series for the ring $R_{n}$. We will need the following bijection between elements of this basis and elements of the symmetric group.
Lemma 12.3. There exists a bijection $S_{n} \sim\left\{\left(i_{1}, \ldots, i_{n}\right): i_{k} \leq n-k\right\}$, sending $w$ to code( $\left.w\right)$.
Here the code of an element $w \in S_{n}$ is defined to be the vector:

$$
\operatorname{code}(w):=(\mid * s \text { on row } 1|,| * s \text { on row } 2|, \ldots,| * s \text { on row } n \mid)
$$

in the matrix obtained by crossing out entries below and to the right of each entry of $w$. For example, code $(316425)=(2,0,3,1,0,0)$.

Using this we derive the Hilbert series for $R_{n}$ :

$$
\sum \operatorname{dim} R_{n}^{(i)} t^{i}=\sum_{w \in S_{n}} t^{\ell(n)}=\prod_{i=1}^{n-1}\left(1+t+\cdots+t^{i}\right)
$$

We now return to the issue of determining ? ${ }_{w}$, the image of $\left[X_{w}\right]$ in the isomorphism of Theorem 12.1. Using the properties that we so far have derived, we consider the case $n=3$ and make a naive guess for what these values might be.

Recall that the set of generators for $A^{*}\left(F l_{3}\right)$ contains a single element $\left[X_{i d}\right]$ of rank 0 , two elements $\left[X_{213}\right],\left[X_{132}\right]$ of rank 1, two elements $\left[X_{312}\right],\left[X_{231}\right]$ of rank 2, and a single element $\left[X_{321}\right]$ of rank $\binom{3}{2}=3$. To define a map to $R_{n}$ we choose to map these generators to elements of the Artin basis according to $\left[X_{i d}\right] \mapsto x_{1}^{2} x_{2},\left[X_{213}\right] \mapsto x_{1}^{2},\left[X_{132}\right] \mapsto x_{1} x_{2},\left[X_{312}\right] \mapsto x_{1},\left[X_{231}\right] \mapsto x_{2}$, and $\left[X_{321}\right] \mapsto 1$.

One can check that some of the conditions on our list of requirements for $?_{w}$ are indeed satisfied (e.g., it is a graded group map that sends $\left[X_{w_{0}}\right]$ to 1 ).
Exercise 12.4. Show that $x_{i}^{3}=0 \bmod I_{3}$ and $x_{1} x_{2}^{2}=o \bmod I_{3}$. For instance $x_{1}^{3}=x_{1}^{2}\left(x_{1}+x_{2}+\right.$ $\left.x_{3}\right)-x_{1}\left(x_{1} x_{2}+x_{1} x_{3}+x_{2} x_{3}\right)+x_{1} x_{2} x_{3}$.

However, in checking that our map is indeed a ring morphism, we note that $x_{1} \cdot x_{2}=x_{1} x_{2}$, and hence are asked to verify the formula:

$$
\left[X_{312}\right]\left[X_{231}\right]=\left[X_{132}\right] .
$$

However, we consider the intersection of flags in question and compute that in fact $\left[X_{312}\right]\left[X_{231}\right]=$ $\left[X_{132}\right]+\left[X_{312}\right]$. Accordingly, we tweak our map to $R_{n}$ described above to now have $\left[X_{231}\right] \mapsto x_{1}+x_{2}$. One can check that this indeed defines the desired isomorphism.

We seek a method for determining $?_{w}$ in a uniform way for all $n$. This is provided by the following procedure due to Bernstein, Gelfand, and Gelfand (BGG).

Observation. We note that $\operatorname{dim}\left(R_{n}^{\binom{n}{2}}\right)=1$, and so for ? $w_{0}$ we can choose and homogeneous polynomial $\sigma_{w_{0}}$ of degree $\binom{n}{2}$ that is not in $I_{n}$. We follow BGG and choose $\sigma_{w_{0}}=\prod_{1 \leq i<j \leq n}\left(x_{i}-x_{j}\right)$ to be the so-called 'Vandermonde determinant'.

To determine the other ? ${ }_{w}$, we will need the notion of a divided difference operator on $\mathbb{Z}\left(x_{1}, \ldots, x_{n}\right)$. If $f \in \mathbb{Z}\left(x_{1}, \ldots, x_{n}\right)$, define $s_{i} f\left(x_{1}, \ldots, x_{n}\right):=f\left(x_{1}, \ldots, x_{i+1}, x_{i}, \ldots, x_{n}\right)$. The divided difference operator $\partial_{i}$ is then defined to be

$$
\partial_{i}(f):=\left(f-s_{i} f\right) /\left(x_{i}-x_{i+1}\right) .
$$

Observations. We point out some properties of the $\partial_{i}$.

- $\partial_{i}(f)=0$ if $f=s_{i} f$.
- $\partial_{i}\left(s_{i} f\right)=-\partial_{i} f$
- $\partial_{i}\left(x_{i}^{r} x_{i+1}^{s}\right)=x_{i}^{r} x_{i+1}^{r}\left(x_{i+1}^{s-r-1}+x_{i} x_{i+1}^{s-r-2}+\cdots+x_{i}^{s-r-1}\right.$ if $s>r$.
- $\partial_{i}\left(x_{i}^{r} x_{i+1}^{s} \cdot f\left(x_{1}, \ldots, \hat{x_{i}}, x_{i+1}, \ldots, x_{n}\right)\right)=f \cdot \partial_{i}\left(x_{i}^{r} x_{i+1}^{s}\right)$.

These last two facts imply that if $f$ is polynomial then so is $\partial_{i}(f)$. The $\partial_{i}$ also satisfy some Coxeter like relations.

- $\partial_{i} \partial_{j}=\partial_{j} \partial_{i}$ if $|i-j|>1$.
- $\partial_{i} \partial_{i+1} \partial_{i}=\partial_{i+1} \partial_{i} \partial_{i+1}$.
- $\partial_{i}^{2}=0$.

Exercise 12.5 (Leibnitz rule for divided difference operators). Show that $\partial_{i}(f g)=\left(\partial_{i} f\right) g+\left(s_{i} f\right)\left(\delta_{i} g\right)$.
Theorem 12.6 (Bernstein, Gelfand, Gelfand). Suppose $\left[X_{w w_{0}}\right] \mapsto \sigma_{w} \in R_{n}$ under the isomorphism of 12.1. Then we have the following:

$$
\partial_{i}\left(\sigma_{w}\right)= \begin{cases}\sigma_{w s_{i}} \bmod I_{n} & \text { if } w_{i}>w_{i+1} \\ 0 \bmod I_{n} & \text { otherwise }\end{cases}
$$

The proof of this theorem relies on the following key lemma.
Lemma 12.7. If $s_{a_{1}} \cdots s_{a_{p}}=s_{b_{1}} \cdots s_{b_{p}}=w$ are both reduced expressions for $w \in S_{n}$ (so that $p=$ $\ell(w))$, then $\partial_{a_{1}} \cdots \partial_{a_{p}}=\partial_{b_{1}} \cdots \partial_{b_{p}}$.

To prove the lemma we introduce the notion of a string diagram associated to an element $w \in S_{n}$.
Definition 12.8. If $w \in S_{n}$, we define the following sets of inversions:

$$
\begin{gathered}
I(w):=\{(i, j): i<j, w(i)>w(j)\} . \\
J(w):=I\left(w^{-1}\right)=\{(i, j) \text { :string } i \text { crosses string } j\} .
\end{gathered}
$$

Every element $w \in S_{n}$ can be written $s_{a_{1}} \cdots s_{a_{p}}$ such that $p=|I(w)|=|J(w)|=\ell(w)$. In this case, we will say that $\left(a_{1} \cdots a_{p}\right)$ is a reduced word for $w$. The collection of all reduced words will be denoted $R(w):=\left\{\left(a_{1} \cdots a_{p}\right): s_{a_{1}} \cdots s_{a_{p}}, p=\ell(w)\right\}$. For example, $R(4132)=\{3213,3231,2321\}$.
Lemma 12.9. Suppose $a_{1} \cdots a_{p}$ is a reduced word for $w$. Then we have $J(w)=\left\{\left(a_{1}, a_{1}+1\right), s_{a_{1}}\left(a_{2}, a_{2}+\right.\right.$ 1), $\left.\ldots, s_{a_{1}} \cdots s_{a_{p-1}}\left(a_{p}, a_{p}+1\right)\right\}$.

Proof. We prove this by induction on the length $\ell(w)$. We see $J(i d)=\emptyset$ and also $J\left(s_{i}\right)=\{(i, i+$ $1)\}$ so that result holds for $\ell(w)=0,1$. We assume the result holds for all $v$ with $\ell(v)<\ell(w)$. If $v=s_{a_{1}} \cdots s_{a_{p-1}}$ then all string crossings for $w$ are string crossings for $v$ except the last crossing $\left(v\left(a_{p}\right), v\left(a_{p}+1\right)\right)$. This proves the lemma and also gives us the following formula:

$$
J\left(w s_{r}\right)= \begin{cases}J(w) \cup w(r, r+1) & \text { if } \ell\left(w s_{r}\right)>\ell(w) \\ J(w) \backslash w(r, r+1) & \text { if } \ell\left(w s_{r}\right)<\ell(w)\end{cases}
$$

This formula provides us with an interpretation of the so-called weak Bruhat order on $S_{n}$. By definition these are given by:

$$
\begin{aligned}
& \text { Right order: } v \leq w \Leftrightarrow J(v) \subseteq J(w), v s_{i}>w \Leftrightarrow \ell\left(v s_{i}\right)>\ell(w) \\
& \text { Left order: } v \geq w \Leftrightarrow I(v) \subseteq I(w), s_{i} v>v \Leftrightarrow \ell\left(s_{i} v\right)>\ell(w) .
\end{aligned}
$$

Lemma 12.10 (Exchange Lemma). Suppose $\left(a_{1} \cdots a_{p}\right),\left(b_{1} \cdots b_{p}\right) \in R(w)$ are both reduced expressions for $w$. Then there exists an integer $i$ such that $\left(b_{1} a_{1} \cdots \hat{a}_{i} \cdots a_{p}\right) \in R(w)$.
Proof. From Lemma 12.9, we have $\left(b_{i}, b_{i}+1\right) \in J(w)$ and hence $\left(b_{i}, b_{i}+1\right)=s_{a_{1}} \cdots s_{a_{i-1}}\left(a_{i}, a_{i}+1\right)$. This implies that $s_{b_{1}}=s_{a_{1}} \cdots s_{a_{i-1}} s_{a_{i}} s_{a_{i-1}} \cdots s_{a_{1}}$ and hence $s_{b_{1}} s_{a_{1}} \cdots s_{a_{i-1}}=s_{a_{1}} \cdots s_{a_{i}}$. We then have $w=s_{a_{1}} \cdots s_{a_{p}}=s_{b_{1}} s_{a_{1}} \cdots s_{a_{i-1}} s_{a_{i+1}} \cdots s_{a_{p}}$, as desired.
Definition 12.11. Given an element $w \in S(n)$, we define $G(w)$ to be the graph with vertex set $R(w)$, and with edges given by $\left(a_{1} \cdots a_{p}\right) \sim\left(b_{1} \cdots b_{p}\right) \Leftrightarrow$ the expressions differ by some $s_{i}$.
Theorem 12.12 (Tits). The graph $G(w)$ is connected.
Proof. We prove this by induction on $\ell(w)$, the length of $w$. First, it's clear that $G\left(s_{i}\right)$ is connected. Next assume $a=\left(a_{1} \cdots a_{p}\right)$ and $b=\left(b_{1} \cdots b_{p}\right)$ are elements of $R(w)$. From the previous lemma we have that $\left(b_{1} a_{1} \cdots \hat{a_{i}} \cdots a_{p}\right) \in R(w)$, and hence by induction $\left(b_{1} \cdots b_{p}\right)$ is connected to $\left(b_{1} a_{1} \cdots \hat{a_{i}} \cdots a_{p}\right)$. If $i \neq$ $p$ then we also have $\left(b_{1} a_{1} \cdots \hat{a_{i}} \cdots a_{p}\right)$ adjacent to $\left(a_{1} \cdots a_{p}\right)$. If $i=p$ then we have $\left(b_{1} a_{1} \cdots a_{p-1}\right) \in$ $R(w)$, and by the Exchange Lemma we have ( $\left.a_{1} b_{1} a_{1} \cdots \hat{a_{j}} \cdots a_{p-1}\right) \in R(w)$, which is in turn adjacent to $\left(b_{1} a_{1} b_{1} \cdots \hat{a_{j}} \cdots a_{p-1}\right) \in R(w)$. By induction, the former of these is adjacent to $a$, while the latter is adjacent to $b$, and hence the claim follows.

## 13. Schubert Polynomials (May 4, 2007)

Last time we saw that there is a map

$$
A^{*} F \ell_{n} \rightarrow \mathbb{Z}\left[x_{1}, \ldots, x_{n}\right] /<e_{1}, \ldots, e_{n}>
$$

and the Bernstein, Gelfand, Gelfand Theorem says that the classes corresponding Schubert varieties form an additive basis for the Chow ring. This map sends $\left[X_{w}\right] \mapsto \sigma_{w w_{0}}$ so that modulo the ideal $\left.<e_{1}, \ldots, e_{n}\right\rangle$, the divided difference operators act by

$$
\partial_{i}\left(\sigma_{w}\right)= \begin{cases}\sigma_{w s_{i}} & \text { if } w_{i}>w_{i+1} \\ 0 & \text { otherwise }\end{cases}
$$

Last time, we had the following candidates for the $\sigma_{w}$ :

$$
\begin{aligned}
& {\left[X_{321}\right] \mapsto \sigma_{123}=1} \\
& {\left[X_{312}\right] \mapsto \sigma_{213}=x_{1}} \\
& {\left[X_{231}\right] \mapsto \sigma_{132}=x_{1}+x_{2}} \\
& {\left[X_{213}\right] \mapsto \sigma_{312}=x_{1}^{2}} \\
& {\left[X_{132}\right] \mapsto \sigma_{231}=x_{1} x_{2}} \\
& {\left[X_{123}\right] \mapsto \sigma_{321}=x_{1}^{2} x_{2} .}
\end{aligned}
$$

Let's check that this actually works. Start at the bottom of this list with $\left[X_{123}\right] \mapsto \sigma_{321}$. Here, $3>2$ and $2>1$, so we expect that $\partial_{1} \sigma_{321}$ and $\partial_{2} \sigma_{321}$ are both nonzero. In particular,

$$
[321] \cdot s_{1}=[321][213]=[231]
$$

and

$$
[321] \cdot s_{2}=[321][132]=[312]
$$

so we expect that $\partial_{1} \sigma_{321}=\sigma_{231}$ and $\partial_{2} \sigma_{321}=\sigma_{312}$. Indeed,

$$
\partial_{1}\left(x_{1}^{2} x_{2}\right)=\frac{x_{1}^{2} x_{2}-x_{1} x_{2}^{2}}{x_{1}-x_{2}}=x_{1} x_{2}=\sigma_{231}
$$

and

$$
\partial_{2}\left(x_{1}^{2} x_{2}\right)=\frac{x_{1}^{2} x_{2}-x_{1}^{2} x_{3}}{x_{2}-x_{3}}=x_{1}^{2}=\sigma_{312}
$$

Next, we look at $\left[X_{132}\right] \mapsto \sigma_{231}$. Since $1<3$, we expect that $\partial_{1} \sigma_{231}=0$; and since $3>1$, we should get that $\partial_{2} \sigma_{213}=\sigma_{[213] \cdot s_{2}}=\sigma_{231}=x_{1} x_{2}$. Indeed, $\partial_{1}\left(x_{1} x_{2}\right)=0$ because $x_{1} x_{2}$ is symmetric in $x_{1}$ and $x_{2}$, and $\partial_{2}\left(x_{1} x_{2}\right)=\frac{x_{1} x_{2}-x_{1} x_{3}}{x_{2}-x_{3}}=x_{1}$. The remaining cases follow similarly: $\partial_{1} \sigma_{312}=$ $x_{1}+x_{2}=\sigma_{[312] \cdot s_{1}}$ and $\partial_{2} \sigma_{312}=0 ; \partial_{1} \sigma_{132}=0$ and $\partial_{2} \sigma_{132}=1=\sigma_{[132] \cdot s_{2}} ; \partial_{1} \sigma_{213}=1=\sigma_{[213] \cdot s_{1}}$ and $\partial_{2} \sigma_{213}=0$.

Now it might be natural to ask what would happen if we took a different choice for the polynomial $\sigma_{321}$ ? Suppose for example we took $\sigma_{321}=x_{3}^{2} x_{2}$. Then we compute the following:

$$
\partial_{1} \partial_{2} \partial_{1}\left(x_{3}^{2} x_{2}\right)=\partial_{1} \partial_{2}\left(-x_{3}^{2}\right)=\partial_{1}\left(x_{2}+x_{3}\right)=-1
$$

and

$$
\partial_{2} \partial_{1} \partial_{2}\left(x_{3}^{2} x_{2}\right)=\partial_{2} \partial_{1}\left(-x_{2} x_{3}\right)=\partial_{2}\left(x_{3}\right)=-1 .
$$

This doesn't quite work because we want $\left[X_{w_{0}}\right]\left[X_{u}\right]=\left[X_{u}\right]$ for all $u \in S_{n}$, and in this case we get $\left[X_{321}\right] \mapsto \sigma_{123}=\partial_{1} \partial_{2} \partial_{1} \sigma_{321}=-1$. However, if we take $\sigma_{321}=-x_{3}^{2} x_{2}$, everything works. However, modulo the ideal generated by the elementary symmetric functions, $-x_{3}^{2} x_{2}$ and $x_{1}^{2} x_{2}$ are equal:

$$
\begin{aligned}
x_{1}^{2} x_{2}+x_{2} x_{3}^{2}= & x_{1} x_{2}\left(x_{1}+x_{2}+x_{3}\right)-x_{1} x_{2}^{2}-x_{1} x_{2} x_{3} \\
& +x_{2} x_{3}\left(x_{1}+x_{2}+x_{3}\right)-x_{1} x_{2} x_{3}-x_{2}^{2} x_{3} \\
\equiv & x_{1} x_{2}^{2}+x_{2}^{2} x_{3} \\
= & x_{2}^{2}\left(x_{1}+x_{2}+x_{3}\right)-x_{2}^{3} \\
\equiv & x_{2}^{3} .
\end{aligned}
$$

In the last lecture, we saw that $x_{2}^{3} \equiv 0$ modulo the ideal.
However, our original choice, with $\sigma_{321}=x_{1}^{2} x_{2}$ had the advantage that each monomial in $\sigma_{w}$ had nonnegative integer coefficients. In fact, we could have taken $\sigma_{w_{0}}$ to be the Vandermonde determinant $\prod_{i<j}\left(x_{i}-x_{j}\right)$, but the resulting system is a mess. The niceness of our original choice motivates the following definition from Lascoux and Schutzenberger.
Definition 13.1. Let $\mathfrak{S}_{w_{0}}=x_{1}^{n-1} x_{2}^{n-2} \cdots x_{n-1}^{1} x_{n}^{0}$. For $w<w_{0}$, we define the Schubert polymomial $\mathfrak{S}_{w}$ as

$$
\mathfrak{S}_{w}=\partial_{w^{-1} w_{0}} \cdot \mathfrak{S}_{w_{0}}
$$

We claim that this is a "good" definition of the Schubert polynomials in the following proposition. First we need some notation: let $\delta=(n-1, n-2, \ldots, 1,0) \in[n]^{n}$. We can partially order elements $\alpha=\left(\alpha_{1}, \alpha_{2}, \ldots, \alpha_{n}\right) \in[n]^{n}$ by containment, decreeing that $\alpha \subset \beta$ if $\alpha_{i} \leq \beta_{i}$ for all $i$. For such an $\alpha$, let $x^{\alpha}=x_{1}^{\alpha_{1}} x_{2}^{\alpha_{2}} \cdots x_{n}^{\alpha_{n}}$.
Proposition 13.2. The Schubert polynomials $\mathfrak{S}_{w}(x)$ have the following properties:
(1) $\mathfrak{S}_{w_{0}}=x^{\delta}$ and $\mathfrak{S}_{i d}=1$. Moreover, $\mathfrak{S}_{w}$ is a homogeneous polynomial of degree $\ell(w)$ in $\mathbb{Z}\left[x_{1}, \ldots, x_{n}\right]$.
(2) The monomial terms of $\mathfrak{S}_{w}$ are contained in the set of Artin monomials, i.e. $\mathfrak{S}_{w}$ is contained in the $\mathbb{Z}$-span of monomial terms of the form $x^{\alpha}$ for $\alpha \subset \delta$. This implies that the $\mathfrak{S}_{w}$ are nonzero.
(3) If $r$ is the last descent of $w \in S_{n}$, i.e. $w_{r}>w_{r+1}<w_{r+2}<\ldots<w_{n}$, then $\mathfrak{S}_{w}$ is a polynomial in $x_{1}, \ldots, x_{r}$.
(4) The polynomial $\mathfrak{S}_{w}$ is symmetric in $x^{i}$ and $x^{i+1}$ if and only if $w_{i}<w_{i+1}$.
(5) $\mathfrak{S}_{s_{i}}=x_{1}+\ldots+x_{i}$
(6) (Stability) Let i : $S_{n} \hookrightarrow S_{n+1}$ be the natural inclusion $\left[w_{1}, \ldots w_{n}\right] \mapsto\left[w_{1}, \ldots, w_{n}, n+1\right]$. Then $\mathfrak{S}_{w}=\mathfrak{S}_{i(w)}$.
Proof. (1) By definition, $\mathfrak{S}_{w_{0}}=x^{\delta}$. Let $w_{0}^{(n)}=[n, n-1, \ldots, 1]$ and $w_{0}^{(n-1)}=[n-1, n-$ $2, \ldots, 1, n]$. Certainly $w_{0}^{(n)}=w_{0}^{(n-1)} s_{n-1} s_{n-2} \cdots s_{2} s_{1}$ so

$$
\begin{aligned}
\mathfrak{S}_{i d} & =\partial_{w_{0}^{(n)}} \cdot \mathfrak{S}_{w_{0}} \\
& =\partial_{w_{0}^{(n-1)}} \partial_{n-1} \partial_{n-2} \cdots \partial_{1}\left(x^{\delta}\right) \\
& =\partial_{w_{0}^{(n-1)}} \partial_{n-1} \partial_{n-2} \cdots \partial_{1}\left(x_{1}^{n-1} x_{2}^{n-2} \cdots x_{n-1}^{1}\right) \\
& =\partial_{w_{0}^{(n-1)}} \partial_{n-1} \partial_{n-2} \cdots \partial_{2}\left(x_{1}^{n-2} x_{2}^{n-2} \cdots x_{n-1}^{1}\right) \\
& =\cdots \\
& =\partial_{w_{0}^{(n-1)}}\left(x_{1}^{n-2} x_{2}^{n-3} \cdots x_{n-1}^{0}\right) .
\end{aligned}
$$

But this is the Schubert polynomial indexed by the identity in $S_{n-1}$. Inductively, this is equal to 1.

Finally, to see that the polynomials are homogeneous, simply note that $\mathfrak{S}_{w_{0}}$ is homogeneous of degree $n(n-1) / 2$ and that the divided difference operators preserve homogeneity and decrease degree by 1 .
(2) If $x^{\alpha}$ is a monomial term with $\alpha \subseteq \delta$, then in computing $\partial_{i}\left(x^{\alpha}\right)$, we only care about the terms $x_{i}^{\alpha_{i}}$ and $x_{i+1}^{\alpha_{i+1}}$. Last time, we saw that $\partial_{i}\left(x_{i}^{\alpha_{i}} x_{i+1}^{\alpha_{i+1}}\right)$ is an element of the span of monomial terms of the form $x_{i}^{\beta_{i}} x_{i+1}^{\beta_{i+1}}$ such that $\beta_{i}$ and $\beta_{i+1}$ are at most $\max \left\{\alpha_{i}, \alpha_{i+1}\right\}-1$. In particular, if $x_{i}^{\alpha_{i}} x_{i+1}^{\alpha_{i+1}} \subseteq x^{\delta}$, then so are all terms $x_{i}^{\beta_{i}} x_{i+1}^{\beta_{i+1}}$ with the above restrictions.
(3) This will follow from (4). If $r$ is the last descent of $w$, then $\mathfrak{S}_{w}$ is symmetric in $x_{r+1}, \ldots, x_{n}$. However there isn't an $x_{n}$ term in $\mathfrak{S}_{w_{0}}$, so there can't be an $x_{n}$ term in $\mathfrak{S}_{w}$. Thus none of the terms $x_{r+1}, \ldots, x_{n}$ appear.
(4) By Bernestein-Gelfand-Gelfand, $w_{i}<w_{i+1}$ if and only if $\partial_{i}\left(\mathfrak{S}_{w}\right)=0$ if and only if $\mathfrak{S}_{w}$ is symmetric in $x_{i}$ and $x_{i+1}$.
(5) Since $s_{i}=[1, \ldots, i-1, i+1, i, i+2, \ldots, n]$, by (4) $\mathfrak{S}_{s_{i}}$ is symmetric in $x_{1}, \ldots, x_{i}$; and by (3), $\mathfrak{S}_{s_{i}}$ is a polynomial in only these variables. By (1), $\mathfrak{S}_{s_{i}}$ is homogeneous of degree $\ell\left(s_{i}\right)=1$. Thus $\mathfrak{S}_{s_{i}}=c\left(x_{1}+\ldots+x_{i}\right)$ for some $c \in \mathbb{Z}$. If we apply $\partial_{i}$ to $\mathfrak{S}_{s_{i}}$, on the one hand, we must get $\mathfrak{S}_{s_{i}} \cdot s_{i}=\mathfrak{S}_{i d}=1$. On the other hand, by direct computation, we get

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

Thus $c=1$.
(6) By direct computation, as in part (1).

$$
\begin{aligned}
\mathfrak{S}_{w} & =\partial_{w^{-1} w_{0}}\left(x^{\delta}\right) \\
& =\partial_{w^{-1} w_{0}} \partial_{n} \partial_{n-1} \cdots \partial_{1}\left(x_{1}^{n} x_{2}^{n-1} \cdots x_{n-1}^{2} x_{n}\right) \\
& =\mathfrak{S}_{i(w)}
\end{aligned}
$$

Remark 13.3.

- In geometry, the variables $x_{i}$ are the Chern classes of certain line bundles $E_{i} / E_{i-1}$, but this is special to type A.
- Stability gives us a sequence on cohomology:

$$
\ldots \hookrightarrow H^{*}\left(F \ell_{n}\right) \hookrightarrow H^{*}\left(F \ell_{n+1}\right) \hookrightarrow \ldots
$$

that sends $\left[X_{w w_{0}}\right] \mapsto\left[X_{i\left(w w_{0}\right)}\right]$.
Proposition 13.4. (1) The set $\left\{\mathfrak{S}_{w}: w \in S_{n}\right\}$ forms a basis for $H_{n}=\operatorname{span}\left\{x^{\alpha}: \alpha \subseteq \delta\right\}$.
(2) If we order the Artin monomials in reverse lexicographic order, then each Schubert polynomial $\mathfrak{S}_{w}$ can be written as a sum whose highest monomial term is $x^{c(w)}$.
(3) Let $S_{\infty}=\bigcup_{n=1}^{\infty} S_{n}$. Then $\left\{\mathfrak{S}_{w}: w \in S_{\infty}\right\}$ form a basis for the polynomials in $x_{1}, x_{2}, \ldots$ over $\mathbb{Z}$.

Proof. (1) Since $\left|S_{n}\right|=n!=\operatorname{dim}\left(H_{n}\right)$, we need only show that the Schubert polynomials are independent over $\mathbb{Z}$. Suppose we have some linear combination $\sum_{w \in S_{n}} a_{w} \mathfrak{S}_{w}=0$. In particular, if we look at the homogeneous terms of degree $k$, we must have

$$
\sum_{\substack{w \in S_{n}: \\ l(w)=k}} a_{w} \mathfrak{S}_{w}=0 .
$$

Now for any $u$ of length $k$, if we apply the divided difference operator $\partial_{u}$ to the above sum, we expect to get 0 because each summand $\mathfrak{S}_{w}$ is homogeneous of degree $k$, and applying $\partial_{u}$ lowers the degree by $k$. On the other hand, $\partial_{u}\left(\sum a_{w} \mathfrak{S}_{w}\right)=a_{u}$, hence $a_{u}=0$.
(2) We induct on $\ell(w)$, starting from the top. Certainly $\mathfrak{S}_{w_{0}}=x^{c\left(w_{0}\right)}$. Suppose inductively that $\mathfrak{S}_{w}$ can be written as $x^{\operatorname{code}(w)}$ plus lower terms in reverse lexicographic order, and that $w_{i}>w_{i+1}$. In the diagram $D(w)$, there are 1 's in positions $\left(i, w_{i}\right)$ and $\left(i+1, w_{i+1}\right)$, and the former lies above and to the right of the latter. Thus the number of free entries in the $i$ th row is at least equal to the number of free entries in the $(i+1)$ th row. In terms of the code of the permutation, $c(w)=\left(c_{1}, \ldots, c_{n}\right), c_{i}=c_{i+1}+r$ for some nonnegative $r$. Now we compute the following:

$$
\begin{aligned}
\partial_{i}\left(x^{c(w)}\right) & =\partial_{i}\left(x_{1}^{c_{1}} \cdots x_{i}^{c_{i+1}+r} x_{i+1}^{c_{i+1}} \cdots x_{n}^{c_{n}}\right) \\
& =\left(x_{1}^{c_{1}} \cdots x_{i}^{c_{i+1}} x_{i+1}^{c_{i+1}} \cdots x_{n}^{c_{n}}\right) \cdot \partial_{i}\left(x_{i+1}^{r}\right) \\
& =\left(x_{1}^{c_{1}} \cdots x_{i}^{c_{i+1}} x_{i+1}^{c_{i+1}} \cdots x_{n}^{c_{n}}\right) \cdot\left(x_{i+1}^{r-1}+\text { lower terms }\right) \\
& =\left(x_{1}^{c_{1}} \cdots x_{i}^{c_{i+1}} x_{i+1}^{c_{i+1}+1+1} \cdots x_{n}^{c_{n}}\right)+\text { lower terms } \\
& =x^{c\left(w s_{i}\right)}+\text { lower terms. }
\end{aligned}
$$

(3) Any $x^{\alpha}$ has the property that there is a $w \in S_{n}$ for $n$ large enough such that $c(w)=\alpha$. Thus $x^{\alpha}$ can be written as a linear combination of $\mathfrak{S}_{w}$ and a sum of terms $\mathfrak{S}_{v}$ where $c(v) \subset c(w)$ by (2) and induction. Since the Artin monomials form a basis for the ring of polynomials in the $x_{i}$, so too must the $\mathfrak{S}_{w}$.

Corollary 13.5. We can write $\mathfrak{S}_{u} \mathfrak{S}_{v}==\sum b_{u v}^{w} \mathfrak{S}_{w}$ where the coefficients $b_{u v}^{w}=C_{u w_{0}, v w_{0}}^{w w_{0}}$ count intersection multiplicities. In particular, these coefficients are nonnegative integers.

Proposition 13.6. Let $f=\sum c_{i} x_{i}$ with $c_{i} \in \mathbb{Z}$. Then

$$
f \cdot \mathfrak{S}_{w}=\sum\left(c_{i}-c_{j}\right) \mathfrak{S}_{w t_{i j}}
$$

where the sum is taken over all $t_{i j}$ such that $\ell(w)=\ell\left(w t_{i j}\right)-1$.

Proof. Since $f \cdot \mathfrak{S}_{w}$ is homogeneous of degree $\ell(w)+1$, we can write

$$
f \cdot \mathfrak{S}_{w}=\sum_{v: \ell(v)=\ell(w)+1} \partial_{v}\left(f \cdot \mathfrak{S}_{w}\right) \cdot \mathfrak{S}_{v} .
$$

Let $a_{1} \cdots a_{p}$ be a reduced expression for $v$ and compute

$$
\begin{aligned}
\partial_{v}\left(f \cdot \mathfrak{S}_{w}\right) & =\partial_{a_{1}} \cdots \partial_{a_{p}}\left(f \cdot \mathfrak{S}_{w}\right) \\
& =\partial_{a_{1}} \cdots \partial_{a_{p-1}}\left(s_{a_{p}}(f) \partial_{a_{p}}\left(\mathfrak{S}_{w}\right)+\partial_{a_{p}}(f) \mathfrak{S}_{w}\right) \\
& =\cdots \\
& =s_{a_{1}} \cdots s_{a_{p}}(f) \partial_{a_{1}} \cdots \partial_{a_{p}}\left(\mathfrak{S}_{w}\right)+\sum_{r=1}^{p} s_{a_{1}} \cdots \partial_{a_{r}} \cdots s_{a_{p}}(f) \cdot \partial_{a_{1}} \cdots \hat{\partial}_{a_{r}} \cdots \partial_{a_{p}}\left(\mathfrak{S}_{w}\right) .
\end{aligned}
$$

Since $p=\ell(v)=\ell(w)+1, \partial_{a_{1}} \cdots \partial_{a_{p}}\left(\mathfrak{S}_{w}\right)=0$. In the second summation, the only nonzero terms will be those in which $s_{a_{1}} \cdots \widehat{s}_{a_{r}} \cdots s_{a_{p}}=w$. This implies that $s_{a_{1}} \cdots \widehat{s}_{a_{r}} \cdots s_{a_{p}} \cdot s_{a_{p}} \cdots s_{a_{r}} \cdots s_{a_{p}}=$ $w t_{i j}$ for $i<j$ such that $t_{i j}=s_{a_{p}} \cdots s_{a_{r}} \cdots s_{a_{p}}$. Thus $(i, j)=s_{a_{p}} \cdots s_{a_{r+1}}\left(a_{r}, a_{r+1}\right)$ and hence $s_{a_{1}} \cdots \partial_{a_{r}} \cdots s_{a_{p}}(f)=c_{i}-c_{j}$.
Corollary 13.7. (Monk's formula)

$$
\mathfrak{S}_{s_{r}} \cdot \mathfrak{S}_{w}=\sum_{\substack{1 \leq r<j \\ \ell(w)=\ell\left(w t_{i j}\right)-1}} \mathfrak{S}_{w t_{i j}}
$$

Proof. We know that $\mathfrak{S}_{s_{r}}=x_{1}+\ldots+x_{r}$ so we just apply the previous proposition and note that the coefficient $c_{i} \neq c_{j}$ if and only if $i \leq r<j$.
Example 13.8. Compute $\mathfrak{S}_{s_{3}} \cdot \mathfrak{S}_{3146527}$.
To do this, we need only determine which transpositions that interchange one of $\{3,1,4\}$ with one of $\{6,5,2,7\}$ increase the length of the permutation [3146527] by 1 . The only such transpositions are $4 \leftrightarrow 5,4 \leftrightarrow 6$, and $1 \leftrightarrow 2$. Thus

$$
\mathfrak{S}_{s_{3}} \cdot \mathfrak{S}_{3146527}=\mathfrak{S}_{3164527}+\mathfrak{S}_{3156427}+\mathfrak{S}_{3246517}
$$

Exercise 13.9. Do any six applications of Monk's formula.
Exercise 13.10. Prove that $\left[X_{w w_{0}}\right]\left[X_{s_{r} w_{0}}\right]=\sum_{1 \leq r<j}\left[X_{w t_{i j} w_{0}}\right]$.
Open Question 13.11. Expand $\mathfrak{S}_{u} \mathfrak{S}_{v}$ via some tree using iterated applications of Monk's formula.
Corollary 13.12. (Transition Equation) Suppose $w=v t_{r s}$ where $r$ is the position of the last descent of $w$ and $s$ is the largest position of $w$ such that $w_{s}<w_{r}$. Then

$$
\mathfrak{S}_{w}=x_{r} \cdot \mathfrak{S}_{v}+\sum \mathfrak{S}_{w^{\prime}}
$$

where the sum is taken over all $w^{\prime}=v t_{i r}$ such that $i<r$ and $\ell(w)=\ell\left(w^{\prime}\right)$.
Proof. The polynomial $x_{r}$ is the simplest linear homogeneous polynomial of degree 1 that we can ask for, so by the product formula,

$$
x_{r} \cdot \mathfrak{S}_{v}=\mathfrak{S}_{w}-\sum \mathfrak{S}_{v t_{i r}}
$$

where the latter sum is taken over all $i<r$ such that $\ell(v)=\ell\left(v t_{i r}\right)-1$.
Remark 13.13. This is a very efficient way to multiply out Schubert polynomials without eating up as much computer memory; it is a sort of "depth-first" multiplication.

## 14. On "Flags, Schubert Polynomials, Degeneracy Loci, and Determinantal Formulas" By William Fulton.

This lecture was given by Ashesh Bakshi and Steve Klee.
14.1. Fiber Bundles. We need to introduce some definitions and facts about fiber bundles. More detailed information can be found in Milnor and Stasheff's Characteristic Classes.

Definition 14.1. Let $B$ be a topological space. A real vector bundle $\xi$ of rank $n$ (or n-plane bundle) over $B$ consists of a topological space $E=E(\xi)$ and a continuous surjective map $p: E \rightarrow B$ with the following properties:
(1) For each $b \in B$, the set $p^{-1}(b)$ has the structure of a real $n$-dimensional vector space.
(2) The bundle is locally trivial, i.e. each point $b \in B$ has a neighborhood $U$ and a homeomorphism $h: U \times \mathbb{R}^{n} \rightarrow p^{-1}(U)$ such that for each $b^{\prime} \in U$, the correspondence $v \mapsto h\left(b^{\prime}, v\right)$ defines an isomorphism from $\mathbb{R}^{n}$ to $p^{-1}\left(b^{\prime}\right)$ and the following diagram commutes: where $\pi_{1}$ is projection onto the first coordinate.
We say $B$ is the base space of the bundle and $E$ is the total space of the bundle. The map $p$ is called the projection map and the set $p^{-1}(b)$ is called the fiber over $b$, and is occasionally denoted $F_{b}(\xi)$.

Similarly, if $F$ is a topological space, we can define a fiber bundle with fiber $F$ to be a continuous surjection $p: E \rightarrow B$ such that for each $b \in B, p^{-1}(b)$ is homeomorphic to $F$ and such that local triviality is satisfied.

Remark 14.2. For general fiber bundles, we may require additional structure such as a Lie group action on $F$. We omit these requirements at this point for the sake of simplicity. For example, in the case of a vector bundle, the Lie group $G L_{n}(\mathbb{R})$ acts transitively on each fiber.

Remark 14.3. Depending on what we want to emphasize, we may simply say that $p: E \rightarrow B$ is a fiber bundle with fiber $F$ rather than saying $\xi$ is a fiber bundle with total space $E$, base space $B$, projection map $p$, and fiber $F$.

Example 14.4. Any topological space $F$ can be viewed as a fiber bundle over a space consisting of a single point $\{*\}$.

Example 14.5. If $\xi$ is an $n$-plane bundle with projection $p: E \rightarrow B$ such that $E$ is homeomorphic to $B \times \mathbb{R}^{n}$, we say that $\xi$ is a trivial bundle .

Example 14.6. If $M$ is a smooth $n$-manifold, the tangent bundle $T M$ is a smooth $n$-plane bundle over M.

Definition 14.7. If $p: E \rightarrow B$ is a bundle, a section of the bundle is a map $s: B \rightarrow E$ such that $p \circ s$ is the identity map on $B$.

Sometimes it is important to construct new vector bundles from known vector bundles. The following constructions will be of particular importance:

Definition 14.8. Suppose $\xi$ is a fiber bundle with projection $p: E \rightarrow B$. If $U$ is a subset of $B$, we define the restriction of $\xi$ to $U$, denoted $\xi \mid U$ to be the bundle with total space $E(\xi \mid U)=p^{-1}(U)$, base space $U$ (both with the subspace topology), and projection map given by the restriction of $p$ to $E(\xi \mid U)$. It is trivial to check that this is indeed a fiber bundle.

For our purposes, we will focus primarily on bundles whose fibers are vector spaces, projective spaces, and flag manifolds.

Definition 14.9. Suppose $\xi$ is a fiber bundle with fiber $F$ and projection $p: E \rightarrow B$. Let $X$ be any space. Given a continuous map $f: X \rightarrow B$, we construct the induced bundle or pullback bundle $f^{*} \xi$ over $X$ as follows. The total space $E\left(f^{*} \xi\right)$ is the subset of $X \times E$ consisting of all pairs $(x, v)$ such that $f(x)=p(v)$. The projection map $p_{1}: E\left(f^{*} \xi\right) \rightarrow X$ sends $(x, v) \mapsto x$. There is also an induced map $\hat{f}: E\left(f^{*} \xi\right) \rightarrow E(\xi)$ sending $(x, v) \mapsto v$ so that the following diagram commutes.

In the specific case that $\xi$ is an $n$-plane bundle, notice that the map $\hat{f}$ is a linear isomorphism on each fiber of $f^{*} \xi$. This motivates the following definition.

Definition 14.10. Suppose $p: E \rightarrow B$ is an $n$-plane bundle and $p^{\prime}: E^{\prime} \rightarrow B^{\prime}$ is an $m$-plane bundle. A bundle map is a pair of continuous maps $F: E \rightarrow E^{\prime}$ and $f: B \rightarrow B^{\prime}$ such that $p \circ F=f \circ p^{\prime}$ and such that $F$ is $\mathbb{R}$-linear when restricted to any fiber of $p$. We will usually only describe the map $F: E \rightarrow E^{\prime}$ and use the requisite commutative diagram to induce the map $f: B \rightarrow B^{\prime}$.

If $p: E(\xi) \rightarrow B$ is a vector bundle with fiber $\mathbb{R}^{n}$, it may be natural to ask questions about vector subspaces of the fibers of $p$. A subbundle $\eta \subset \xi$ is a $k$-plane bundle over $B$ whose total space $E(\eta)$ is a subspace of $E(\xi)$ such that over each point $b \in B, F_{b}(\eta)$ is a $k$-dimensional subspace of $F_{b}(\xi)$. This gives rise to one additional construction:

Definition 14.11. Suppose $\xi$ is an $n$-plane bundle over $B$ and $\eta \subset \xi$ is a $k$-plane subbundle. The quotient bundle $\xi / \eta$ is an $(n-k)$-plane bundle over $B$. Over each point $b \in B, F_{b}(\xi / \eta)$ is the vector space quotient $F_{b}(\xi) / F_{b}(\eta)$.

Remark 14.12. As before, if we do not wish to emphasize the subbundle $\eta \subset \xi$, we may simply say that $E^{\prime} \subset E$ is a $k$-plane subbundle of $E$. In saying this, we mean to indicate that there is a subbundle $\eta$ of $\xi$ whose total space is $E^{\prime}$, and whose projection map is $p \mid E^{\prime}$.

Example 14.13. Suppose $p: E \rightarrow B$ is an $n$-plane bundle. We form a new projective bundle $\mathbb{P}(E)$ over $B$ such that $F_{b}(\mathbb{P}(E))=\mathbb{P}\left(F_{b}(E)\right)$, i.e. the fiber over $b$ in $\mathbb{P}(E)$ is the real projective space on the vector space $F_{b}(E)$. We can also view $\mathbb{P}(E)$ as a rank $(n-1)$-vector bundle over $E$.

Definition 14.14. Suppose $\xi: E \xrightarrow{p} B$ and $\eta: E^{\prime} \xrightarrow{p^{\prime}} B^{\prime}$ are $n$ - and $m$-plane bundles over a common base space $B$. We obtain the total space of new bundle $\operatorname{Hom}(\xi, \eta)$ by decreeing that for each $b \in B$, $F_{b}(\operatorname{Hom}(\xi, \eta))$ is the $m n$-dimensional vector space of linear maps from $p^{-1}(b) \rightarrow p^{\prime-1}(b)$. (It takes a fair amount of work to show that this space can be given a topology that gives it the structure of a vector bundle.) Similarly we can obtain bundles $\xi \otimes \eta$ (respectively $\xi \oplus \eta$ ) by applying the tensor product (resp. direct sum) functors to the fiber over each point. (The latter bundle is called the Whitney sum of $\xi$ and $\eta$.)
14.2. Degeneracy Loci. In general, subbundles of vector bundles are not as easy to find as one might expect. Using cohomology theory, it is possible to show that the tangent space of real projective space $\mathbb{P}^{n}$ admits a line subbundle (i.e. a subbundle of rank 1 ) if and only if $n$ is odd. Moreover, the tangent space to complex projective space never admits a line subbundle.

Now, however, we will bestow upon ourselves a most fortunate situation that is analogous to a decomposition of a vector space into a flag of subspaces.
Definition 14.15. Suppose $p: E \rightarrow B$ is a vector bundle of rank $n$. We say $E$ is a filtered vector bundle (over $B$ ) if there is a chain of subbundles $E_{1} \subset E_{2} \subset \ldots \subset E_{n}=E$ such that $E_{i}$ has rank $i$ over $B$.

Dually, if $E=E_{\bullet}$ is a filtered vector bundle, there is a descending chain of surjective bundle maps

$$
E=Q_{n} \rightarrow Q_{n-1} \rightarrow \ldots \rightarrow Q_{2} \rightarrow Q_{1}
$$

where $Q_{i}=E / E_{n-i}$ has rank $i$. As a means of notation, in this case we occasionally abuse notation and write $E=E_{n} \rightarrow E_{n-1} \rightarrow \ldots \rightarrow E_{1}$ to indicate this chain of quotient bundles.

With this notation in place, suppose $X$ is a variety, and $E_{\bullet}$ and $F_{\bullet}$ are filtered vector bundles of rank $n$ over $X$. If $h: E \rightarrow F$ is a bundle map then we can form a chain of maps

$$
E_{1} \subset E_{2} \subset \ldots \subset E_{n} \xrightarrow{h} F_{n} \rightarrow F_{n-1} \rightarrow \ldots \rightarrow F_{1} .
$$

For any $1 \leq p, q \leq n$, we will denote the composition of bundle maps

$$
E_{p} \hookrightarrow E_{p+1} \hookrightarrow \ldots \hookrightarrow E_{n} \xrightarrow{h} F_{n} \rightarrow F_{n-1} \rightarrow \ldots \rightarrow F_{q}
$$

as $E_{p} \rightarrow F_{q}$.
Definition 14.16. For a permutation $w \in S_{n}$, we define the degeneracy locus $\Omega_{w}=\Omega_{w}(h)=\Omega_{w}\left(h, E_{\bullet}, F_{\bullet}\right)$ to be the subset of points $x \in X$ such that the restricted map $E_{p}\left|x \rightarrow F_{q}\right| h(x)$ has rank at most $r_{w}(q, p)$ for all $1 \leq p, q \leq n$. As a means of convenience we abbreviate this by saying that $\Omega_{w}$ is defined by the conditions that $\operatorname{rank}\left(E_{p} \rightarrow F_{q}\right) \leq r_{w}(q, p)$ for all $p, q$.

As in the case of flag varieties, we can use Fulton's essential set to determine a smaller generating set for $\Omega_{w}$.
Definition 14.17. In the theorem below we are using the definition that the Essential Set is the collection of southeast corners of connected components in the permutation matrix for $w$ after canceling south and east from the entries $(i, w(i))$.

Proposition 14.18. The variety $\Omega_{w}$ is defined by the conditions

$$
\operatorname{rank}\left(E_{p} \rightarrow F_{q}\right) \leq r_{w}(q, p)
$$

for all $p$ and $q$ in $\{1, \ldots n\}$ such that $(q, p) \in \operatorname{Ess}(w)$.
Proof. If we pick a point $x \in \Omega_{w}$, then we know the condition that the rank of the map $\left.\left.E_{p}\right|_{x} \rightarrow F_{q}\right|_{x}$ is bounded by $r_{w}(q, p)$ for all $p, q$ is equivalent to the rank being bounded for all pairs $(q, p)$ in the essential set of $w$. If $U \subset X$ is a trivializing neighborhood of $x$ for all $E_{i}$ (i.e. $E_{i} \mid U$ is homeomorphic to the trivial bundle), then the result certainly holds over $U \cap \Omega_{w}$ where $E_{p}$ and $F_{q}$ have the same behavior over all fibers. Thus the assertion holds in a neighborhood of each point in $\Omega_{w}$.
14.3. Flag Bundles. Suppose that $p: E \rightarrow X$ is a rank $n$ vector bundle over a variety $X$. The flag bundle of $E$, denoted $F l(E)$ is a vector bundle of rank $n(n-1) / 2$ over $X$ with projection map $\rho$ : $F l(E) \rightarrow X$. Since $\rho$ is a continuous map to the base space of a bundle, using the pullback construction of definition 14.9 we can form an induced bundle $\rho^{*} E$. The flag bundle comes equipped with a universal flag of subbundles $U_{\bullet}$ of $\rho^{*} E$, with each $U_{i}$ a vector bundle of rank $i$. The situation is as follows:

The flag bundle is universal in the following sense: if $f: Y \rightarrow X$ is a continuous map such that $f_{\tilde{f}}^{*} E$ has a complete flag of subbundles $V_{1} \subset V_{2} \subset \ldots \subset V_{n-1} \subset f^{*} E$, then there is a unique map $\tilde{f}: Y \rightarrow F l(E)$ such that $\tilde{f}^{*} U_{i}=V_{i}$ as subbundles of $f^{*} E$ for all $i$. Finally, the universal flag of subbundles of $F l(E)$ gives a universal flag of quotient bundles

$$
\rho^{*} E \rightarrow Q_{n-1} \rightarrow \ldots \rightarrow Q_{1}
$$

where $Q_{i}=\rho^{*} E / U_{n-i}$ for all $i$.
We can construct $F l(E)$ as a sequence of projective bundles (fiber bundles whose fiber is some projective space). Start with the canonical projection map $\rho_{1}: \mathbb{P}(E) \rightarrow X$. The $n$-plane bundle $\rho_{1}^{*}(E)$ admits a universal line subbundle $M_{1}$ : for $L \in \mathbb{P}(E)$, let $b=\rho_{1}(L)$. As a point in $\mathbb{P}(E), L$ represents a line
through the origin in $F_{b}(E)$, and the fiber over $L$ in $\rho_{1}^{*}(E)$ is canonically isomorphic to $F_{b}(E)$. Over $L$, take $M_{1}$ to be the collection of points $x \in F_{b}(E)$ such that $x$ lies on $L$. Now we can take $\mathbb{P}\left(\rho_{1}^{*} E / M_{1}\right)$, which is a $\mathbb{P}^{n-2}$-bundle over $\mathbb{P}(E)$ with projection map $\rho_{2}$. As before, the pullback bundle $\rho_{2}^{*}\left(\rho_{1}^{*} E\right)$ admits a universal line subbundle $M_{2}$. Iterating this process, we obtain $F l(E)$.

We attempt to gain some intuition for the universal flag bundle with the following example, which will be useful later.
Example 14.19. Suppose $V$ is an $n$-dimensional vector space, viewed as the trivial bundle over the one point space $\left\{b_{0}\right\}$ with projection $p: V \rightarrow\left\{b_{0}\right\}$. The flag bundle $F l(V)$ is the flag manifold. Moreover, we can explicitly describe the universal flag of subbundles $U_{\bullet}$ of $\rho^{*} E$ : for a complete flag $W_{\bullet}$ on $V$, $U_{i} \mid W_{\bullet}=W_{i}$.
Proof. We run through the above algorithm to construct the universal flag bundle $F l(V)$. We start by taking $\mathbb{P}(V)$, the space of lines $V_{1}$ through the origin in $V$ with projection $\rho_{1}$ onto $\left\{b_{0}\right\}$. In this case, $\rho_{1}^{*} V$ is the trivial bundle $\mathbb{P}(V) \times V$ with the universal line subbundle $M_{1}=\left\{\left(V_{1}, x\right) \in \mathbb{P}(V) \times V: x \in V_{1}\right\}$.

Next, we look at the projective space on the quotient bundle $(\mathbb{P}(V) \times V) / M_{1}$, and consider the projection map $\rho_{2}: \mathbb{P}\left(\rho^{*} V / M_{1}\right) \rightarrow \mathbb{P}(V)$. For $V_{1} \in \mathbb{P}(V)$, we can explicitly compute $\rho_{2}^{-1}\left(V_{1}\right)$ : the fiber over $V_{1}$ in $\rho^{*} V / M_{1}$ is the vector space $V / V_{1}$, so the fiber over $V_{1}$ by $\rho_{2}$ is $\mathbb{P}\left(V / V_{1}\right)$. But a line $V_{2}$ through the origin in $V / V_{1}$ can be identified with a plane in $V_{2} \subset V$ that contains $V_{1}$. Thus we can identify $\mathbb{P}\left(\rho_{1}^{*} V / M_{1}\right)$ with the collection of partial flags $V_{1} \subset V_{2}$ in $V$ with projection map $\left(V_{1} \subset V_{2}\right) \stackrel{\rho_{1}}{\mapsto} V_{1}$.

Under this identification, $\rho_{2}^{*}\left(\rho_{1}^{*}(V)\right)$ can be viewed as the set of triples $\left(V_{2}, V_{1}, x\right)$ where $V_{1} \in \mathbb{P}(V)$ is a line through the origin in $V$ and $V_{2}$ is a 2-plane containing $V$. This has a universal two-plane subbundle $M_{2}$ whose fiber over the pair $\left(V_{2}, V_{1}\right) \in \mathbb{P}\left(\rho_{1}^{*}(V) / M_{1}\right)$ is the set of points $x \in V_{2}$, which can be used to define a universal line subbundle $M_{2} / \rho_{2}^{*} M_{1}$ over $\mathbb{P}\left(\rho^{*} V / M_{1}\right)$.

From here, we appeal to the induction gods.
Suppose our bundle $E \xrightarrow{p} X$ admits a complete flag of subbundles $E_{\bullet}=E_{1} \subset E_{2} \subset \ldots \subset E_{n-1} \subset$ $E$. Then $\rho^{*} E$ admits a complete flag of subbundles $\rho^{*} E_{1} \subset \ldots \subset \rho^{*} E_{n}$ is a complete flag of subbundles. Thus, over $F l(E)$ we have a map of filtered vector bundles

$$
\rho^{*} E_{1} \subset \ldots \subset \rho^{*} E_{n-1} \subset \rho^{*} E=\rho^{*} E \rightarrow Q_{n-1} \rightarrow \ldots \rightarrow Q_{1} .
$$

For $w \in S_{n}$, we can form the degeneracy locus $\Omega_{w}$ in $F l(E)$, denoted $\Omega_{w}\left(E_{\bullet}\right)$. We will show that if $X$ is an irreducible variety, then $\Omega_{w}\left(E_{\bullet}\right)$ is an irreducible subvariety of $F l(E)$ of codimension $l(w)$.
14.4. Schubert Varieties. We will now work in the specific case that the variety $X$ is a single point. Here a bundle over $X$ is an $n$-dimensional vector space $V$; and, as noted in example 14.19, the flag bundle $F l(V)=F l(n)$ is just the flag manifold with universal subbundle $U_{\text {• }}$. Suppose we fix a flag $V_{\bullet}$ of vector spaces in $V$. For a permutation $w$ in $S_{n}$, we look at the open Schubert cell

$$
X_{w}^{o}=X_{w}^{o}\left(V_{\bullet}\right)=\left\{W_{\bullet} \in F l(V): \operatorname{dim}\left(W_{q} \cap V_{p}\right)=r_{w}(q, p) ; p, q \in[n]\right\}
$$

and the closed Schubert variety

$$
X_{w}=\left\{W_{\bullet} \in F l(V): \operatorname{dim}\left(W_{q} \cap V_{p}\right) \geq r_{w}(q, p) ; p, q \in[n]\right\} .
$$

Now we consider $\rho: F l(V) \rightarrow X=\{*\}$. As before, $\rho^{*} V$ is the trivial bundle $F l(V) \times V$. Using our fixed flag $V_{\bullet}$, we get a flag of subbundles of $F l(V) \times V$ :

$$
F l(V) \times V_{1} \subset F l(V) \times V_{2} \subset \ldots \subset F l(V) \times V_{n-1} \subset F l(V) \times V .
$$

We also have a universal flag of quotient bundles

$$
F l(V) \times V \rightarrow Q_{n-1} \ldots \rightarrow Q_{1}
$$

where $Q_{i}$ is the quotient bundle $(F l(V) \times V) / U_{n-i}$. As in the previous section, $F l(V) \times V_{i}=\rho^{*} V_{i}$, so we can look at the degeneracy locus $\Omega_{w}\left(V_{\bullet}\right)$.

Proposition 14.20. For any $w \in S_{n}, \Omega_{w}\left(V_{\bullet}\right)=X_{w \cdot w_{0}}$.
Proof. Recall that the universal subbundle $U_{\bullet}$ of $F l(V) \times V$ satisfies $U_{i} \mid W_{\bullet}=W_{i}$ for $W_{\bullet} \in F l(V)$. To compute the degeneracy locus $\Omega_{w}\left(V_{\bullet}\right)$, we need to compute the ranks of the maps $\left(\rho^{*} V\right)_{p} \rightarrow Q_{q}$. For a complete flag $W_{\bullet}$, we can compute $\left(\rho^{*} V\right)_{p} \mid W_{\bullet}=V_{p}$ and

$$
Q_{q}\left|W_{\bullet}=\left(\rho^{*} V / U_{n-q}\right)\right| W_{\bullet}=\rho^{*} V\left|W_{\bullet} / U_{n-q}\right| W_{\bullet}=V / W_{n-q} .
$$

Thus over $W_{\bullet}$, the map $\left(\rho^{*} V\right)_{p} \rightarrow Q_{q}$ is simply the restriction of the quotient map $V_{p} \rightarrow V / W_{n-q}$.
Now $\Omega_{w}\left(V_{\bullet}\right)$ is the subvariety of $F l(V)$ of all flags $W_{\bullet}$ such that the map $V_{p} \rightarrow Q_{q}$ has rank at most $r_{w}(q, p)$. Over a single flag, this is just a map out of $V_{p}$, so saying its rank is at most $r_{w}(q, p)$ is equivalent to saying that its kernel has dimension at least $p-r_{w}(q, p)$. The kernel of the map $V_{p} \rightarrow V / W_{n-q}$ is clearly $V_{p} \cap W_{n-q}$. Thus $\Omega_{w}\left(V_{\bullet}\right)$ is the collection of complete flags such that for all $p$ and $q$,

$$
\begin{aligned}
\operatorname{dim}\left(W_{q} \cap V_{p}\right) & \geq p-r_{w}(n-q, p) \\
& =p-\#\{i \leq n-q: w(i) \leq p\} \\
& =p-\#\left\{i \leq p: w^{-1}(i) \leq n-q\right\} \\
& =\#\left\{i \leq p: w^{-1}(i)>n-q\right\} \\
& =\#\left\{i \leq p: w_{0} w^{-1}(i) \leq q\right\} \\
& =\#\left\{i \leq q: w \cdot w_{0}(i) \leq p\right\} \\
& =r_{w \cdot w_{0}}(q, p) .
\end{aligned}
$$

But this is precisely $X_{w} \cdot w_{0}$.
14.5. Double Schubert polynomials. To state the main theorem in this paper we will need a slight generalization of Schubert polynomials.

Definition 14.21. Let $A$ be a commutative ring. Recall that for each $1 \leq i \leq n-1$, the divided difference operator on $A\left[x_{1}, \ldots, x_{n}\right]$ takes a polynomial $P$ to

$$
\begin{equation*}
\partial_{i} P=\frac{P\left(x_{1}, \ldots, x_{n}\right)-P\left(x_{1}, \ldots, x_{i-1}, x_{x+i}, x_{i}, x_{i+2}, \ldots, x_{n}\right)}{x_{i}-x_{i+1}} \tag{14.1}
\end{equation*}
$$

Definition 14.22. Let $s_{i}$ denote the transposition exchanging $i$ and $i+1$. The double Schubert polynomial for the permutation $w$ is a homogeneous polynomial in $2 n$ variables of degree $\ell(w)$. Write $w$ as the product $w_{0} \cdot s_{i_{1}} \cdots \cdot s_{i_{r}}$, where $r=\ell\left(w_{0}\right)-\ell(w)$, and set

$$
\begin{equation*}
\mathfrak{S}_{w}\left(x_{1}, \ldots, x_{n}, y_{1}, \ldots, y_{n}\right)=\partial_{i_{r}} \circ \cdots \circ \partial_{i_{1}}\left(\prod_{i+j \leq n}\left(x_{i}-y_{j}\right)\right) \tag{14.2}
\end{equation*}
$$

Here the $y_{i}$ are to be regarded as constants, i.e., the operators are defined on the ring $A\left[x_{1} \ldots, x_{n}\right]$, where $A$ is the ring $\mathbb{Z}\left[y_{1}, \ldots, y_{n}\right]$.
Remark 14.23. The above definition is independent of the choice of $s_{i_{1}}, \ldots, s_{i_{j}}$ and we can recover the ordinary Schubert polynomials as $\mathfrak{S}_{w}\left(x_{1}, \ldots, x_{n}\right)=\mathfrak{S}_{w}\left(x_{1}, \ldots, x_{n}, 0, \ldots, 0\right)$.

Conversely, the double Schubert polynomials can be expressed in terms of ordinary Schubert polynomials. For brevity, write $x$ and $y$ in place of $x_{1}, \ldots, x_{n}$ and $y_{1} \ldots, y_{n}$ (respectively.)
Fact 14.24. For any permutation $w \in \mathbb{S}_{n}$,

$$
\mathfrak{S}_{w}(x, y)=\sum(-1)^{\ell(v)} \mathfrak{S}_{u}(x) \mathfrak{S}_{v}(y)
$$

where the sum is over pairs of permutations $(u, v)$ such that $v^{-1} \cdot u=w$ and $\ell(u)+\ell(v)=\ell(w)$.

Remark 14.25. Suppose $A$ is a ring, $c_{1}, \ldots, c_{n}$ are elements of $A$, and $I$ is the ideal of $A[x]$ generated by the polynomials $e_{i}(x)-c_{i}$ where the $e_{i}$ are the elementary symmetric functions. Notice that the operators $\partial_{w}$ map $I$ to itself, and hence descend to well-defined operators, also denoted by $\partial_{w}$ from $A[x] / I \rightarrow A[x] / I$.
14.6. Chern classes. Let $\omega: E \xrightarrow{p} B$ be a complex n-plane bundle. There are certain characteristic classes in the singular cohomology ring of the base space. Working with $\mathbb{Z} /(2)$ coefficients, one has the Stiefel-Whitney classes, which can be used to study diverse questions: When can a smooth closed $m$-manifold be realized as the boundary of an $(m+1)$-manifold? For which $m$ is the projective space $\mathbb{P}^{m}$ parallelizable? For which $m$ can a manifold be immersed in $\mathbb{R}^{m}$ ?

Working instead with integer coefficients, we have for each $0 \leq i \leq n$ a Chern class $c_{i}(\omega) \in$ $H^{2 i}(B ; \mathbb{Z})$. As is the case for the Stiefel-Whitney classes, they play an important role in the study of the topology of vector bundles and manifolds, and as we'll see below they crop up when we study the strucutre of the Chow ring as well. A proof of the existence of Chern classes is beyond the scope of this lecture; see Fulton's Young Tableax, chapter three of Fulton's Intersection Theory, or Milnor and Stasheff's Characteristic Classes for different constructions. For now we content ourselves with listing a few of their properties, not all of which will be used here:

Fact 14.26. (Vanishing) For any complex $n$-plane bundle $\omega, c_{i}(\omega)=0$ for $i>n$.
Fact 14.27. (Naturality) Given complex $n$-plane bundles $E \rightarrow B$ and $E^{\prime} \rightarrow B^{\prime}$, and maps $F: E \rightarrow E^{\prime}$ and $f: B \rightarrow B^{\prime}$ so that the pair $(F, f)$ is a bundle map (cf. definition 14.10), then for each $i, c_{i}(\omega)=$ $f^{*} c_{i}\left(\omega^{\prime}\right)$.

Definition 14.28. A short exact sequence of vector bundles is a sequence

$$
0 \rightarrow \omega^{\prime} \xrightarrow{f} \omega \xrightarrow{g} \omega^{\prime \prime} \rightarrow 0
$$

of bundles over a common base space $B$ such that: $f$ and $g$ induce the identity map on the base space, and for each $b \in B$, the restriction to the fiber over $b$

$$
\left.\left.\left.0 \rightarrow \omega^{\prime}\right|_{b} \xrightarrow{\left.f\right|_{b}} \omega\right|_{b} \xrightarrow{\left.g\right|_{b}} \omega^{\prime \prime}\right|_{b} \rightarrow 0
$$

is a short exact sequence of linear maps between vector spaces.
Fact 14.29. (Whitney formula) If $0 \rightarrow \omega^{\prime} \rightarrow \omega \rightarrow \omega^{\prime \prime} \rightarrow 0$ is a short exact sequence of vector bundles over $B$, where $B$ is paracompact (or, if you prefer to avoid lengthy topological definitions, where $B$ is a variety), then $c_{k}(\omega)=\sum_{i+j=k} c_{i}\left(\omega^{\prime}\right) \cup c_{j}\left(\omega^{\prime \prime}\right)$.

It is often convenient to work in the ring $H^{\mathrm{H}}(B ; \mathbb{Z})$ of all formal infinite series $a=a_{0}+a_{1}+\cdots$ with $a_{i} \in H^{i}(B)$, with the product operation $H^{\Pi}(B ; \mathbb{Z}) \times H^{\Pi}(B ; \mathbb{Z}) \xrightarrow{*} H^{\Pi}(B ; \mathbb{Z})$ given by:

$$
\left(a_{0}+a_{1}+\cdots\right) *\left(b_{0}+b_{1}+\cdots\right)=\sum_{l=0}^{\infty} \sum_{j+k=l}\left(a_{j} \cup b_{k}\right)=\left(a_{0} \cup b_{0}\right)+\left(a_{1} \cup b_{0}+a_{0} \cup b_{1}\right)+\cdots
$$

In this ring define the total Chern class $c(\omega)=1+c_{1}(\omega)+c_{2}(\omega)+\cdots$. Two useful consequences are the following: the total Chern class of a vector bundle is a unit in the ring and the Whitney formula becomes simply $c(\omega)=c\left(\omega^{\prime}\right) * c\left(\omega^{\prime \prime}\right)$.
Definition 14.30. Let $G_{k}\left(\mathbb{C}^{\infty}\right)$ denote the complex Grassmannian of $k$-planes through the origin in $\mathbb{C}^{\infty}$. Notice in particular that when $k=1, G_{k}\left(\mathbb{C}^{\infty}\right)=\mathbb{C} P^{\infty}$. The universal $k$-plane bundle $\gamma^{k}$ over $G_{k}\left(\mathbb{C}^{\infty}\right)$ is the bundle whose total space is the subset of $G_{k}\left(\mathbb{C}^{\infty}\right) \times \mathbb{C}^{\infty}$ consisting of all pairs $(V, x)$ such that $V$ is a $k$-plane in $\mathbb{C}^{\infty}$ and $x \in V$.

Fact 14.31. If $p: E \rightarrow B$ is a complex $k$-plane bundle, then there is a map $f: B \rightarrow G_{k}\left(\mathbb{C}^{\infty}\right)$ such that $E=f^{*} \gamma^{k}$.
Fact 14.32. The cohomology ring $H^{*}\left(\mathbb{C} P^{\infty} ; \mathbb{Z}\right)$ is a polynomial ring generated in degree 2 by $c_{1}\left(\gamma^{1}\right)$. As a corollary, $H^{i}\left(\mathbb{C} P^{\infty}\right)=0$ for all odd integers $i$.

We need to introduce one final piece of notation. Suppose $\xi$ is an $n$-plane bundle over $X$ and $\eta$ is an $m$-plane bundle over $Y$. The external tensor product of $\xi$ and $\eta$ is the bundle $\xi \hat{\otimes} \eta$ over $X \times Y$ whose fiber over $(x, y)$ is $F_{x}(\xi) \otimes F_{y}(\eta)$.
Lemma 14.33. Suppose $\xi$ and $\eta$ are line bundles over a common base space $B$. Then $c_{1}(\xi \otimes \eta)=$ $c_{1}(\xi)+c_{1}(\eta)$.

Proof. (The reader is invited to skip this proof if so inclined; only the statement will be of use below.) Recall the Kunneth formula for cohomology: if $X$ and $Y$ are spaces, then the cross product gives an isomorphism $H^{k}(X \times Y ; \mathbb{Z}) \rightarrow \bigoplus_{i+j=k} H^{i}(X ; \mathbb{Z}) \otimes H^{j}(Y ; \mathbb{Z})$ in cases (like this one) where all the cohomology groups are free.

Let $d: B \rightarrow B \times B$ be the diagonal map $x \mapsto(x, x)$. By Fact 14.31, there are maps $f, g: B \rightarrow \mathbb{C} P^{\infty}$ such that $f^{*} \gamma^{1}=\xi$ and $g^{*} \gamma^{1}=\eta$. We can certainly map $f \times g: B \times B \rightarrow \mathbb{C} P^{\infty} \times \mathbb{C} P^{\infty}$, and it is clear that $(f \times g)^{*}\left(\gamma^{1} \hat{\otimes} \gamma^{1}\right)=\xi \hat{\otimes} \eta$. Moreover, $d^{*}(\xi \hat{\otimes} \eta)=\xi \otimes \eta$. This is illustrated in the following commutative diagram:

By the Kunneth formula, $H^{2}\left(\mathbb{C} P^{\infty} \times \mathbb{C} P^{\infty} ; \mathbb{Z}\right)$ is isomorphic to $\left(H^{2}\left(\mathbb{C} P^{\infty}\right) \otimes H^{0}\left(\mathbb{C} P^{\infty} ; \mathbb{Z}\right)\right) \oplus$ $\left(H^{0}\left(\mathbb{C} P^{\infty} ; \mathbb{Z}\right) \otimes H^{2}\left(\mathbb{C} P^{\infty} ; \mathbb{Z}\right)\right) ; H^{1}\left(\mathbb{C} P^{\infty}\right)=0$ by Fact 14.32. Thus $c_{1}\left(\gamma^{1} \hat{\otimes} \gamma^{1}\right)$ can be written as $a c_{1}\left(\gamma^{1}\right) \times 1+1 \times b c_{1}\left(\gamma^{1}\right)$ for some $a, b \in \mathbb{Z}$. Fix any point $*$ in $\mathbb{C} P^{\infty}$. To determine the constants $a$ and $b$, we consider the maps $\iota_{1}: \mathbb{C} P^{\infty} \times\{*\} \rightarrow \mathbb{C} P^{\infty} \times \mathbb{C} P^{\infty}$ sending $(P, *) \mapsto(P, *)$ and $\iota_{2}:\{*\} \times \mathbb{C} P^{\infty} \rightarrow \mathbb{C} P^{\infty} \times \mathbb{C} P^{\infty}$ defined similarly. We can identify $\iota_{1}^{*}\left(\gamma^{1} \hat{\otimes} \gamma^{1}\right)$ with $\gamma^{1}$ and hence $\iota_{1}^{*}\left(c_{1}\left(\gamma^{1} \hat{\otimes} \gamma^{1}\right)\right)=c_{1}\left(\gamma^{1}\right) \times 1$. Using the equation for $c_{1}\left(\gamma^{1} \hat{\otimes} \gamma^{1}\right)$ from the Kunneth formula, $\iota_{1}^{*}\left(a c_{1}\left(\gamma^{1}\right) \times\right.$ $\left.1+1 \times b c_{1}\left(\gamma^{1}\right)\right)=a c_{1}\left(\gamma^{1}\right) \times 1+0$. Thus $a=1$. Similarly, $b=1$.

Now by naturality, we can compute $c_{1}(\xi \otimes \eta)$ :

$$
\begin{aligned}
c_{1}(\xi \otimes \eta) & =d^{*}\left((f \times g)^{*}\left(c_{1}\left(\gamma^{1} \hat{\otimes} \gamma^{1}\right)\right)\right) \\
& =d^{*}\left((f \times g)^{*}\left(c_{1}\left(\gamma^{1}\right) \times 1+1 \times c_{1}\left(\gamma^{1}\right)\right)\right) \\
& =d^{*}\left(f^{*}\left(c_{1}\left(\gamma^{1}\right)\right) \times g^{*} 1+f^{*} 1 \times g^{*}\left(c_{1}\left(\gamma^{1}\right)\right)\right) \\
& =d^{*}\left(c_{1}(\xi) \times 1+1 \times c_{1}(\eta)\right) \\
& =c_{1}(\xi) \cup 1+1 \cup c_{1}(\eta) \\
& \left.=c_{1}(\xi)+c_{1}(\eta)\right)
\end{aligned}
$$

14.7. A Giambelli formula for flag bundles. The key result in the paper is the following:

Theorem 14.34. For any complete flag $E_{\bullet}$ of subbundles in a vector bundle $E$ of rank $n$ on a nonsingular variety $X$ and any $w \in S_{n}$, the class of $\Omega_{w}\left(E_{\bullet}\right)$ in $A^{\ell(w)}(F l(E))$ is given by the formula

$$
\begin{equation*}
\left[\Omega_{w}\left(E_{\bullet}\right)\right]=\mathfrak{S}_{w}\left(x_{1}, \ldots, x_{n} \cdot y_{1}, \ldots, y_{n}\right) \tag{14.3}
\end{equation*}
$$

where $x_{i}$ is the first Chern class of $\operatorname{ker}\left(Q_{i} \rightarrow Q_{i-1}\right)$ and $y_{i}$ is the first Chern class of $E_{i} / E_{i-1}$.

This proof depends on a few key facts, which we will state without proof.
First of all, let $A=A^{*} X$ be the Chow ring of $X$. There is a canonical (surjective) homomorphism from $A\left[x_{1}, \ldots, x_{n}\right] \rightarrow A^{*}(F l(E))$ whose kernel is the ideal $I$ generated by the polynomials $e_{i}\left(x_{1}, \ldots, x_{n}\right)-c_{i}(E)$. As noted in 14.25 , this means that the divided difference operators descend to well-defined operators on $A\left[x_{1}, \ldots, x_{n}\right] / I$, and hence also on $A^{*}(F l(E))$. We'll also need:

Lemma 14.35. Let $w \in S_{n}, 1 \leq k<n$, and set $w^{\prime}=w \cdot s_{k}$. Then

$$
\partial_{k}\left(\left[\Omega_{w}\left(E_{\bullet}\right)\right]\right)= \begin{cases}{\left[\Omega_{w^{\prime}}\left(E_{\bullet}\right)\right]} & \text { if } w(k)>w(k+1) \\ 0 & \text { if } w(k)<w(k+1)\end{cases}
$$

Lemma 14.36. Let $E$ be a vector bundle of rank $d$ on a nonsingular variety $X$, and let $s: X \rightarrow E$ be a section. If the image of $s$ intersects the zero-section transversally, then $\left[s^{-1}(0)\right] \in H^{2 d} X$ is equal to $c_{d}(E)$.

In words, the top Chern class of a complex vector bundle is Poincaré dual to the zero set of a generic section. This theorem is well known to algebraic geometers working in the area but a concise statement is hard to find. One reference containing a (somewhat) accesible proof is: W. Fulton, Equivariant Intersection Theory, Notes by Dave Anderson, Michigan University, 2005-2006, available at:

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http://www.math.lsa.umich.edu/~ danderson/notes.html
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Proof. (of Theorem 14.34) Our first goal is to establish the formula for the longest permutation $w_{0}$. The essential set for this permutation is the collection of pairs $\{(p, n-p): 1 \leq p<n\}$, so the degeneracy locus of $\Omega_{w_{0}}\left(E_{\bullet}\right)$ is determined by the condition that $\operatorname{rank}\left(\rho^{*} E_{p} \rightarrow Q_{n-p}\right) \leq r_{w_{0}}(p, n-p)=0$ for $1 \leq p<n$, i.e. all these maps have rank 0 .

Let $H=\bigoplus_{p=1}^{n-1} \operatorname{Hom}\left(\rho^{*} E_{p}, Q_{n-p}\right)$ and $H^{\prime}=\bigoplus_{p=1}^{n-2} \operatorname{Hom}\left(\rho^{*} E_{p}, Q_{n-p-1}\right)$. Over each element $W_{\bullet} \in F l(E)$, the fiber $\left.H\right|_{W_{\bullet}}$ is a direct sum of complex vector spaces $\bigoplus_{p=1}^{n-1} \operatorname{Hom}\left(\left.\rho^{*} E_{p}\right|_{W_{\bullet}},\left.Q_{n-p}\right|_{W_{\bullet}}\right)$. An element of the fiber can be written as an $(n-1)$-tuple $\left(\alpha_{1}, \ldots \alpha_{n-1}\right)$, with each $\alpha_{p}$ a $\mathbb{C}$-linear map from $\left.\rho^{*} E_{p}\right|_{W_{\bullet}}$ to $\left.Q_{n-p}\right|_{W_{\bullet}}$. Let $K_{W_{\bullet}}$ be the collection of $(n-1)$-tuples making the following diagram commute for each $p$ :

I claim that $K=\bigcup_{W_{\bullet} \in F l(E)} K_{W_{\bullet}}$ is a subbundle of $H$. In fact $K$ is the kernel of a bundle map $g: H \rightarrow H^{\prime}$ carrying $\left(\alpha_{1}, \ldots, \alpha_{n-1}\right)$ to $\left(\beta_{1}, \ldots, \beta_{n-2}\right)$ where, for each $p, \beta_{p}=\alpha_{p+1} \circ \iota_{p}-j_{n-p} \circ \alpha_{p}$. One can check that this map is surjective, so we have an exact sequence of bundles

$$
0 \rightarrow K \hookrightarrow H \xrightarrow{g} H^{\prime} \rightarrow 0
$$

Since $\operatorname{Hom}\left(\left.\rho^{*} E_{p}\right|_{W_{\bullet}},\left.Q_{n-p}\right|_{W_{\bullet}}\right)$ has rank $p(n-p), H$ has rank $\sum_{p=1}^{n-1} p(n-p)$. Similarly, $H^{\prime}$ has rank $\sum_{p=1}^{n-2} p(n-p-1)$, and by dimension considerations it follows that $K$ has rank $N=n(n-1) / 2$. Define a section $s: F l(E) \rightarrow K, W_{\bullet} \mapsto\left(\alpha_{1}, \ldots, \alpha_{n-1}\right)$, where $\alpha_{p}$ is the linear map $\left.\left(\rho^{*} E_{p} \rightarrow Q_{n-p}\right)\right|_{W_{\bullet}}$. The zero locus is an old friend: $s^{-1}(0)=\left\{W_{\bullet} \in F l(E):\left.\left.\rho^{*} E_{p}\right|_{W_{\bullet}} \rightarrow Q_{n-p}\right|_{W_{\bullet}} \equiv 0\right\}=\Omega_{w_{0}}$. Therefore by lemma 14.36, we can relate the class of the degeneracy locus to the Chern classes of $K$ :

$$
\left[\Omega_{w_{0}}\left(E_{\bullet}\right)\right]=c_{N}(K)
$$

Our goal now is to compute $c_{N}(K)$. Recall that since we have an exact sequence of bundles

$$
0 \rightarrow K \rightarrow H \rightarrow H^{\prime} \rightarrow 0
$$

, by the Whitney formula $c(H)=c\left(H^{\prime}\right) c(K)$. In particular, the top Chern class of $H$ is the product of the top Chern class of $H^{\prime}$ with $c_{N}(K)$. Noting the similarity between the bundles $H$ and $H^{\prime}$, we will study the bundle $\operatorname{Hom}\left(\rho^{*} E_{k}, Q_{\ell}\right)$. By some basic homological algebra,

$$
\operatorname{Hom}\left(\rho^{*} E_{k}, Q_{\ell}\right)=\left(\rho^{*} E_{k}\right)^{*} \otimes Q_{\ell}
$$

where $\left(\rho^{*} E_{k}\right)^{*}$ denotes the dual bundle $\operatorname{Hom}\left(\rho^{*} E_{k}, \mathbb{C}\right)$. One final fact we must cite is that for a complex bundle $\omega, c_{i}\left(\omega^{*}\right)=(-1)^{i} c_{i}(\omega)$.

Since we have a filtration

$$
Q_{\ell} \rightarrow Q_{\ell-1} \rightarrow \ldots \rightarrow Q_{1}
$$

and for all $i$ the sequence

$$
0 \rightarrow \operatorname{ker}\left(Q_{i} \rightarrow Q_{i-1}\right) \rightarrow Q_{i} \rightarrow Q_{i-1} \rightarrow 0
$$

is exact, we can decompose $Q_{\ell}$ as a sum of line bundles:

$$
Q_{\ell}=\bigoplus_{i=1}^{\ell} \operatorname{ker}\left(Q_{i} \rightarrow Q_{i-1}\right)
$$

Similarly, by the filtration

$$
\rho^{*} E_{1} \rightarrow \rho^{*} E_{2} \rightarrow \ldots \rightarrow \rho^{*} E_{k},
$$

for all $j \leq k$ we get a short exact sequence

$$
0 \rightarrow \rho^{*} E_{j-1} \rightarrow \rho^{*} E_{j} \rightarrow \rho^{*} E_{j} / \rho^{*} E_{j-1} \rightarrow 0
$$

Hence we can decompose $\rho^{*} E_{k}$ as a sum of line bundles

$$
\rho^{*} E_{k}=\bigoplus_{j=1}^{k} \rho^{*} E_{j} / \rho^{*} E_{j-1} .
$$

Thus (since the Hom functor commutes with finite direct sums)

$$
\left(\rho^{*} E_{k}\right)^{*}=\bigoplus_{j=1}^{k}\left(\rho^{*} E_{j} / \rho^{*} E_{j-1}\right)^{*}
$$

This means we have reduced our study to

$$
\operatorname{Hom}\left(\rho^{*} E_{k}, Q_{\ell}\right)=\underset{\substack{1 \leq i \leq k \\ 1 \leq j \leq \ell}}{\bigoplus}\left(\left(\rho^{*} E_{j} / \rho^{*} E_{j-1}\right)^{*} \otimes \operatorname{ker}\left(Q_{i} \rightarrow Q_{i-1}\right)\right)
$$

Since each factor of the above sum is a tensor product of line bundles, Lemma 14.33 implies that

$$
c\left(\rho^{*} E_{j} / \rho^{*} E_{j-1} \otimes \operatorname{ker}\left(Q_{i} \rightarrow Q_{i-1}\right)\right)=1+x_{i}-y_{j} .
$$

Thus by the Whitney product theorem,

$$
c\left(\operatorname{Hom}\left(\rho^{*} E_{k}, Q_{\ell}\right)\right)=\prod_{\substack{1 \leq i \leq k \\ 1 \leq j \leq \ell}}\left(1+x_{i}-y_{j}\right)
$$

Since the top Chern class of a product bundle is the product of the top Chern classes of each of its factors, the top Chern class of $\operatorname{Hom}\left(\rho^{*} E_{k}, Q_{\ell}\right)$ is

$$
\prod_{\substack{1 \leq i \leq k \\ 1 \leq j \leq \ell}}\left(x_{i}-y_{j}\right) .
$$

In particular, this means that the top Chern class of $H$ is the product

$$
\prod_{p=1}^{n-1} \prod_{\substack{1 \leq i \leq p \\ 1 \leq j \leq n-p}}\left(x_{i}-y_{j}\right)
$$

and the top Chern class of $H^{\prime}$ is

$$
\prod_{p=1}^{n-2} \prod_{\substack{1 \leq i \leq p \\ 1 \leq j \leq n-p-1}}\left(x_{i}-y_{j}\right)
$$

Since $c_{N}(K)$ is the quotient of the top Chern classes of $H$ and $H^{\prime}$,

$$
\begin{aligned}
c_{N}(K) & =\prod_{1 \leq i \leq n-1}\left(x_{i}-y_{1}\right) \prod_{p=1}^{n-1} \prod_{1 \leq i \leq p}\left(x_{i}-y_{n-p}\right) \\
& =\prod_{i+j \leq n}\left(x_{i}-y_{j}\right) \\
& =\mathfrak{S}_{w_{0}}(x, y) .
\end{aligned}
$$

We have established the following identities:

$$
\left[\Omega_{w_{0}}\left(E_{\bullet}\right)\right]=\prod_{i+j \leq n}\left(x_{i}-y_{j}\right)=\mathfrak{S}_{w_{0}}(x, y)
$$

Now, for any $w \in S_{n}$, write $w=w_{0} \cdot s_{k_{1}} \cdots \cdot s_{k_{r}}$, where $r=n(n-1) / 2-\ell(w)$. Since $s_{k}(k)>s_{k}(k+1)$ for any $k$, applying Lemma $14.35 r$ times gives

$$
\begin{aligned}
{\left[\Omega_{w}\left(E_{\bullet}\right)\right] } & =\partial_{k_{r}} \circ \cdots \circ \partial_{k_{1}}\left(\left[\Omega_{w_{0}}\left(E_{\bullet}\right)\right]\right) \\
& =\partial_{k_{r}} \circ \cdots \circ \partial_{k_{1}}\left(\mathfrak{S}_{w_{0}}(x, y)\right) \\
& =\mathfrak{S}_{w}(x, y) .
\end{aligned}
$$

From this theorem Fulton deduces a more general statement that generalizes the above in a number of ways:
(1) Instead of applying to a flag of subbundles inside a single bundle, it applies in the setting of a morphism of filtered bundles $A_{1} \subset \cdots \subset A_{s} \xrightarrow{h} B_{t} \rightarrow \cdots \rightarrow B_{1}$.
(2) Instead of requiring a complete flag, it allows for partial flags.
(3) It permits $X$ to be singular.

The proof (and even the full statement) of the generalization are out of reach; they require additional tools from algebraic geometry and the theory of characteristic classes. (The idea of the proof is to set $E=A_{s} \oplus B_{t}$, embed $A_{s}$ in $E$ as the graph of $h$, apply the above theorem, and pull back the resulting formula for a cohomology class in $H^{*}(F l(E))$ to one in $H^{*}(X)$.) Let's content ourselves with a weaker statement, still too hard to prove but strong enough to work an interesting example:
Corollary 14.37. Given a generic morphism of filtered complex vector bundles $A_{1} \subset \cdots \subset A_{s} \xrightarrow{h} B_{t} \rightarrow$ $\cdots \rightarrow B_{1}$ on a smooth variety $X$, the class of the degeneracy locus in $A^{*}(X)$ is given by

$$
\left[\Omega_{w}(h)\right]=\mathfrak{S}_{w}\left(x_{1}, \ldots, x_{n} \cdot y_{1}, \ldots, y_{n}\right)
$$

where $x_{i}=c_{1}\left(\operatorname{ker}\left(B_{i} \rightarrow B_{i-1}\right)\right)$ and $y_{i}=c_{1}\left(A_{i} / A_{i-1}\right)$.

Example 14.38. Let us work out one explicit computation of the class of a degeneracy locus in the Chow ring. Suppose $h: E \rightarrow F$ is a generic map of 4-plane bundles on a variety, and $E_{1} \subset E_{2} \subset \cdots \subset$ $E_{4}=E$ and $F=F_{4} \rightarrow F_{3} \rightarrow \cdots \rightarrow F_{1}$ are complete flags of subbundles and quotient bundles. Let $w=[2431] \in S_{4}$. The the essential set is $\{(3,1),(2,3)\}$, and corresponding rank matrix is:

$$
\left[\begin{array}{llll}
0 & 1 & 1 & 1 \\
0 & 1 & \mathbf{1} & 2 \\
\mathbf{0} & 1 & 2 & 3 \\
1 & 2 & 3 & 4
\end{array}\right] .
$$

Therefore $\Omega_{w}$ consists of those points of $x$ over which $E_{1} \rightarrow F_{3}$ is zero and $E_{3} \rightarrow F_{2}$ has rank $\leq 1$. Since $w=w_{0} \cdot s_{2} \cdot s_{1}$,

$$
\begin{aligned}
\mathfrak{S}_{w}(x, y) & =\partial_{1} \circ \partial_{2}\left(\left(x_{1}-y_{1}\right)\left(x_{1}-y_{2}\right)\left(x_{1}-y_{3}\right)\left(x_{2}-y_{1}\right)\left(x_{2}-y_{2}\right)\left(x_{3}-y_{1}\right)\right) \\
& =\partial_{1}\left(\left(x_{1}-y_{1}\right)\left(x_{1}-y_{2}\right)\left(x_{1}-y_{3}\right)\left(x_{2}-y_{1}\right)\left(x_{3}-y_{1}\right)\right) \\
& =\left(x_{1}-y_{1}\right)\left(x_{2}-y_{1}\right)\left(x_{3}-y_{1}\right)\left(x_{1}+x_{2}-y_{2}-y_{3}\right) .
\end{aligned}
$$

Setting $x_{i}=c_{1}\left(\operatorname{ker}\left(F_{i} \rightarrow F_{i-1}\right)\right)$ and $y_{i}=c_{1}\left(E_{i} / E_{i-1}\right)$, our formula yields:

$$
\left[\Omega_{w}(h)\right]=\left(x_{1}-y_{1}\right)\left(x_{2}-y_{1}\right)\left(x_{3}-y_{1}\right)\left(x_{1}+x_{2}-y_{2}-y_{3}\right) .
$$

