# EULER CHARACTERISTICS OF HILBERT SCHEMES OF POINTS ON SIMPLE SURFACE SINGULARITIES 

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#### Abstract

We study the geometry and topology of Hilbert schemes of points on the orbifold surface $\left[\mathbb{C}^{2} / G\right]$, respectively the singular quotient surface $\mathbb{C}^{2} / G$, where $G<\operatorname{SL}(2, \mathbb{C})$ is a finite subgroup of type $A$ or $D$. We give a decomposition of the (equivariant) Hilbert scheme of the orbifold into affine space strata indexed by a certain combinatorial set, the set of Young walls. The generating series of Euler characteristics of Hilbert schemes of points of the singular surface of type $A$ or $D$ is computed in terms of an explicit formula involving a specialized character of the basic representation of the corresponding affine Lie algebra; we conjecture that the same result holds also in type $E$. Our results are consistent with known results in type $A$, and are new for type $D$.


## Contents

1. Orbifold singularities and their Hilbert schemes
1.1. Quotient surface singularities, Hilbert schemes and generating series
1.2. Simple surface singularities
1.3. Some terminology and structure of the paper

Acknowledgements
2. Type $A_{n}$
2.1. Type $A$ basics
2.2. Partitions, torus-fixed points and decompositions
2.3. Abacus of type $A_{n}$
2.4. Relating partitions to 0 -generated partitions
3. Type $D_{n}$ : ideals and Young walls
3.1. The binary dihedral group 10
3.2. Young wall pattern and Young walls
3.3. Decomposition of $\mathbb{C}[x, y]$ and the transformed Young wall pattern
3.4. Subspaces and operators

11
$\square$
3.5. Cell decompositions of equivariant Grassmannians 13
3.6. The Young wall associated to a homogeneous ideal
4. Type $D_{n}$ : decomposition of the orbifold Hilbert scheme
4.1. The decomposition 16
4.2. Incidence varieties 16
4. Prof of
4.4. Preparation for the proof of the incidence propositions
4.5. Profs of propositions about incidence vas
5.5. Proofs of propositions about incidence vareties
5. Type $D_{n}$ : special loci
5.1. Support blocks
5.2. Special loci in orbifold strata and the supporting rules

31
32
5.3. Special loci in Grassmannians
5.4. Proof of Theorem 5.2

35
6. Type $D_{n}$ : decomposition of the coarse Hilbert scheme
6.1. Distinguished 0-generated Young walls
6.2. The decomposition of the coarse Hilbert scheme 39
6.3. Possibly and almost invariant ideals
6.4. Euler characteristics of strata and the coarse generating series

[^0]7. Type $D_{n}$ : abacus combinatorics
7.1. Young walls and abacus of type $D_{n}$ 46
7.2. Core Young walls and their abacus representation
7.3. 0-generated Young walls and their abacus representations
7.4. The generating series of distinguished 0 -generated walls
Appendix A. Background on representation theory
A.1. Affine Lie algebras and extended basic representations
A.2. Affine crystals
A.3. Affine Lie algebras and Hilbert schemes
Appendix B. Joins
References

## 1. Orbifold singularities and their Hilbert schemes

1.1. Quotient surface singularities, Hilbert schemes and generating series. Let $G<$ $\mathrm{GL}(2, \mathbb{C})$ be a finite subgroup and denote by $\mathbb{C}^{2} / G$ the corresponding quotient variety. There are two different types of Hilbert scheme attached to this data. First, there is the classical Hilbert scheme $\operatorname{Hilb}\left(\mathbb{C}^{2} / G\right)$ of the quotient space. This is the moduli space of ideals in $\mathcal{O}_{\mathbb{C}^{2} / G}\left(\mathbb{C}^{2} / G\right)=$ $\mathbb{C}[x, y]^{G}$ of finite colength. We call this the coarse Hilbert scheme of points. It decomposes as

$$
\operatorname{Hilb}\left(\mathbb{C}^{2} / G\right)=\bigsqcup_{m \in \mathbb{N}} \operatorname{Hilb}^{m}\left(\mathbb{C}^{2} / G\right)
$$

into components that are quasiprojective but singular varieties indexed by "the number of points", the codimension $m$ of the ideal. Second, there is the moduli space of $G$-invariant finite colength subschemes of $\mathbb{C}^{2}$, the invariant part of $\operatorname{Hilb}\left(\mathbb{C}^{2}\right)$ under the lifted action of $G$. This Hilbert scheme is also well known and is variously called the orbifold Hilbert scheme [43] or equivariant Hilbert scheme [18]. We denote it by $\operatorname{Hilb}\left(\left[\mathbb{C}^{2} / G\right]\right)$. This space also decomposes as

$$
\operatorname{Hilb}\left(\left[\mathbb{C}^{2} / G\right]\right)=\bigsqcup_{\rho \in \operatorname{Rep}(G)} \operatorname{Hilb}^{\rho}\left(\left[\mathbb{C}^{2} / G\right]\right)
$$

where

$$
\operatorname{Hilb}^{\rho}\left(\left[\mathbb{C}^{2} / G\right]\right)=\left\{I \in \operatorname{Hilb}\left(\mathbb{C}^{2}\right)^{G}: H^{0}\left(\mathcal{O}_{\mathbb{C}^{2}} / I\right) \simeq_{G} \rho\right\}
$$

for any finite-dimensional representation $\rho \in \operatorname{Rep}(G)$ of $G$; here $\operatorname{Hilb}\left(\mathbb{C}^{2}\right)^{G}$ is the set of $G$-invariant ideals of $\mathbb{C}[x, y]$, and $\simeq_{G}$ means $G$-equivariant isomorphism. Being components of fixed point sets of a finite group acting on smooth quasiprojective varieties, the orbifold Hilbert schemes themselves are smooth and quasiprojective (5].

There is a natural pushforward map between the two kinds of Hilbert scheme: each $J \in$ $\operatorname{Hilb}\left(\left[\mathbb{C}^{2} / G\right]\right)$ can be mapped to its $G$-invariant part, giving a morphism [4, 3.4]

$$
\begin{array}{cl}
p_{*}: \operatorname{Hilb}\left(\left[\mathbb{C}^{2} / G\right]\right) & \rightarrow \operatorname{Hilb}\left(\mathbb{C}^{2} / G\right) \\
J & \mapsto J^{G}=J \cap \mathbb{C}[x, y]^{G}
\end{array}
$$

called the quotient-scheme map. There is also a set-theoretic pullback map, which however does not preserve flatness in families, so it is not a morphism between the Hilbert schemes: the inclusion $i: \mathbb{C}[x, y]^{G} \subset \mathbb{C}[x, y]$ induces a pullback map on the ideals, and its image is contained in the set of $G$-equivariant ideals, leading to a map of sets

$$
\begin{array}{cc}
i^{*}: \quad \operatorname{Hilb}\left(\mathbb{C}^{2} / G\right)(\mathbb{C}) & \rightarrow \quad \operatorname{Hilb}\left(\left[\mathbb{C}^{2} / G\right]\right)(\mathbb{C}) \\
I & \mapsto \quad i^{*} I=\mathbb{C}[x, y] . I
\end{array}
$$

Since for $I \triangleleft \mathbb{C}[x, y]^{G}$, we clearly have $(\mathbb{C}[x, y] . I)^{G}=I$, the composite $p_{*} \circ i^{*}$ is the identity on the set of ideals of the invariant ring.

We collect the topological Euler characteristics of the two versions of the Hilbert scheme into two generating functions. Let $\rho_{0}, \ldots, \rho_{n} \in \operatorname{Rep}(G)$ denote the (isomorphism classes of) irreducible representations of $G$, with $\rho_{0}$ the trivial representation.

Definition 1.1. (a) The orbifold generating series of the orbifold $\left[\mathbb{C}^{2} / G\right]$ is

$$
Z_{\left[\mathbb{C}^{2} / G\right]}\left(q_{0}, \ldots, q_{n}\right)=\sum_{m_{0}, \ldots, m_{n}=0}^{\infty} \chi\left(\operatorname{Hilb}^{m_{0} \rho_{0}+\ldots+m_{n} \rho_{n}}\left(\left[\mathbb{C}^{2} / G\right]\right)\right) q_{0}^{m_{0}} \cdot \ldots \cdot q_{n}^{m_{n}}
$$

(b) The coarse generating series of the singularity $\mathbb{C}^{2} / G$ is

$$
Z_{\mathbb{C}^{2} / G}(q)=\sum_{m=0}^{\infty} \chi\left(\operatorname{Hilb}^{m}\left(\mathbb{C}^{2} / G\right)\right) q^{m}
$$

Remark 1.2. For a smooth variety $X$, the generating series

$$
Z_{X}(q)=\sum_{m=0}^{\infty} \chi\left(\operatorname{Hilb}^{m}(X)\right) q^{m}
$$

of the Euler characteristics of Hilbert schemes of points of $X$, as well as various refinements of this series, have been extensively studied. In particular, for a nonsingular curve $C$, we have MacDonald's result 32]

$$
Z_{C}(q)=(1-q)^{-\chi(C)}
$$

whereas for a nonsingular surface $S$ we have (a specialization of) Göttsche's formula [16]

$$
\begin{equation*}
Z_{S}(q)=\left(\prod_{m=1}^{\infty}\left(1-q^{m}\right)^{-1}\right)^{\chi(S)} \tag{1}
\end{equation*}
$$

There are also results for higher-dimensional varieties [6].
For singular varieties $X$, the series $Z_{X}(q)$ is much less studied. For a singular curve $C$ with a finite set $\left\{P_{1}, \ldots, P_{k}\right\}$ of planar singularities however, we have the beautiful conjecture of Oblomkov and Shende [39], proved by Maulik [33], which takes the form

$$
\begin{equation*}
Z_{C}(q)=(1-q)^{-\chi(C)} \prod_{j=1}^{k} Z^{\left(P_{i}, C\right)}(q) \tag{2}
\end{equation*}
$$

Here $Z^{\left(P_{i}, C\right)}(q)$ are highly nontrivial local terms that depend only on the embedded topological type of the link of the singularity $P_{i} \in C$.
1.2. Simple surface singularities. In this paper we are only concerned with finite subgroups $G<\mathrm{SL}(2, \mathbb{C})$. See [19] for some partial results for some other finite groups. As it is well known, finite subgroups of $\operatorname{SL}(2, \mathbb{C})$ are classified into three types: type $A_{n}$ for $n \geq 1$, type $D_{n}$ for $n \geq 4$ and type $E_{n}$ for $n=6,7,8$. The type of the singularity can be parametrized by a simply-laced irreducible Dynkin diagram with $n$ nodes, arising from an irreducible simply laced root system $\Delta$. We denote the corresponding group by $G_{\Delta}<\mathrm{SL}(2, \mathbb{C})$; all other data corresponding to the chosen type will also be labelled by the subscript $\Delta$. Irreducible representations $\rho_{0}, \ldots, \rho_{n}$ of $G_{\Delta}$ are then labelled by vertices of the affine Dynkin diagram associated with $\Delta$. The singularity $\mathbb{C}^{2} / G_{\Delta}$ is known as a simple (Kleinian, surface) singularity; we will refer to the corresponding orbifold $\left[\mathbb{C}^{2} / G_{\Delta}\right]$ as the simple singularity orbifold.

As we recall in Appendix A.3 the following result is known.
Theorem 1.3 ([37]). Let $\left[\mathbb{C}^{2} / G_{\Delta}\right]$ be a simple singularity orbifold. Then its orbifold generating series can be expressed as

$$
\begin{equation*}
Z_{\left[\mathbb{C}^{2} / G_{\Delta}\right]}\left(q_{0}, \ldots, q_{n}\right)=\left(\prod_{m=1}^{\infty}\left(1-q^{m}\right)^{-1}\right)^{n+1} \cdot \sum_{\bar{m}=\left(m_{1}, \ldots, m_{n}\right) \in \mathbb{Z}^{n}} q_{1}^{m_{1}} \cdots \cdots q_{n}^{m_{n}}\left(q^{1 / 2}\right)^{\bar{m}^{\top} \cdot C_{\Delta} \cdot \bar{m}} \tag{3}
\end{equation*}
$$

where $q=\prod_{i=0}^{n} q_{i}^{d_{i}}$ with $d_{i}=\operatorname{dim} \rho_{i}$, and $C_{\Delta}$ is the finite type Cartan matrix corresponding to $\Delta$.
Our first main result is a strengthening of this theorem. Given a Dynkin diagram $\Delta$ of type $A$ or $D$, we will recall below in [2.2 respectively 3.2 the definition of a certain combinatorial set, the set of Young walls $\mathcal{Z}_{\Delta}$ of type $\Delta$.

Theorem 1.4. Let $\left[\mathbb{C}^{2} / G_{\Delta}\right]$ be a simple singularity orbifold, where $\Delta$ is of type $A_{n}$ for $n \geq 1$ or $D_{n}$ for $n \geq 4$. Then there exists a decomposition

$$
\operatorname{Hilb}\left(\left[\mathbb{C}^{2} / G_{\Delta}\right]\right)=\bigsqcup_{Y \in \mathcal{Z}_{\Delta}} \operatorname{Hilb}\left(\left[\mathbb{C}^{2} / G_{\Delta}\right]\right)_{Y}
$$

into locally closed strata indexed by the set of Young walls $\mathcal{Z}_{\Delta}$ of the appropriate type. Each stratum is isomorphic to an affine space of a certain dimension, and in particular has Euler characteristic $\chi\left(\operatorname{Hilb}\left(\left[\mathbb{C}^{2} / G_{\Delta}\right]\right)_{Y}\right)=1$.

For type $A$, the set of Young walls is simply the set of finite partitions, represented as Young diagrams, equipped with a diagonal labelling. In this case, Theorem 1.4 is well known; the decomposition in type $A$ is not unique, but depends on a choice of a one-dimensional subtorus of the full torus $\left(\mathbb{C}^{*}\right)^{2}$ acting on the affine plane $\mathbb{C}^{2}$. For completeness, we summarize the details in 2.2 , On the other hand, the type $D$ case appears to be new; in this case, our decomposition is unique, there is no further choice to make.
Remark 1.5. The orbifold Hilbert schemes of points for $G<\mathrm{SL}(2, \mathbb{C})$ are well known to be Nakajima quiver varieties for the corresponding affine quiver. As it was shown in [40], certain Lagrangian subvarieties in Nakajima quiver varieties are isomorphic to quiver Grassmannians for the preprojective algebra of the same type, parametrizing submodules of certain fixed modules. On the other hand, results of the recent papers [30, 31] imply that every quiver Grassmannian of a representation of a quiver of affine type $D$ has a decomposition into affine spaces. The relation between this decomposition and ours deserves further investigation.

As we will explain combinatorially in 2.3, respectively 7.2, and via representation theory in A.1. A.2, the right hand side of (3) enumerates the set of Young walls $\mathcal{Z}_{\Delta}$ of the appropriate type. Thus Theorem 1.4 implies Theorem 1.3 ,
Remark 1.6. In type $A$, it is easy to refine formula (3) to a formula involving the Betti numbers [13], or the motives [18], of the orbifold Hilbert schemes. We leave the study of such a refinement in type $D$ to future work; compare Remark 4.4.

The second main result of our paper is the following formula, which says that the coarse generating series is a very particular specialization of the orbifold one.
Theorem 1.7. Let $\mathbb{C}^{2} / G_{\Delta}$ be a simple singularity, where $\Delta$ is of type $A_{n}$ for $n \geq 1$ or $D_{n}$ for $n \geq 4$. Let $h^{\vee}$ be the (dual) Coxeter number of the corresponding finite root system (one less than the dimension of the corresponding simple Lie algebra divided by $n$ ). Then

$$
Z_{\mathbb{C}^{2} / G_{\Delta}}(q)=\left(\prod_{m=1}^{\infty}\left(1-q^{m}\right)^{-1}\right)^{n+1} \cdot \sum_{\bar{m}=\left(m_{1}, \ldots, m_{n}\right) \in \mathbb{Z}^{n}} \zeta^{m_{1}+m_{2}+\cdots+m_{n}}\left(q^{1 / 2}\right)^{m^{\top} \cdot C_{\Delta} \cdot \bar{m}},
$$

where $\zeta=\exp \left(\frac{2 \pi i}{1+h^{v}}\right)$ and $C_{\Delta}$ is the finite type Cartan matrix corresponding to $\Delta$.
Thus $Z_{\mathbb{C}^{2} / G_{\Delta}}(q)$ is obtained from $Z_{\left[\mathbb{C}^{2} / G_{\Delta}\right]}\left(q_{0}, \ldots, q_{n}\right)$ by the substitutions

$$
q_{1}=\cdots=q_{n}=\exp \left(\frac{2 \pi i}{1+h^{\vee}}\right), \quad q_{0}=q \exp \left(-\frac{2 \pi i}{1+h^{\vee}} \sum_{i \neq 0} \operatorname{dim} \rho_{i}\right)
$$

In type $A$, the formula in Theorem 1.7 is not new: it was proved directly (in a slight disguise) by Dijkgraaf and Sulkowski in [8] and also recently, using completely different methods, by Toda in [42]. Our main contribution is the general Lie-theoretic formulation, as well as a proof in type $D$; we also provide a direct combinatorial proof in type $A$, which appears to be new.

One can check directly that the generating series in Theorem 1.7 has also integer coefficients for $E_{6}, E_{7}$ and $E_{8}$ to a high power in $q$. This motivates the following.

Conjecture 1.8. Let $\mathbb{C}^{2} / G_{\Delta}$ be a simple singularity of type $E_{n}$ for $n=6,7,8$. Let $h^{\vee}$ be the (dual) Coxeter number of the corresponding finite root system. Then, as for other types,

$$
Z_{\mathbb{C}^{2} / G_{\Delta}}(q)=\left(\prod_{m=1}^{\infty}\left(1-q^{m}\right)^{-1}\right)^{n+1} \cdot \sum_{\bar{m}=\left(m_{1}, \ldots, m_{n}\right) \in \mathbb{Z}^{n}} \zeta^{m_{1}+m_{2}+\cdots+m_{n}}\left(q^{1 / 2}\right)^{\bar{m}^{\top} \cdot C_{\Delta} \cdot \bar{m}}
$$

where $\zeta=\exp \left(\frac{2 \pi i}{1+h^{\vee}}\right)$ and $C_{\Delta}$ is the finite type Cartan matrix corresponding to $\Delta$.
The key tool in our proof of Theorem 1.7 for types $A$ and $D$ is the combinatorics of Young walls, in particular their abacus representation. We are not aware of such explicit combinatorics in type $E$. We hope to return to this question in later work.

Remark 1.9. We are dealing here with Hilbert schemes, parametrizing rank $r=1$ sheaves on the orbifold or singular surface. In the relationship between the instantons on algebraic surfaces and affine Lie algebras, level equals rank [17]. Indeed the (extended) basic represenatation underlying the Young wall combinatorics (see Appendix) has level $l=1$. Thus the substitution above is by the root of unity $\zeta=\exp \left(\frac{2 \pi i}{l+h^{\vee}}\right)$, with $l=1$ and $h^{\vee}$ the (dual) Coxeter number. There is an intriguing analogy here with the Verlinde formula, which uses a similar substitution, into characters of Lie algebras, by a root of unity $\zeta=\exp \left(\frac{2 \pi i}{l+h^{\vee}}\right)$, where $l$ again is the level, and $h^{\vee}$ the (dual) Coxeter number of the root system of the Lie algebra of the gauge group. The geometric significance of this observation, if any, is left for future research.

Remark 1.10. Given the results above, it is easy to write down a global formula analogous to (2) for a singular surface with canonical singularities. This formula, as well as its modularity, are discussed in the announcement [21].
1.3. Some terminology and structure of the paper. We work over the field $\mathbb{C}$ of complex numbers. We call a regular map $f: X \rightarrow Y$ a trivial affine fibration with fibre $\mathbb{A}^{k}$, if there is an isomorphism $X \cong Y \times \mathbb{A}^{k}$ with $f$ being the first projection.

The structure of the rest of the paper is as follows. In Section 2, we give a new proof of Theorem [1.7 in type $A$, which has the advantage that it generelizes away from that case. The rest of the paper treats the case of type $D$. In Section 3, we introduce Schubert-style cell decompositions of Grassmannians of homogeneous summands of $\mathbb{C}[x, y]$. In Section 4 we give a cell decomposition of the orbifold Hilbert scheme, proving Theorem [1.4] In Section 5, we discuss some special subsets of the strata and their geometry. A decomposition of the coarse Hilbert scheme is given in Section 6. In Section 7 the proof of Theorem [1.7 is completed using combinatorial enumeration. The representation theoretic background is briefly summarized in Appendix A. Some relevant facts on joins of projective varieties are discussed in Appendix B.

Acknowledgements. The authors would like to thank Gwyn Bellamy, Alastair Craw, Eugene Gorsky, Ian Grojnowski, Kevin McGerty, Iain Gordon, Tomas Nevins and Tamás Szamuely for helpful comments and discussions. Á.Gy. was partially supported by the Lendület program (Momentum Programme) of the Hungarian Academy of Sciences and by ERC Advanced Grant LDTBud (awarded to András Stipsicz). A.N. was partially supported by OTKA Grants 100796 and K112735. B.Sz. was partially supported by EPSRC Programme Grant EP/I033343/1.

## 2. Type $A_{n}$

2.1. Type $A$ basics. Let $\Delta$ be the root system of type $A_{n}$. Choosing a primitive ( $n+1$ )-st root of unity $\omega$, the corresponding subgroup $G_{\Delta}$ of $S L(2, \mathbb{C})$, a cyclic subgroup of order $n+1$, is generated by the matrix

$$
\sigma=\left(\begin{array}{cc}
\omega & 0 \\
0 & \omega^{-1}
\end{array}\right)
$$

All irreducible representations of $G_{\Delta}$ are one dimensional, and they are simply given by $\rho_{j}: \sigma \mapsto \omega^{j}$, for $j \in\{0, \ldots, n\}$. The corresponding McKay quiver is the cyclic Dynkin diagram of type $\widetilde{A}_{n}^{(1)}$.

The group $G_{\Delta}$ acts on $\mathbb{C}^{2} ;$ the quotient variety $\mathbb{C}^{2} / G_{\Delta}$ has an $A_{n}$ singularity at the origin. The matrix $\sigma$ clearly commutes with the diagonal two-torus $T=\left(\mathbb{C}^{*}\right)^{2}$, and so $T$ acts on the quotient $\mathbb{C}^{2} / G_{\Delta}$ and the orbifold $\left[\mathbb{C}^{2} / G_{\Delta}\right]$. Consequently $T$ also acts on the orbifold Hilbert scheme $\operatorname{Hilb}\left(\left[\mathbb{C}^{2} / G_{\Delta}\right]\right)$ and the (reduced) coarse Hilbert scheme $\operatorname{Hilb}\left(\mathbb{C}^{2} / G_{\Delta}\right)$ as well.
2.2. Partitions, torus-fixed points and decompositions. Consider the set $\mathbb{N} \times \mathbb{N}$ of pairs of non-negative integers; we will draw this set as a set of blocks on the plane, occupying the nonnegative quadrant. Label blocks diagonally with $(n+1)$ labels $0, \ldots, n$ as in the picture; the block with coordinates $(i, j)$ is labelled with $(i-j) \bmod (n+1)$. We will call this the pattern of type $A_{n}$.

| $\vdots$ |  |  |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $\vdots$ |  |  |  |  |  |  |
| 0 | 1 |  |  |  |  |  |
| 1 | 2 |  |  |  |  |  |
|  | $\vdots$ | $\vdots$ |  |  |  |  |
| $n$ | 0 |  | $n-2$ | $n-1$ | $n$ | 0 |
| 0 | 1 | $\cdots$ | $n-1$ | $n$ | 0 | 1 |$\ldots$

Let $\mathcal{P}$ denote the set of partitions. Given a partition $\lambda=\left(\lambda_{1}, \ldots, \lambda_{k}\right) \in \mathcal{P}$, with $\lambda_{1} \geq \ldots \geq \lambda_{k}$ positive integers, we consider its Young (or Ferrers) diagram, the subset of $\mathbb{N} \times \mathbb{N}$ which consists of the $\lambda_{i}$ lowest blocks in column $i-1$. The blocks in $\lambda$ also get labelled by the $n+1$ labels. Let $\mathcal{Z}_{\Delta}$ denote the resulting set of $(n+1)$-labelled partitions, including the empty partition. For a labelled partition $\lambda \in \mathcal{Z}_{\Delta}$, let $\mathrm{wt}_{j}(\lambda)$ denote the number of blocks in $\lambda$ labelled $j$, and define the multiweight of $\lambda$ to be $\underline{\mathrm{wt}}(\lambda)=\left(\mathrm{wt}_{0}(\lambda), \ldots, \mathrm{wt}_{n}(\lambda)\right)$.

Proposition 2.1. The torus $T$ acts with isolated fixed points on $\operatorname{Hilb}\left(\left[\mathbb{C}^{2} / G_{\Delta}\right]\right)$, parametrized by the set $\mathcal{Z}_{\Delta}$ of $(n+1)$-labelled partitions. More precisely, for $k_{0}, \ldots, k_{n}$ non-negative integers and $\rho=\oplus_{j=0}^{n} \rho_{i}^{\oplus k_{i}}$, the $T$-fixed points on $\operatorname{Hilb}^{\rho}\left(\left[\mathbb{C}^{2} / G_{\Delta}\right]\right)$ are parametrized by $(n+1)$-labelled partitions of multiweight $\left(k_{0}, \ldots, k_{n}\right)$.
Proof. We just sketch the proof, which is well known 10, 13]. It is clear that the $T$-fixed points on $\operatorname{Hilb}\left(\left[\mathbb{C}^{2} / G_{\Delta}\right]\right)$, which coincide with the $T$-fixed points on $\operatorname{Hilb}\left(\mathbb{C}^{2}\right)$, are the monomial ideals in $\mathbb{C}[x, y]$ of finite colength. The monomial ideals are enumerated in turn by Young diagrams of partitions. The labelling of each block gives the weight of the $G_{\Delta}$-action on the corresponding monomial, proving the refined statement.

Corollary 2.2. There exist a locally closed decomposition, depending on a choice specified below, of $\operatorname{Hilb}\left(\left[\mathbb{C}^{2} / G_{\Delta}\right]\right)$ into strata indexed by the set of $(n+1)$-labelled partitions. Each stratum is isomorphic to an affine space.

Proof. Again, this is well known 38]. Fixing a representation $\rho$, choose a sufficiently general onedimensional subtorus $T_{0} \subset T$ which has positive weight on both $x$ and $y$. For general $T_{0} \subset T$, the fixed point set on $\operatorname{Hilb}^{\rho}\left(\left[\mathbb{C}^{2} / G_{\Delta}\right]\right)$ is unchanged and in particular consists of a finite number of isolated points. Choosing positive weights on $x, y$ ensures that all limits of $T_{0}$-orbits at $t=0$ in $\operatorname{Hilb}\left(\left[\mathbb{C}^{2} / G_{\Delta}\right]\right)$ exist, even though $\operatorname{Hilb}^{\rho}\left(\left[\mathbb{C}^{2} / G_{\Delta}\right]\right)$ is non-compact. Since $\operatorname{Hilb}^{\rho}\left(\left[\mathbb{C}^{2} / G_{\Delta}\right]\right)$ is smooth, the result follows by taking the Białynicki-Birula decomposition of $\operatorname{Hilb}^{\rho}\left(\left[\mathbb{C}^{2} / G_{\Delta}\right]\right)$ given by the $T_{0}$-action.

Denote by

$$
Z_{\Delta}\left(q_{0}, \ldots, q_{n}\right)=\sum_{\lambda \in \mathcal{Z}_{\Delta}} \underline{q}^{\underline{\operatorname{wt}}(\lambda)}
$$

the generating series of $(n+1)$-labelled partitions, where we used multi-index notation

$$
\underline{q}^{\underline{\mathrm{wt}}(\lambda)}=\prod_{i=0}^{n} q_{i}^{\mathrm{wt}_{i}(\lambda)}
$$

From either of the previous two statements, we immediately deduce the following.
Corollary 2.3. Let $\left[\mathbb{C}^{2} / G_{\Delta}\right]$ be a simple singularity orbifold of type $A$. Then its orbifold generating series can be expressed as

$$
\begin{equation*}
Z_{\left[\mathbb{C}^{2} / G_{\Delta}\right]}\left(q_{0}, \ldots, q_{n}\right)=Z_{\Delta}\left(q_{0}, \ldots, q_{n}\right) \tag{4}
\end{equation*}
$$

According to [13] the generating series of $(n+1)$-labelled partitions has the following form:

$$
\begin{equation*}
Z_{\Delta}\left(q_{0}, \ldots, q_{n}\right)=\frac{\sum_{\bar{m}=\left(m_{1}, \ldots, m_{n}\right) \in \mathbb{Z}^{k}}^{\infty} q_{1}^{m_{1}} \cdots \cdots q_{n}^{m_{n}}\left(q^{1 / 2}\right)^{\bar{m}^{\top}} \cdot C \cdot \bar{m}}{\prod_{m=1}^{\infty}\left(1-q^{m}\right)^{n+1}} \tag{5}
\end{equation*}
$$

where $q=q_{0} \cdots q_{n}$ and $C$ is the (finite) Cartan matrix of type $A_{n}$; for a sketch proof, see the end of 2.3 below. In particular, (4) and (5) imply Theorem 1.3 for type $A$.

We now turn to the coarse Hilbert scheme. Let us define a subset $\mathcal{Z}_{\Delta}^{0}$ of the set of $(n+1)$ labelled partitions $\mathcal{Z}_{\Delta}$ as follows. An $(n+1)$-labelled partition $\lambda \in \mathcal{Z}_{\Delta}$ will be called 0 -generated (a slight misnomer, this should be really be "complement-0-generated") if the complement of $\lambda$ inside $\mathbb{N} \times \mathbb{N}$ can be completely covered by translates of $\mathbb{N} \times \mathbb{N}$ to blocks labelled 0 contained in this complement. Equivalently, an $(n+1)$-labelled partition $\lambda$ is 0 -generated, if all its addable blocks (blocks whose addition gives another partition) are labelled 0 . It is immediately seen that this condition is equivalent to the corresponding monomial ideal $I \triangleleft \mathbb{C}[x, y]$ being generated by its invariant part $I \cap \mathbb{C}[x, y]^{G \Delta}$. Indeed, we have the following.
Proposition 2.4. The torus $T$ acts with isolated fixed points on $\operatorname{Hilb}\left(\mathbb{C}^{2} / G_{\Delta}\right)$, which are in bijection with the set $\mathcal{Z}_{\Delta}^{0}$ of 0 -generated $(n+1)$-labelled partitions. More precisely, for a nonnegative integer $k$, the $T$-fixed points on $\operatorname{Hilb}^{k}\left(\mathbb{C}^{2} / G_{\Delta}\right)$ are parametrized by 0 -generated $(n+1)$ labelled partitions $\lambda$ with 0 -weight $\mathrm{wt}_{0}(\lambda)=k$.
Proof. This is immediate from the above discussion. The $T$-fixed points of $\operatorname{Hilb}\left(\mathbb{C}^{2} / G_{\Delta}\right)$ are the monomial ideals $I$ of $\mathbb{C}[x, y]^{G_{\Delta}}$ of finite colength. Inside $\mathbb{C}[x, y]$, the ideals they generate correspond to partitions which are 0 -generated. The ring $\mathbb{C}[x, y]^{G_{\Delta}}$ has a basis consisting of monomials with corresponding blocks labelled 0 inside $\mathbb{C}[x, y]$; thus the codimension of a monomial ideal $I$ inside $\mathbb{C}[x, y]^{G_{\Delta}}$ is simply the number of blocks denoted 0 .

Denoting by

$$
Z_{\Delta}^{0}(q)=\sum_{\lambda \in \mathcal{Z}_{\Delta}^{0}} q^{\mathrm{wt}_{0}(\lambda)}
$$

the corresponding specialization of the generating series of 0 -generated $(n+1)$-labelled partitions, we deduce the following.
Corollary 2.5. Let $\left[\mathbb{C}^{2} / G_{\Delta}\right]$ be a simple singularity orbifold of type $A$. Then the coarse generating series can be expressed as

$$
\begin{equation*}
Z_{\mathbb{C}^{2} / G_{\Delta}}(q)=Z_{\Delta}^{0}(q) \tag{6}
\end{equation*}
$$

Proof of Theorem 1.7 for the $A_{n}$ case. The (dual) Coxeter number of the type $A_{n}$ root system is $h^{\vee}=n+1$. Thus Theorem 1.7 for this case follows from Corollary 2.5, formula (5), and the combinatorial Proposition 2.7 below, which computes the series $Z_{\Delta}^{0}(q)$.
Remark 2.6. The single variable generating series $Z_{\mathbb{C}^{2} / G_{\Delta}}$ in type $A$ was calculated by Toda in 42] using threefold machinery including a flop formula for Donaldson-Thomas invariants of certain Calabi-Yau threefolds. He does not mention any connection to Lie theory. The combinatorics, and the one-variable formula for $Z_{\Delta}^{0}(q)$, were already known to Dijkgraaf and Sulkowski [8]. They do not give the interpretation of the combinatorial formula in terms of Hilbert schemes, though they are clearly motivated by closely related ideas. Their proof is different, using the method of Andrews [2] in place of the abacus combinatorics we use below. We believe that already in type $A$, our new proof is preferable since it directly exhibits the clear connection between the orbifold and coarse generating series. Also, as we show later, this method generalizes away from type $A$.
2.3. Abacus of type $A_{n}$. We now introduce some standard combinatorics related to the type $A$ root system, which will allow us to relate the generating series $Z_{\Delta}$ of $(n+1)$-labelled partitions to the specialized series $Z_{\Delta}^{0}$ of 0 -generated partitions. We follow the notations of 29].

The abacus of type $A_{n}$ is the arrangement of the set of integers in $(n+1)$ columns according to the following pattern.

$$
\begin{array}{ccccc}
\vdots & \vdots & & \vdots & \vdots \\
-2 n-1 & -2 n & \ldots & -n-2 & -n-1 \\
-n & -n+1 & \ldots & -1 & 0 \\
1 & 2 & \ldots & n & n+1 \\
n+2 & n+3 & \ldots & 2 n+1 & 2 n+2 \\
\vdots & \vdots & & \vdots & \vdots
\end{array}
$$

Each integer in this pattern is called a position. For any integer $1 \leq k \leq n+1$ the set of positions in the $k$-th column of the abacus is called the $k$-th runner. An abacus configuration is a set of beads, denoted by $\bigcirc$, placed on the positions, with each position occupied by at most one bead.

To an $(n+1)$-labelled partition $\lambda=\left(\lambda_{1}, \ldots, \lambda_{k}\right) \in \mathcal{Z}_{\Delta}$ we associate its abacus representation (sometimes also called Maya diagram) as follows: place a bead in position $\lambda_{i}-i+1$ for all $i$, interpreting $\lambda_{i}$ as 0 for $i>k$. Alternatively, the abacus representation can be described by tracing the outer profile of the Young diagram of a partition: the occupied positions occur where the profile moves "down", whereas the empty positions are where the profile moves "right". In the abacus representation of a partition, the number of occupied positive positions is always equal to the number of absent nonpositive positions; we call such abacus configurations balanced. Conversely, it is easy to see that any balanced configuration represents a unique $(n+1)$-labelled partition, an element of $\mathcal{Z}_{\Delta}$.

For $n=0$, we obtain a representation of partitions on a single runner; this is sometimes called the Dirac sea representation of partitions.

The $(n+1)$-core of a labelled partition $\lambda \in \mathcal{Z}_{\Delta}$ is the partition obtained from $\lambda$ by successively removing border strips of length $n+1$, leaving a partition at each step, until this is no longer possible. Here a border strip is a skew Young diagram which does not contain $2 \times 2$ blocks and which contains exactly one $j$-labelled block for all labels $j$. The removal of a border strip corresponds in the abacus representation to shifting one of the beads up on its runner, if there is an empty space on the runner above it. In this way, the core of a partition corresponds to the bead configuration in which all the beads are shifted up as much as possible; this in particular shows that the ( $n+1$ )-core of a partition is well-defined. We denote by $\mathcal{C}_{\Delta}$ the set of $(n+1)$-core partitions, and

$$
c: \mathcal{Z}_{\Delta} \rightarrow \mathcal{C}_{\Delta}
$$

the map which takes an $(n+1)$-labelled partition to its $(n+1)$-core.
Given an $(n+1)$-core $\lambda$, we can read the $(n+1)$ runners of its abacus rrepresentation separately. These will not necessarily be balanced. The $i$-th one will be shifted from the balanced position by a certain integer number $a_{i}$ steps, which is negative if the shift is toward the negative positions (upwards), and positive otherwise. These numbers satisfy $\sum_{i=0}^{n} a_{i}=0$, since the original abacus configuration was balanced. The set $\left\{a_{1}, \ldots, a_{n}\right\}$ completely determines the partition, so we get a bijection

$$
\begin{equation*}
\mathcal{C}_{\Delta} \longleftrightarrow\left\{\sum_{i=0}^{n} a_{i}=0\right\} \subset \mathbb{Z}^{n+1} \tag{7}
\end{equation*}
$$

We will represent an $(n+1)$-core partition by the corresponding $(n+1)$-tuple $\underline{a}=\left(a_{0}, \ldots, a_{n}\right)$.
On the other hand, for an arbitrary partition, on each runner we have a partition up to shift, so we get a bijection

$$
\mathcal{Z}_{\Delta} \longleftrightarrow \mathcal{C}_{\Delta} \times \mathcal{P}^{n+1}
$$

This corresponds to the structure of formula (5) above; its denominator is the generating series of ( $n+1$ )-tuples of (unlabelled) partitions, whereas its numerator (after eliminating a variable) is exactly a sum over $\underline{a} \in \mathcal{C}_{\Delta}$. The multiweight of a core partition corresponding to an element $\underline{a}$ is given by the quadratic expression $Q(\underline{a})$ in the exponent of the numerator of (5). For more details, see Bijections 1-2 in [14, §2].
2.4. Relating partitions to $\mathbf{0}$-generated partitions. The purpose of this section is to prove the following, completely combinatorial statement.

Proposition 2.7. Let $\Delta$ be of type $A_{n}$, and let $\xi$ be a primitive $(n+2)$-nd root of unity. Then the generating series of 0 -generated partitions can be computed from that of all $(n+1)$-labelled ones by the following substitution:

$$
Z_{\Delta}^{0}(q)=\left.Z_{\Delta}\left(q_{0}, \ldots, q_{n}\right)\right|_{q_{0}=\xi^{-n} q, q_{1}=\cdots=q_{n}=\xi}
$$

We start by combinatorially relating partitions to 0 -generated partitions. $\mathcal{Z}_{\Delta}^{0}$ is clearly a subset of $\mathcal{Z}_{\Delta}$, but there is also a map

$$
p: \mathcal{Z}_{\Delta} \rightarrow \mathcal{Z}_{\Delta}^{0}
$$

defined as follows: for an arbitrary partition $\lambda$, let $p(\lambda)$ be the smallest 0 -generated partition containing it. Since the set of 0 -generated partitions is closed under intersection, $p(\lambda)$ is welldefined, and it can be constructed as follows: $p(\lambda)$ is the complement of the unions of the translates of $\mathbb{N} \times \mathbb{N}$ to 0-labelled blocks in the complement of $\lambda$. It is clear that $p(\lambda)$ can equivalently be obtained by adding all possible addable blocks to $\lambda$ of labels different from 0 .
Remark 2.8. The map $p$ can also be described in the language of ideals. If the monomial ideal $I \triangleleft \mathbb{C}[x, y]$ corresponds to the partition $\lambda$, then the monomial ideal $i^{*} p_{*} I=\left(I \cap \mathbb{C}[x, y]^{G \Delta}\right) . \mathbb{C}[x, y] \triangleleft$ $\mathbb{C}[x, y]$ corresponds to the partition $p(\lambda)$.
Lemma 2.9. The bead configurations corresponding to 0 -generated partitions are exactly those which have all beads right-justified on each row, with no empty position to the right of a filled position. The map $p: \mathcal{Z}_{\Delta} \rightarrow \mathcal{Z}_{\Delta}^{0}$ can be described in the abacus representation by the process of pushing all beads of an abacus configuration as far right as possible.

Proof. This follows from the description of the map from a partition to its abacus representation using the profile of the partition. Indeed, a 0 -generated partition has a profile which only turns from "down" to "right" at 0-labelled blocks. In other words, the only time when a string of filled positions can be followed by an empty position is when the last filled position is on the rightmost runner. In other words, there cannot be empty positions to the right of filled positions in a row. The proof of the second statement is similar.
Remark 2.10. As explained above, the maps $c: \mathcal{Z}_{\Delta} \rightarrow \mathcal{C}_{\Delta}$ and $p: \mathcal{Z}_{\Delta} \rightarrow \mathcal{Z}_{\Delta}^{0}$ have natural descriptions on abacus configurations: corresponds to pushing beads all the way up within their column, whereas $p$ corresponds to pushing beads all the way to the right within their row. It is then clear that there is also a third map $\mathcal{Z}_{\Delta} \rightarrow^{0} \mathcal{Z}_{\Delta} \subset \mathcal{Z}_{\Delta}$, dual to $p$, defined on the abacus by pushing beads all the way to the left. On labelled partitions this corresponds to the operation of removing all possible blocks with labels different from 0 . This dual constuction occured in the literature earlier in [15].

Proof of Proposition 2.7. We will prove the substitution formula on the fibres of the map $p: \mathcal{Z}_{\Delta} \rightarrow$ $\mathcal{Z}_{\Delta}^{0}$. In other words, we need to show that for any given $\lambda_{0} \in \mathcal{Z}_{\Delta}^{0}$, we have

$$
\begin{equation*}
\left.\sum_{\mu \in p^{-1}\left(\lambda_{0}\right)} \underline{q}^{\underline{\mathrm{wt}}(\mu)}\right|_{q_{1}=\cdots=q_{n}=\xi, q_{0}=\xi^{-n} q}=q^{\mathrm{wt}_{0}\left(\lambda_{0}\right)} \tag{8}
\end{equation*}
$$

As a first step, we reduce the computation to 0 -generated cores. Given an arbitrary 0 -generated partition $\lambda$, by the first part of Lemma 2.9 its core $\nu=c(\lambda)$ is also 0 -generated, and the corresponding abacus configuration can be obtained by permuting the rows of the configuration of $\lambda$. Fix one such permutation $\sigma$ of the rows. Then, using the second part of Lemma 2.9, we can use the row permutation $\sigma$ to define a bijection

$$
\tilde{\sigma}: p^{-1}(\lambda) \rightarrow p^{-1}(\nu)
$$

between (abacus representations of) partitions in the fibres, mapping $\lambda$ itself to $\nu$.
The difference between the partitions $\lambda$ and $\nu$ is a certain number of border strips, each removal represented by pushing up one bead on some runner by one step. Each border strip contains one block of each label, so the total number of times we need to push up a bead by one step on the different runners is $N=\mathrm{wt}_{0}(\lambda)-\mathrm{wt}_{0}(\nu)$. Thus, with $q=q_{0} \cdot \ldots \cdot q_{n}$ as in the substitution above, we can write

$$
\underline{q}^{\underline{\mathrm{wt}}(\lambda)}=q^{\mathrm{wt}_{0}(\lambda)-\mathrm{wt}_{0}(\nu)} \underline{q}^{\underline{\mathrm{wt}}(\nu)} .
$$

On the other hand, it is easy to see that in fact for any $\mu \in p^{-1}(\lambda)$, the corresponding $\tilde{\sigma}(\mu)$ can also be obtained by pushing up beads exactly $N$ times, one step at a time, the difference being just in the runners on which these shifts are performed. This means that each $\mu$ differs from $\tilde{\sigma}(\mu)$ by the same number $N=\mathrm{wt}(\lambda)-\mathrm{wt}(\nu)$ of border strips. Therefore, we have

$$
\sum_{\mu \in p^{-1}(\lambda)} q^{\underline{\mathrm{wt}^{-}(\mu)}}=q^{\mathrm{wt}_{0}(\lambda)-\mathrm{wt}_{0}(\nu)} \sum_{\mu \in p^{-1}(\nu)} q^{\underline{\mathrm{wt}}(\mu)}
$$

This is clearly compatible with (8) and reduces the argument to 0 -generated core partitions.
Fix a 0 -generated core $\lambda \in \mathcal{Z}_{\Delta}^{0} \cap \mathcal{C}_{\Delta}$; using Lemma 2.9 again, the corresponding ( $n+1$ )-tuple is a set of nondecreasing integers $\underline{a}=\left(a_{0}, \ldots, a_{n}\right)$ summing to 0 . The fibre $p^{-1}(\lambda)$ consists of partitions
whose abacus representation contains the same number of beads in each row as $\lambda$. The shift of one bead to the left results in the removal in the partition of a block labelled $i$, with $1 \leq i \leq n$. After substitution, this multiplies the contribution of the diagram on the right hand side of (8) by $\xi^{-1}$. If we fix all but one row, which contains $k$ beads, then these contributions add up to

$$
\sum_{n_{1}=0}^{n-k+1} \sum_{n_{2}=0}^{n_{1}} \cdots \sum_{n_{k}=0}^{n_{k-1}}\left(\xi^{-1}\right)^{n_{1}+\cdots+n_{k}}=\binom{n+1}{k}_{\xi^{-1}}
$$

where $\binom{m}{r}_{z}=\frac{[m]_{z}!}{[r]_{z}![m-r]_{z}!}$ is the Gaussian binomial coefficient, with $[m]_{z}=\frac{1-z^{m}}{1-z}$.
The number of rows containing exactly $k$ beads in the configuration corresponding to $\lambda$ is $a_{n+1-k}-a_{n-k}$. Therefore, the total contribution of the preimages, the left hand side of (8), is

$$
\begin{aligned}
\left.\sum_{\mu \in p^{-1}(\lambda)} \underline{q}^{\underline{\mathrm{wt}}(\mu)}\right|_{q_{1}=\cdots=q_{n}=\xi, q_{0}=\xi^{-n} q} & =\left.\prod_{k=1}^{n}\binom{n+1}{k}_{\xi^{-1}}^{a_{n+1-k}-a_{n-k}} \underline{q^{\underline{\mathrm{wt}}}(\lambda)}\right|_{q_{1}=\cdots=q_{n}=\xi, q_{0}=\xi^{-n} q} \\
& =\left.\prod_{l=0}^{n}\left(\frac{\binom{n+1}{n+1-l}_{\xi^{-1}}}{\binom{n+1}{n-l}_{\xi^{-1}}}\right)^{a_{l}} \underline{q}^{\underline{\mathrm{wt}}(\lambda)}\right|_{q_{1}=\cdots=q_{n}=\xi, q_{0}=\xi^{-n} q} \\
& =\left.\prod_{l=0}^{n}\left(\frac{1-\xi^{-l-1}}{1-\xi^{l-n-1}}\right)^{a_{l}} \underline{\underline{q}^{\underline{\mathrm{wt}}}(\lambda)}\right|_{q_{1}=\cdots=q_{n}=\xi, q_{0}=\xi^{-n} q} \\
& =\left.\prod_{l=1}^{n}\left(\frac{1-\xi^{-n-1}}{1-\xi^{-1}} \frac{1-\xi^{-l-1}}{1-\xi^{l-n-1}}\right)^{a_{l}} \underline{q^{\underline{\mathrm{wt}}}(\lambda)}\right|_{q_{1}=\cdots=q_{n}=\xi, q_{0}=\xi^{-n} q} \\
& =\left.\xi^{-\sum_{l=1}^{n} l a_{l}} \underline{q}^{\underline{\mathrm{wt}}(\lambda)}\right|_{q_{1}=\cdots=q_{n}=\xi, q_{0}=\xi^{-n} q}
\end{aligned}
$$

where in the second equality we used $\binom{n+1}{0}_{z}=\binom{n+1}{n+1}_{z}=1$, in the penultimate equality we used $a_{0}=-a_{1}-\cdots-a_{n}$, and in the last equality we used

$$
\frac{1-\xi^{-n-1}}{1-\xi^{-1}} \frac{1-\xi^{-l-1}}{1-\xi^{l-n-1}}=\xi^{-l}
$$

which can be checked to hold for $\xi$ a primitive $(n+2)$-nd root of unity. Incidentally, as the multiplicative order of $\xi$ is exactly $n+2$, all the denominators appearing above are non-vanishing. Finally, according to [14, §2], we have

$$
\underline{q}^{\underline{\mathrm{wt}}(\lambda)}=q^{\underline{Q(\underline{a})}} q_{1}^{a_{1}+\cdots+a_{n}} \cdot \ldots \cdot q_{n}^{a_{n}}
$$

where again $q=q_{0} \cdot \ldots \cdot q_{n}$ and $Q: \mathbb{Z}^{n} \rightarrow \mathbb{Z}$ is the quadratic form associated to $C_{\Delta}$. Since $q_{0}$ appears only in $q$ on the right hand side, it is clear that $\frac{Q(\underline{a})}{2}=\mathrm{wt}_{0}(\lambda)$. Hence,

$$
\left.q^{\frac{Q(a)}{2}} q_{1}^{a_{1}+\cdots+a_{n}} \ldots q_{n}^{a_{n}}\right|_{q_{1}=\cdots=q_{n}=\xi}=q^{\mathrm{wt}_{0}(\lambda)} \xi^{\sum_{l=1}^{n} l a_{l}} .
$$

This concludes the proof.

## 3. Type $D_{n}$ : ideals and Young walls

3.1. The binary dihedral group. Fix an integer $n \geq 4$, and let $\Delta$ be the root system of type $D_{n}$. For $\varepsilon$ a fixed primitive $(2 n-4)$-th root of unity, the corresponding subgroup $G_{\Delta}$ of $S L(2, \mathbb{C})$ can be generated by the following two elements $\sigma$ and $\tau$ :

$$
\sigma=\left(\begin{array}{cc}
\varepsilon & 0 \\
0 & \varepsilon^{-1}
\end{array}\right), \quad \tau=\left(\begin{array}{cc}
0 & 1 \\
-1 & 0
\end{array}\right)
$$

The group $G_{\Delta}$ has order $4 n-8$, and is often called the binary dihedral group. We label its irreducible representations as shown in Table 1 There is a distinguished 2-dimensional representation, the defining representation $\rho_{\text {nat }}=\rho_{2}$. See [23, 7] for more detailed information.

We will often meet the involution on the set of representations of $G_{\Delta}$ which is given by tensor product with the sign representation $\rho_{1}$ : on the set of indices $\{0, \ldots, n\}$, this is the involution $j \mapsto$ $\kappa(j)$ which swaps 0 and 1 and $n-1$ and $n$, fixing other values $\{2, \ldots, n-2\}$. Given $j \in\{0, \ldots, n\}$, we denote $\kappa(j, k)=\kappa^{k n}(j)$; this is an involution which is nontrivial when $k$ and $n$ are odd, and trivial otherwise. The special case $k=1$ will also be denoted as $\underline{j}=\kappa^{n}(j)$.

| $\rho$ | $\operatorname{Tr}(1)$ | $\operatorname{Tr}(\sigma)$ | $\operatorname{Tr}(\tau)$ |
| :---: | :---: | :---: | :---: |
| $\rho_{0}$ | 1 | 1 | 1 |
| $\rho_{1}$ | 1 | 1 | -1 |
| $\rho_{2}$ | 2 | $\varepsilon+\varepsilon^{-1}$ | 0 |
| $\vdots$ |  | $\vdots$ |  |
| $\rho_{n-2}$ | 2 | $\varepsilon^{n-3}+\varepsilon^{-(n-3)}$ | 0 |
| $\rho_{n-1}$ | 1 | -1 | $-i^{n}$ |
| $\rho_{n}$ | 1 | -1 | $i^{n}$ |

TABLE 1. Labelling the representations of the group $G_{\Delta}$

The following identities will be useful:

$$
\begin{equation*}
\rho_{n-1}^{\otimes 2} \cong \rho_{n}^{\otimes 2} \cong \rho_{\underline{0}}, \quad \rho_{n-1} \otimes \rho_{n} \cong \rho_{\underline{1}}, \quad \rho_{1} \otimes \rho_{n-1} \cong \rho_{n}, \quad \rho_{1} \otimes \rho_{n} \cong \rho_{n-1}, \quad \rho_{1}^{\otimes 2} \cong \rho_{0} \tag{9}
\end{equation*}
$$

3.2. Young wall pattern and Young walls. We describe here the type $D$ analogue of the set of labelled partitions used in type $A$, following [25, 27]. In this section, we only describe the combinatorics; see Appendix A for the representation-theoretic significance of this set.

First we define the Young wall pattern of typ $\rrbracket^{17} D_{n}$, the analogue of the $(n+1)$-labelled positive quadrant lattice of type $A_{n}$ used above. This is the following infinite pattern, consisting of two types of blocks: half-blocks carrying possible labels $j \in\{0,1, n-1, n\}$, and full blocks carrying possible labels $1<j<n-1$ :


Next, we define the set of Young wall $\&^{2}$ of type $D_{n}$. A Young wall of type $D_{n}$ is a subset $Y$ of the infinite Young wall of type $D_{n}$, satisfying the following rules.
(YW1) $Y$ contains all grey half-blocks, and a finite number of the white blocks and half-blocks.
(YW2) $Y$ consists of continuous columns of blocks, with no block placed on top of a missing block or half-block.
(YW3) Except for the leftmost column, there are no free positions to the left of any block or halfblock. Here the rows of half-blocks are thought of as two parallel rows; only half-blocks of the same orientation have to be present.

[^1](YW4) A full column is a column with a full block or both half-blocks present at its top; then no two full columns have the same height ${ }^{3}$.
Let $\mathcal{Z}_{\Delta}$ denote the set of all Young walls of type $D_{n}$. For any $Y \in \mathcal{Z}_{\Delta}$ and label $j \in\{0, \ldots, n\}$ let $w t_{j}(Y)$ be the number of white half-blocks, respectively blocks, of label $j$. These are collected into the multi-weight vector $\underline{w t}(Y)=\left(w t_{0}(Y), \ldots, w t_{n}(Y)\right)$. The total weight of $Y$ is the sum
$$
|Y|=\sum_{j=0}^{n} w t_{j}(Y)
$$
and for the formal variables $q_{0}, \ldots, q_{n}$,
$$
\underline{q}^{\underline{w t}(Y)}=\prod_{j=0}^{n} q_{j}^{w t_{j}(Y)}
$$
3.3. Decomposition of $\mathbb{C}[x, y]$ and the transformed Young wall pattern. The group $G_{\Delta}$ acts on the affine plane $\mathbb{C}^{2}$ via the defining representation $\rho_{\text {nat }}=\rho_{2}$. Let $S=\mathbb{C}[x, y]$ be the coordinate ring of the plane, then $S=\oplus_{m \geq 0} S_{m}$ where $S_{m}$ is the $m$ th symmetric power of $\rho_{\text {nat }}$, the space of homogeneous polynomials of degree $m$ of the coordinates $x, y$.

We further decompose

$$
S_{m}=\bigoplus_{j=0}^{n} S_{m}\left[\rho_{j}\right]
$$

into subrepresentations indexed by irreducible representations. We will also use this notation for linear subspaces: for $U \subset S_{m}$ a linear subspace, $U\left[\rho_{j}\right]=U \cap S_{m}\left[\rho_{j}\right]$. We will call an element $f \in S$ degree homogeneous, if $f \in S_{m}$ for some $m$; we call it degree and weight homogeneous, if $f \in S_{m}\left[\rho_{j}\right]$ for some $m, j$.

The decomposition of $S$ into $G_{\Delta}$-summands can be read off very conveniently from the transformed Young wall pattern. The transformation is an affine one, involving a shear: reflect the original Young wall pattern in the line $x=y$ in the plane, translate the $n$th row by $n$ to the right, and remove the grey triangles of the original pattern. In this way, we get the following picture:


As it can be checked readily, this is a representation of $S$ and its decomposition into $G_{\Delta^{-}}$ representations. The homogeneous components $S_{m}$ are along the antidiagonals. For $1<i<n-1$, a full block labelled $j$ below the diagonal, together with its mirror image, correspond to a 2 dimensional representation $\rho_{j}$. For $j \in\{0,1, n-1, n\}$, a full block labelled $j$ on the diagonal, as well as a half-block labelled $j$ below the diagonal with its mirror image, corresponds to a onedimensional representation. The dimension of $S_{m}\left[\rho_{j}\right]$ is the same as the total number of full blocks labelled $j$ on the $m$ th diagonal in the transformed Young wall pattern, counting mirror images also.

It is easy to translate the conditions (YW1)-(YW4) into the combinatorics of the transformed pattern; see Proposition 3.7 and Remark 3.8 below. Pictures of some small Young walls in the transformed pattern can be found below in Examples 4.54 .9 below.
3.4. Subspaces and operators. For each non-negative integer $m$ and irreducible representation $\rho_{j}$, consider the space $P_{m, j}$ of nontrivial $G_{\Delta}$-invariant subspaces of minimal dimension in $S_{m}\left[\rho_{j}\right]$. Specifically, if $\rho_{j}$ is one-dimensional, then these will be lines, and $P_{m, j}$ is simply the projectivization $\mathbb{P} S_{m}\left[\rho_{j}\right]$. If $\rho_{j}$ is two-dimensional, then $P_{m, j}$ is a closed subvariety of $\operatorname{Gr}\left(2, S_{m}\left[\rho_{j}\right]\right)$. It is easy to see that in this case also, $P_{m, j}$ is isomorphic to a projective space.

[^2]More generally, let $G_{m, j}^{r}$ be the space of $(r-1)$-dimensional projective subspaces of $P_{m, j}$. If $\rho_{j}$ is one-dimensional, then this is the Grassmannian $\operatorname{Gr}\left(r, S_{m}\left[\rho_{j}\right]\right)$. When $\rho_{j}$ is two-dimensional, then $G_{m, j}^{r}$ is a closed subvariety of $\operatorname{Gr}\left(2 r, S_{m}\left[\rho_{j}\right]\right)$ isomorphic to a Grassmannian of rank $r$. Clearly $G_{m, j}^{1}=P_{m, j}$.

For $0 \leq j \leq n$, we introduce operators $L_{j}: \operatorname{Gr}(S) \rightarrow \operatorname{Gr}(S)$ on the Grassmannian $\operatorname{Gr}(S)$ of all linear subspaces of the vector space $S$ as follows: for $v \in \operatorname{Gr}(S)$, we set
(1) $L_{0} v=v$;
(2) $L_{1} v=x y \cdot v$;
(3) for $1<j<n-1, L_{j} v=\left\langle x^{j-1} \cdot v, y^{j-1} \cdot v\right\rangle$;
(4) $L_{n-1} v=\left(x^{n-2}-i^{n} y^{n-2}\right) \cdot v$;
(5) $L_{n} v=\left(x^{n-2}+i^{n} y^{n-2}\right) \cdot v$.

Sometimes we will use the notation $L_{2}=L_{2, x}+L_{2, y}$ for the $x$ - and $y$-component of the operator $L_{2}$, i.e. multiplication with $x$, respectively $y$. The operators above restrict to operators $L_{0}: \operatorname{Gr}\left(S_{m}\right) \rightarrow \operatorname{Gr}\left(S_{m}\right), L_{1}: \operatorname{Gr}\left(S_{m}\right) \rightarrow \operatorname{Gr}\left(S_{m+2}\right), L_{j}: \operatorname{Gr}\left(S_{m}\right) \rightarrow \operatorname{Gr}\left(S_{m+j-1}\right)$ for $1<j<n-1$, and $L_{n-1}, L_{n}: \operatorname{Gr}\left(S_{m}\right) \rightarrow \operatorname{Gr}\left(S_{m+n-2}\right)$ on the Grassmannians of the graded pieces $S_{m}$. To simplify notation, if we do not write the space to which these operators are applied, then application to $\langle 1\rangle$ is meant. So, for example, the symbol $L_{1}^{2}$ standing alone denotes the vector subspace $\left\langle x^{2} y^{2}\right\rangle$ of $S_{4}$, while $L_{2}$ alone denotes the two-dimensional vector subspace $\langle x, y\rangle$ of $S_{1}$. For a linear subspace $v$ of $S$, the sum $\sum_{j \in I} L_{j} v$ denotes the subspace of $S$ generated by the images $L_{j} v$. We use the operator notation also for a set of subspaces; the meaning should be clear from the context.
3.5. Cell decompositions of equivariant Grassmannians. We start this section by defining decompositions of the Grassmannians $P_{m, j}$ of nontrivial $G_{\Delta}$-invariant subspaces of minimal dimension in $S_{m}\left[\rho_{j}\right]$. Given $(m, j)$, let $B_{m, j}$ denote the set of pairs of non-negative integers $(k, l)$ such that $k+l=m, l \geq k$, and the block position $(k, l)$ on the $m$-th antidiagonal on or below the main diagonal contains a block or half block of color $j$. Here $k$ is the row index, $l$ is the column index, and both of them in a nonnegative integer. It clearly follows from our setup that

$$
\operatorname{dim} P_{m, j}=\left|B_{m, j}\right|-1
$$

Proposition 3.1. Given $(m, j)$, there exists a locally closed stratification

$$
P_{m, j}=\bigsqcup_{(k, l) \in B_{m, j}} V_{k, l, j}
$$

which is a standard stratification of the projective space $P_{m, j}$ into affine spaces $V_{k, l, j}$ of decreasing dimension.

We will call $V_{k, l, j}$ the cells of $P_{m, j}$. The decomposition will be defined inductively, based on the following Lemma. Recall that $j \mapsto \kappa(j)$ denotes the involution on $\{0, \ldots, n\}$ which swaps 0 and 1 and $n-1$ and $n$.

Lemma 3.2. For any $l \geq 0$ and any $j \in[0, n]$, we have an injection

$$
L_{1}: P_{l-2, j} \rightarrow P_{l, \kappa(j)}
$$

This map is an isomorphism except in the case when the block or half-block in the bottom row of the transformed Young wall pattern on the l-th antidiagonal has label $j$, in which case the image has codimension one.

Proof. It is clear that multiplication by $L_{1}$ induces an injection, so we simply need to check the dimensions. The statement then clearly follows by looking at the transformed Young wall pattern: multiplication by $L_{1}$ corresponds to shifting the $(l-2)$-nd diagonal up by one diagonal step to the $l$-th diagonal; the number of blocks or half-blocks labelled $j$ is identical, unless the new (half-)block has label $j$, and then the codimension is exactly one.

Remark 3.3. The half-block in the bottom row of the transformed Young wall pattern in the $l$-th antidiagonal has label $j=0,1$ for $l \equiv 0 \bmod (2 n-4)$ except at $(0,0)$ where only 0 occurs. Half-blocks labelled $j=(n-1), n$ occur for $l \equiv n-2 \bmod (2 n-4)$. For $j \in[2, n-2]$, there are full blocks labelled $j$ in the bottom row on antidiagonals for $l \equiv j-1$ or $2 n-3-j \bmod (2 n-4)$.

Proof of Proposition 3.1. Nontrivial cells $V_{0, l, j}$ need to be defined exactly when the block or halfblock in the bottom row of the transformed Young wall pattern in the $l$-th antidiagonal has label $j$. In these cases, we set the cells along the bottom row to be

$$
V_{0, l, j}=P_{l, j} \backslash L_{1} P_{l-2, \kappa(j)}
$$

Once the cells $V_{0, k, j}$ along the bottom row are defined, we define the general cells for all $0 \leq j \leq n$, all $l$ and $k$ by

$$
V_{k, k+l, j}=L_{1}^{k} V_{0, l, \kappa^{k}(j)}
$$

What this says is that the cells are shifted up diagonally by $L_{1}$, taking into account that $L_{1}$ multiplies by the sign representation, so shifts the indices by the appropriate power of the involution $\kappa$. By induction, we obtain a decomposition of $P_{m, j}$ with the stated properties.

As it is well known, a decomposition of a projectivization of a vector space into affine cells is equivalent to giving a flag in the space itself. This induces a natural decomposition of all higher rank Grassmannians into Schubert cells, which are known to be affine. Thus our cell decomposition of $P_{m, j}$ induces cell decompositions of all $G_{m, j}^{r}$. Since the cells in the first decomposition are indexed by the set $B_{m, j}$, the cells in the second will be indexed by subsets of $B_{m, j}$ of size $r$. A Schubert cell of $G_{m, j}^{r}$ corresponding to a subset $S=\left\{\left(k_{1}, l_{1}\right), \ldots\left(k_{r}, l_{r}\right)\right\} \subset B_{m, j}$ will consist of those $(r-1)$ dimensional projective subspaces of $P_{m, j}$ that intersect $V_{k_{i}, l_{i}, j}$ nontrivially for all $1 \leq i \leq r$. We will denote the cell corresponding to $S$ in $G_{m, j}^{r}$ by $V_{S, j}$. We obtain a locally closed decomposition

$$
G_{m, j}^{r}=\bigsqcup_{\substack{S \subseteq B_{m, j} \\|S|=r}} V_{S, j} .
$$

Occasionally, when it is clear from the context that $S$ is a subset of $B_{m, j}$, we will supress the index $j$ and write just $V_{S}$ for the Schubert cells of $G_{m, j}^{r}$.

We will call a Schubert cell maximal if it intersects the maximal dimensional cell of $P_{m, j}$ nontrivially. Such a cell corresponds to subsets $S \subset B_{m, j}$ which contain ( $k_{m i n}, l$ ) where $k_{m i n}$ is minimal among the first components of the elements of $B_{m, j}$. The intersection with $V_{k_{m i n}, l, j}$ of a subspace corresponding to a point in a maximal Schubert cell is an affine subspace of $V_{k_{m i n}, l, j}$. Conversely, to any affine subspace of $V_{k_{m i n}, l, j}$, there corresponds a point in a maximal Schubert cell given by the completion of the subspace in $P_{m, j}$.

For a maximal subset $S$, denote by $\bar{S} \subset B_{m, j}$ the set of indices which we get by deleting ( $k_{m i n}, l$ ) from $S$. $\bar{S}$ is empty, if $|S|=1$. Define the codimension one projective subspace

$$
\bar{P}_{m, j}=\bigsqcup_{\left\{(k, l) \in B_{m, j}: k>k_{m i n}\right\}} V_{k, l, j} \subset P_{m, j}=P_{m, j} \backslash V_{k_{m i n}, l, j}
$$

For each $(r-1)$-dimensional subspace $U \subset P_{m, j}$ intersecting the affine space $V_{k_{m i n}, l, j}$ nontrivially, let $\bar{U}=U \cap \bar{P}_{m, j}$.

Lemma 3.4. The map $\omega: V_{S, j} \rightarrow V_{\bar{S}, j}$ defined by $\omega(U)=\bar{U}$ is a trivial affine fibration with fibre $\mathbb{A}^{\left|B_{m, j}\right|-|S|}$.

We can think of this map as associating to an affine subspace of $V_{k_{m i n}, l, j}$ its set of "ideal points at infinity". We define the mapping $\omega: V_{S, j} \rightarrow V_{\bar{S}, j}$ as the identity for those index sets and the corresponding cells which are not maximal. In such cases, $\bar{S}=S$ considered as a subset of $B_{m, j} \backslash\left\{\left(k_{\min }, l\right)\right\}$.

Consider the fibre $\omega^{-1}(\bar{U})$ over a point $\bar{U} \in V_{\bar{S}}$, which we will also denote by $\left.V_{S}\right|_{\bar{U}}$ below. This fibre consists of those subspaces $U \subseteq P_{m, j}$ which intersect $\bar{P}_{m, j}$ in $\bar{U}$, i.e. when considered as an affine subspace of $V_{k_{m i n}, l, j}$, they have $\bar{U}$ as their set of "points at infinity". We will denote the set of such subspaces also by $V_{k_{m i n}, l, j} / \bar{U}$. This notation means that we take the cosets in $V_{k_{m i n}, l, j}$ of an arbitrary affine subspace $U \subset V_{k_{m i n}, l, j}$ with $\bar{U}=\bar{P}_{m, j} \cap U$. The affine structure on $V_{k_{m i n}, l, j}$ descends to an affine structure on $V_{k_{m i n}, l, j} / \bar{U}$ which does not depend on the particular affine subspace $U$ whose cosets were taken.

We also need a description of the affine subspaces of the cells $V_{k, l, j}$ for $k>k_{\min }$. The relevant Schubert cells in this case are indexed by those subsets $S$ of $B_{m, j}$ which contain $(k, l)$ but do not contain any $\left(k^{\prime}, l^{\prime}\right)$ for $k^{\prime}<k$. Hence, the index set $B_{m, j}$ is first truncated by deleting the pairs
$\left(k^{\prime}, l^{\prime}\right)$ with $k^{\prime}<k$. We will denote the result as $B_{m, j}(k)$. Then, the maximal Schubert cells for $V_{k, l, j}$ correspond to those subsets $S \subseteq B_{m, j}(k)$ of the truncated index set which contain $(k, l)$. For these, $\bar{S}$ is defined by removing $(k, l)$ from $S$. There is still a morphism $\omega: V_{S, j} \rightarrow V_{\bar{S}, j}$ which is defined in the same way as above; its global structure and the description of its fibres is analogous to the previous special case.

Note that the notation $\bar{S}$ is ambiguous at this point: any maximal subset $S \subseteq B_{m, j}(k)$ can also be considered as a nonmaximal subset of $B_{m, j}\left(k^{\prime}\right)$ for $k^{\prime}<k$. If we view $S$ as a subset of $B_{m, j}(k)$, then $\bar{S}=S \backslash\{(k, m-k)\}$. On the other hand, if we view it as a nonmaximal subset of $B_{m, j}\left(k^{\prime}\right)$ for $k^{\prime}<k$, then $\bar{S}=S$. We have decided not to introduce extra notation; when this notation gets used below, we will always specify the reference point $k$ explicitly.
3.6. The Young wall associated to a homogeneous ideal. In this section, we study ideals generated by degree- and weight-homogeneous polynomials; we will call such ideals simply homogeneous ideals. Here is the main definition of this section.

Definition 3.5. Consider a homogeneous $G_{\Delta}$-invariant ideal $I \triangleleft \mathbb{C}[x, y]$. Let $Y_{I}$ denote the following subset of the transformed Young wall pattern of type $D_{n}$ : for each block or half-block $(k, l)$ of label $j$, with $k+l=m$, include this block or half-block in $Y_{I}$ if and only if $I \cap S_{m}\left[\rho_{j}\right]$ does not intersect the preimage in $S_{m}\left[\rho_{j}\right]$ of the stratum $V_{k, l, j} \subset P_{m, j}$ from the stratification of Proposition 3.1. $Y_{I}$ will be called the profile of $I$.

It will be useful to introduce a little bit of extra notation, and to reformulate this definition in the new notation. Given a homogeneous $G_{\Delta}$-invariant ideal $I$, let $I_{m, j}$ be the set of $G_{\Delta}$-invariant subspaces of minimal dimension in $I \cap S_{m}\left[\rho_{j}\right]$. Then as $I \cap S_{m}\left[\rho_{j}\right] \subset S_{m}\left[\rho_{j}\right]$ is a linear subspace, $I_{m, j} \subset P_{m, j}$ is a projective linear subspace. Then the definition simply says that a block or halfblock $(k, l)$ labelled $j$ is included in $Y_{I}$ if and only if $I_{m, j} \cap V_{k, l, j}=\emptyset$, for $m=k+l$ as before. Since $\left\{V_{k, l, j}\right\}$ is a standard stratification of the projective space $P_{m, j}$ into affine spaces, $\left\{I_{m, j} \cap V_{k, l, j}\right\}$ is also a standard stratification of its projective linear subspace $I_{m, j}$ into affine spaces, and so has the same number of strata as its affine dimension. We conclude

Lemma 3.6. For all $m, j$, the number of absent blocks or half-blocks of label $j$ on the $m$-th diagonal equals $\operatorname{dim}\left(I \cap S_{m}\left[\rho_{j}\right]\right)$.
Proposition 3.7. Given a homogeneous $G_{\Delta}$-invariant ideal $I \triangleleft \mathbb{C}[x, y]$, the associated subset $Y_{I}$ of the transformed Young wall pattern of type $D_{n}$ has the following properties.
(1) If a full or half block is missing, then all the blocks above-right from it on the diagonal are missing
(2) If a full block is missing, then all full or half blocks to the right of it are missing, and at least one (full or half) block immediately above it is missing.
(3) If a half block is missing, then the full block to the right of it is missing.
(4) If both half-blocks sharing the same block position are missing, then the full block immediately above this position is missing.
In particular, if $I$ is of finite codimension, then $Y_{I}$ is a Young wall of type $D_{n}$, an element of the set $\mathcal{Z}_{\Delta}$.

Remark 3.8. As it can be checked from the definitions, the relationship between the directions in the original and transformed Young wall patterns is the following: (right, up, diagonally right and down) in the original correspond after transformation to (diagonally right and up, right, up) respectively. This way, it is easy to check the correspondence between the rules for Young walls from 3.2 and this proposition.
Proof of Proposition 3.7. Fix a homogeneous invariant ideal $I \triangleleft \mathbb{C}[x, y]$ and let $Y_{I}$ be the corresponding subset of the Young wall pattern. Property (1) of $Y_{I}$ follows by applying $L_{1}$, recalling the inductive nature of the stratification of $P_{m, j}$ using $L_{1}$. The inductice construction also implies that it suffices to check properties (2)-(4) for blocks missing on the bottom row.

Let us next prove (2) in the general case, when a full block in position ( $0, l$ ) in representation $j \in[3, n-3]$ is missing from $Y_{I}$; by the choice of $j$, both above and to the right of this block there are also full blocks. Since the block at $(0, l)$ is missing, there is an invariant 2-dimensional subspace

[^3]$u \in I_{l, j} \cap V_{0, l, j}$ contained in $I$. Since $u$ is in the lowest stratum $V_{0, l, j}$ of $P_{l, j}$, it has a basis one of whose members at least is not divisible by $x y$; without loss of generality, we may assume that this polynomial is $x^{a} p$ where $a$ is a non-negative integer and $p$ a polynomial in $x, y$ not divisible by $x, y$. Now we can write
$$
L_{2} u=u_{+} \oplus u_{-}
$$
with $u_{+} \in I_{l+1, j+1}$ and $u_{-} \in I_{l+1, j-1}$. Then $u_{+}$must contain a polynomial with $x^{a+1} p$ as nonzero summand, so it cannot be in the image of $L_{1}$; so we have $u_{+} \in V_{0, l+1, j+1}$. Similarly, $u_{-}$must contain a polynomial with $x^{a} y p$ as nonzero summand, so it cannot be in the image of $L_{1}^{2}$ and so $u_{-} \in V_{1, l, j-1}$. Thus indeed both the blocks in positions $(0, l+1)$ and $(1, l)$ are missing as claimed.

Let us now consider what changes if $j$ is chosen such that there are half-blocks around. Suppose first that the half-blocks happen to lie to the right of our block labelled $j$. Then we have a decomposition

$$
L_{2} u=u_{+}^{1} \oplus u_{+}^{2} \oplus u_{-},
$$

with $u_{+}^{i}$ both one-dimensional. In this case, it is easy to check that the polynomial $x^{a+1} p$ cannot itself generate a one-dimensional eigenspace, so both $u_{+}^{i}$ will contain a polynomial with $x^{a+1} p$ as nonzero summand. Thus neither of these subspaces can be in the image of $L_{1}$, and so must lie in the large stratum. Hence both these blocks are missing.

Suppose now that the half-blocks happen to lie above our block labelled $j$. Then $L_{2} u$ is either three- or four-dimensional. In the general case, it has dimension four and there is a decomposition

$$
L_{2} u=u_{-}^{1} \oplus u_{-}^{2} \oplus u_{+}
$$

with $u_{-}^{i}$ both one-dimensional. In special situations $L_{2} u$ is only three dimensional, and one of the $u_{-}^{i}$ 's is missing (see Lemma 4.20 below for a detailed analysis). In any case, $x^{a+1} p$ will lie in $u_{+}$, forcing that subspace to be in the large stratum. The other relevant polynomial $x^{a} y p$ may or may not generate a one-dimensional invariant subspace, depending on the values of $a, p$; so at least one, possibly both, of $u_{-}^{1}, u_{-}^{2}$ lies in the image of $L_{1}$ but not $L_{1}^{2}$, forcing them to lie in the corresponding stratum. So at least one, but possibly both, of the corresponding half-blocks must be missing. We remark here that the other $u_{-}^{i}$, if present, can be divisible by a higher power of $L_{1}$. This implies that in this case the ideal generated by $u$ may have nontrivial intersection with the cells $V_{k, l, j}$ even with $l>0$.

The proofs of (3)-(4) follow the same pattern; we omit the details. Finally if $I$ is of finite codimension, then it contains $S_{m}$ for $m$ large enough, and so $Y_{I}$ contains only finitely many blocks and half-blocks.

## 4. Type $D_{n}$ : decomposition of the orbifold Hilbert scheme

4.1. The decomposition. The aim of this section is to prove the following result, which gives a constructive proof of Theorem 1.4 for type $D_{n}$.
Theorem 4.1. Let $G_{\Delta}$ be the subgroup of $S L(2, \mathbb{C})$ of type $D_{n}$. Then there is a locally closed decomposition

$$
\operatorname{Hilb}\left(\left[\mathbb{C}^{2} / G_{\Delta}\right]\right)=\bigsqcup_{Y \in \mathcal{Z}_{\Delta}} \operatorname{Hilb}\left(\left[\mathbb{C}^{2} / G_{\Delta}\right]\right)_{Y}
$$

of the equivariant Hilbert scheme $\operatorname{Hilb}\left(\left[\mathbb{C}^{2} / G_{\Delta}\right]\right)$ into strata indexed bijectively by the set $\mathcal{Z}_{\Delta}$ of Young walls of type $D_{n}$, with each stratum $\operatorname{Hilb}\left(\left[\mathbb{C}^{2} / G_{\Delta}\right]\right)_{Y}$ a non-empty affine space.
Proof. The affine plane $\mathbb{C}^{2}$ carries the diagonal $T=\mathbb{C}^{*}$-action, which commutes with the $G_{\Delta^{-}}$ action. The action of $T$ lifts to all the equivariant Hilbert schemes $\operatorname{Hilb}^{\rho}\left(\left[\mathbb{C}^{2} / G_{\Delta}\right]\right)$ which are themselves nonsingular. Thus the fixed point set

$$
\operatorname{Hilb}\left(\left[\mathbb{C}^{2} / G_{\Delta}\right]\right)^{T}=\sqcup_{\rho} \operatorname{Hilb}^{\rho}\left(\left[\mathbb{C}^{2} / G_{\Delta}\right]\right)^{T}
$$

is also a union of nonsingular varieties, and it consists of points representing homogeneous invariant ideals. The construction of 3.6 associates a Young wall $Y$ to each homogeneous invariant ideal $I \triangleleft \mathbb{C}[x, y]$. Since the construction uses a locally closed decomposition of the projective spaces $P_{m, j}$, the Young wall $Y$ also depends in a locally closed way on the ideal $I$, and thus we obtain a decomposition

$$
\operatorname{Hilb}\left(\left[\mathbb{C}^{2} / G_{\Delta}\right]\right)^{T}=\bigsqcup_{Y \in Z_{\Delta}} Z_{Y}
$$

into reduced locally closed subvarieties, where $Z_{Y} \subset \operatorname{Hilb}\left(\left[\mathbb{C}^{2} / G_{\Delta}\right]\right)^{T}$ is the locus of homogeneous invariant ideals $I$ with associated Young wall $Y$.

Let $\operatorname{Hilb}\left(\left[\mathbb{C}^{2} / G\right]\right)_{Y} \subset \operatorname{Hilb}\left(\left[\mathbb{C}^{2} / G\right]\right)$ denote the locus of ideals which flow to $Z_{Y}$ under the action of the torus $T$. Then by the Biatynicki-Birula theorem [3], there is a regular map Hilb $\left(\left[\mathbb{C}^{2} / G\right]\right)_{Y} \rightarrow$ $Z_{Y}$ which is a Zariski locally trivial fibration with affine space fibres, and a compatible $T$-action on the fibres. By Theorem4.3] below, the base is an affine space as well. Hence, by [3, Sect.3, Remarks], $\operatorname{Hilb}\left(\left[\mathbb{C}^{2} / G\right]\right)_{Y}$ is an algebraic vector bundle over this base, and hence trivial (Serre-Quillen-Suslin) [28]. Theorem 4.1 follows.

Remark 4.2. (1) As $\operatorname{Hilb}\left(\left[\mathbb{C}^{2} / G_{\Delta}\right]\right)=\operatorname{Hilb}\left(\mathbb{C}^{2}\right)^{G_{\Delta}} \subset \operatorname{Hilb}\left(\mathbb{C}^{2}\right)$ is a smooth subvariety, the universal family over $\operatorname{Hilb}\left(\mathbb{C}^{2}\right)$ restricts to a universal family over the equivariant Hilbert scheme $\operatorname{Hilb}\left(\left[\mathbb{C}^{2} / G_{\Delta}\right]\right)$. This restricts to a universal family of homogeneous invariant ideals $\mathcal{U} \triangleleft \mathcal{O}_{\text {Hilb }\left(\left[\mathbb{C}^{2} / G_{\Delta}\right]\right)^{T}} \otimes \mathbb{C}[x, y]$ over the $T$-fixed point set. Restricting this universal family $\mathcal{U}$ to each of the strata constructed above gives flat families of homogeneous invariant ideals $\mathcal{U}_{Y} \triangleleft \mathcal{O}_{Z_{Y}} \otimes \mathbb{C}[x, y]$ over each stratum $Z_{Y}$. It follows from the construction that the families $\mathcal{U}_{Y}$ are universal for flat families of homogeneous invariant ideals with associated Young wall $Y$. We will have occasion to use the universal property of the strata $Z_{Y}$ below.
(2) The diagonal $T=\mathbb{C}^{*}$-action on $\mathbb{C}^{2}$ induces the usual monomial grading on $\mathbb{C}[x, y]$, and when $I$ is homogeneous, also on the quotient $\mathbb{C}[x, y] / I$. For a $G_{\Delta}$-invariant ideal $I$, this quotient is a $G_{\Delta}$-representation. In the homogeneous $G_{\Delta}$-invariant case, the quotient has a multigraded Hilbert function: a $\mathbb{Z}$-grading induced by the $T=\mathbb{C}^{*}$-action and a grading (or rather labelling) by the set $\left\{\rho_{0}, \ldots, \rho_{n}\right\}$. By Lemma 3.6, the multigraded Hilbert function of a homogeneous invariant ideal $I \triangleleft \mathbb{C}[x, y]$ is determined by its associated Young wall $Y$.
The following is the main technical result of this section.
Theorem 4.3. For each $Y \in \mathcal{Z}_{\Delta}$, the stratum $Z_{Y}$ constructed above is isomorphic to a nonempty affine space.

Remark 4.4. We note that our proof of Theorem 4.3 below certainly provides some information about the dimension of the affine space $Z_{Y}$, and thus of the affine space $\operatorname{Hilb}\left(\left[\mathbb{C}^{2} / G_{\Delta}\right]\right)_{Y}$. We leave the study of these quantities, which could lead to a refinement of Theorem 1.3 in the Grothendieck ring of varieties for type $D_{n}$, for further study.

Our proof of Theorem 4.3, discussed below following some preparation, is a direct inductive proof. We start with a series of examples which exhibit the range of issues our proof will have to tackle; the discussions use results to be proved further below. Throughout we take the simplest example $n=4$, which exhibits all the nontrivial features.
Example 4.5. Let $Y_{1}$ be the triangle of size 3.


An invariant homogeneous ideal $I$ corresponding to this Young wall necessarily has a generator in $V_{0,3,2} \subset P_{3,2}$. The latter is a projective line whose points can be represented by expressions $\alpha_{0} L_{2} L_{3}+\alpha_{1} L_{2} L_{4}$. The affine line $V_{0,3,2}$ is given by $\alpha_{0}+\alpha_{1} \neq 0$. It is straightforward to check that a general point in $V_{0,3,2}$, that is, when $\left[\alpha_{0}: \alpha_{1}\right] \notin\{[1: 0],[0: 1]\}$, indeed generates an ideal $I$ with Young wall $Y_{1}$. However, when $\alpha_{0}$ or $\alpha_{1}$ become zero, then $I$ does not intersect $V_{1,3,4}$, respectively $V_{1,3,3}$, even though it should, so we have to add another generator to $I$ from within the corresponding cell (see Proposition 4.22(5) below). Both these cells are points, so there is no further choice to make and thus the space $Z_{Y_{1}}$ is isomorphic to an affine line. This example already illustrates the fact that even within a single stratum $Z_{Y}$, the minimal number of generators of an ideal $I$ with Young wall $Y$ can vary.

Example 4.6. Let $Y_{2}$ be the triangle of size 4. In this case, we get $Z_{Y_{2}} \cong Z_{Y_{1}} \cong \mathbb{A}^{1}$, the affine line of Example 4.5 due to the isomorphism $V_{0,3,2} \cong V_{0,4,0} \times V_{0,4,1}$, see Proposition 4.22(2) below, or repeat the same argument as above.

Example 4.7. Let $Y_{3}$ be the triangle of size 5.


For each fixed ideal $I$ with associated Young wall $Y_{3}$ must necessarily have an generator $f \in I$ such that $[f] \in V_{0,5,2}$. By construction, $V_{0,5,2}$ is an affine plane. This generator is, up to scalar, unique, since the block of label 2 in position $(0,5)$ is the only one with this label missing from the degree 5 antidiagonal. In other words, $I_{5,2} \cap V_{0,5,2}$ should be a point, otherwise the points at infinity of $I_{5,2} \cap V_{0,5,2}$ would intersect the other cells of degree 5 which is not allowed because of the shape of the Young wall. I must also intersect both $V_{1,5,0}$ and $V_{1,5,1}$, and again in essentially unique points, the corresponding blocks being the only $0 / 1$ blocks missing from the degree 6 antidiagonal. This puts the following constraint on the allowed $[f] \in V_{0,5,2}$. Use the isomorphism $L_{1}^{-1}$ which maps $V_{1,5,0} \sqcup V_{1,5,1}$, a disjoint union of an affine line and a point, to $V_{0,4,1} \sqcup V_{0,4,0}$. Map this locus into $V_{0,5,2}$, an affine plane, to obtain $M_{1}^{0} \sqcup M_{0}^{0} \subset V_{0,5,2}$ by taking the ideals generated by them (the curious notation $M_{1}^{0} \sqcup M_{0}^{0}$ for this locus is used here to be consistent with the general setup later, see the definitions after Lemma 4.231). Then by Proposition4.22(2) below, the ideal $\langle f\rangle$ intersects $P_{6,0}$ and $P_{6,1}$ at most in the correct cells $V_{1,5,0}$ and $V_{1,5,1}$ if and only if $[f] \in V_{0,5,2}$ lies on the linear join of the affine line $M_{0}^{0}$ and the point $M_{1}^{0}$ inside the affine plane $V_{0,5,2}$ but not on $M_{1}^{0} \sqcup M_{0}^{0}$. This join is the plane $V_{0,5,2}$ minus a punctured affine line $M_{1} \backslash M_{1}^{0}$, where $M_{1}$ is the line parallel to $M_{0}^{0}$, going through $M_{1}^{0}$. On this join, but away from the locus $M_{0}^{0} \sqcup M_{1}^{0}$ itself, we can set $I=\langle f\rangle$ to indeed get an ideal with Young wall $Y_{3}$. On the locus $M_{0}^{0} \sqcup M_{1}^{0}$ however, the ideal $\langle f\rangle$ itself will not actually meet both cells, so we have to add an arbitrary generator of the missed cell to $f$ to obtain an ideal of the correct Young wall. Over the affine line $M_{0}^{0}$, there is no choice, since $V_{1,5,0}$ is a point. But over the point $M_{1}^{0}$, we still have $V_{1,5,0}$, in other words an affine line, worth of choices. This sofar tells us that $Z_{Y_{3}}$ is the disjoint union of $V_{0,5,2} \backslash M_{1} \cong \mathbb{A}^{2} \backslash \mathbb{A}^{1}$ and $V_{1,5,0} \cong \mathbb{A}^{1}$.


To fully work out the geometry of $Z_{Y_{3}}$, note that $P_{5,2} \cong \mathbb{P}^{2}$ can be parametrized by expersions $\alpha_{0} L_{2} L_{3}^{2}+\alpha_{1} L_{2} L_{3} L_{4}+\alpha_{2} L_{2} L_{4}^{2}$. The locus $V_{0,5,2} \subset P_{5,2}$ is given by $\alpha_{0}+\alpha_{1}+\alpha_{2} \neq 0$. The image of $V_{0,4,0}$ in these coordinates is $M_{0}^{0}=\left\{\left(\alpha_{0}, 0, \alpha_{1}\right): \alpha_{0}+\alpha_{1} \neq 0\right\}$, while the image of $V_{0,4,1}$ is $M_{1}^{0}=\{(0,1,0)\}$. The linear combinations of the points in $M_{0}^{0}$ and $M_{1}^{0}$ cover the whole affine plane $V_{0,5,2}$, except a punctured line. For a general linear combination $a \cdot\left(\alpha_{0}, 0, \alpha_{1}\right)+b \cdot(0,1,0)$, the ideal generated by the corresponding $f$ intersects $V_{1,5,0} \times V_{1,5,1}$ in ( $L_{1} L_{3} L_{4}, \alpha_{0} L_{1} L_{3}^{2}+\alpha_{1} L_{1} L_{4}^{2}$ ). For $(a, b)=(1,0)$ we have to have an extra generator in $V_{1,5,0}$, while for $(a, b)=(0,1)$ we need an extra generator in $V_{1,5,1}$.

Consider a family of ideals which approaches the point $M_{1}^{0}$ from the direction $\left(\alpha_{0}, 0, \alpha_{1}\right)$. Then it can be shown by explicit calculation that the limit ideal contains the subspace generated by $\alpha_{0} L_{1} L_{3}^{2}+\alpha_{1} L_{1} L_{4}^{2} \in V_{1,5,1}$. This shows that $Z_{Y_{3}}$ can be obtained by blowing up the affine plane $V_{0,5,2}$ in its point $M_{1}^{0}$, and removing the proper transform of the punctured line $M_{1} \backslash M_{1}^{0}$ from this blowup. Thus $Z_{Y_{3}} \cong \mathbb{A}^{2}$. From the blowup construction, we also obtain a canonical morphism $Z_{Y_{3}} \rightarrow \mathbb{A}^{1}$, the restriction of the morphism $\mathrm{Bl}_{0} \mathbb{A}^{2} \rightarrow \mathbb{P}^{1}$ to the exceptional curve of the blowup.

Example 4.8. Let $Y_{3}^{\prime}$ be the Young wall


There is necessarily still a unique generator in $V_{0,5,2}$. The difference compared to $Y_{3}$ is that there is now no intersection with $V_{1,5,1}$ but an intersection with $V_{3,3,1}$. The cells of the $0 / 1$-blocks missing from the degree six diagonal are $V_{1,5,0}$ and $V_{3,3,1}$, both of which are zero dimensional. As before, we pull back these using $L_{1}^{-1}$, take the linear combinations of their images in $P_{5,2}$, and intersect this line with $V_{0,5,2}$. This is exactly the line $M_{1}$ from Example 4.7. This has one special point, $M_{1}^{0}$. If the new generator is (in the subspace represented by) this point, then it will not generate an ideal with shape $Y$, except if we keep the unique element of $V_{3,3,0}$ as a generator. In any case, $Z_{Y_{3}^{\prime}} \cong M_{1} \cong \mathbb{A}^{1}$.

Example 4.9. As a final example, it is well known that the minimal resolution of the singularity $\mathbb{C}^{2} / G_{\Delta}$ is given by the component $\operatorname{Hilb}{ }^{\rho_{\mathrm{reg}}}\left(\left[\mathbb{C}^{2} / \mathrm{G}\right]\right)$ of $\operatorname{Hilb}\left(\left[\mathbb{C}^{2} / \mathrm{G}\right]\right)$ corresponding to the regular representation [26]. The $\mathbb{C}^{*}$-fixed set on the minimal resolution consists of the $\mathbb{P}^{1}$ within the exceptional locus corresponding to the central node in the Dynkin diagram, as well as three isolated points on the other three $\mathbb{P}^{1}$ 's. The following five Young walls contribute for the regular representation.


A quick computation shows that $Z_{Y}$ is a point in each case, except for the last Young wall in the first row, when it is an affine line (as in Example 4.5 except that in this case there is also a generator in $V_{2,2,0}$ ). This affine line contains the point corresponding to the Young wall next to it in its closure, giving the central $\mathbb{P}^{1}$.
4.2. Incidence varieties. The purpose of this section is to introduce some incidence varieties inside products of the Schubert cells defined in 3.5. We state some propositions regarding these incidence varieties and morphisms between them, whose proofs we defer to 4.5 below. We discuss four different cases.
Case 4.2, 1 Assume that $m \equiv 0 \bmod n-2$ is a nonnegative integer, such that at position $(0, m)$ there is a divided block with labels $\left(c_{1}, c_{2}\right)$. Let $c$ be the label of the block at position $(1, m)$. Let $S_{c} \subseteq B_{m+1, c}$ be a nonempty maximal subset. Let $S_{1} \subseteq B_{m, c_{1}}$ and $S_{2} \subseteq B_{m, c_{2}}$ be two maximal subsets which are allowed by $S_{c}$. This means, by definition, that each block above every composite block whose two halves are both contained in $S_{1} \cup S_{2}$, and each block to the right of every composite block at least one half-block of which in $S_{1} \cup S_{2}$, is in $S_{c}$.

Consider the incidence varieties

$$
X_{S_{1}, S_{2}}^{S_{c}}=\left\{\left(U_{1}, U_{2}, U_{c}\right):\left(U_{1}, U_{2}\right) \cap P_{m+1, c} \subseteq U_{c}\right\} \subseteq V_{S_{1}} \times V_{S_{2}} \times V_{S_{c}, c}
$$

and

$$
Y_{\bar{S}_{1}, \bar{S}_{2}}^{S_{c}}=\left\{\left(\bar{U}_{1}, \bar{U}_{2}, U_{c}\right):\left(\bar{U}_{1}, \bar{U}_{2}\right) \cap P_{m+1, c} \subseteq U_{c}\right\} \subseteq V_{\bar{S}_{1}} \times V_{\bar{S}_{2}} \times V_{S_{c}, c} .
$$

Using the maps $\omega$ from Lemma 3.4, these varieties fit into the diagram

$$
\begin{array}{cccc}
X_{S_{1}, S_{2}}^{S_{c}} & \subseteq & V_{S_{1}} \times V_{S_{2}} \times V_{S_{c}, c} \xrightarrow{\operatorname{Id} \times \mathrm{Id} \times \omega} & V_{S_{1}} \times V_{S_{2}} \times V_{\bar{S}_{c}, c} \\
& \downarrow_{\omega \times \omega \times \mathrm{Id}} & \\
\omega \times \omega \times \mathrm{Id}
\end{array}
$$

Proposition 4.10. (1) The image of $X_{S_{1}, S_{2}}^{S_{c}}$ under the vertical morphism $V_{S_{1}} \times V_{S_{2}} \times V_{S_{c}, c} \rightarrow$ $V_{\bar{S}_{1}} \times V_{\bar{S}_{2}} \times V_{S_{c}, c}$ is precisely $Y_{\bar{S}_{1}, \bar{S}_{2}}^{S_{c}}$.
(2) The induced morphism $X_{S_{1}, S_{2}}^{S_{c}} \rightarrow Y_{\bar{S}_{1}, \bar{S}_{2}}^{S_{c}}$ is a trivial fibration over its image with affine fibers of dimension $\left|S_{c}\right|-\left|S_{1}\right|-\left|S_{2}\right|+1$.
(3) The horizontal morphism $X_{S_{1}, S_{2}}^{S_{c}} \rightarrow V_{S_{1}} \times V_{S_{2}} \times V_{\bar{S}_{c}, c}$ is injective.

Case 4.2, 2 Let $m \equiv 0 \bmod n-2$, but this time consider only one half block of label $c_{0}=\kappa(c)$ at the position $(0, m)$. Let $S_{c} \subseteq B_{m+2, c}$ be a nonempty maximal subset, and $S \subseteq B_{m, \kappa(c)}$ be a maximal subset which is allowed by $S_{c}$. This means that for each block in $S$ there is a block in $S_{c}$ at the top right corner. In analogy with the previous case, let

$$
X_{S}^{S_{c}}=\left\{\left(U, U_{c}\right):(U) \cap P_{m+2, c} \subseteq U_{c}\right\} \subseteq V_{S} \times V_{S_{c}, c}
$$

and

$$
Y_{\bar{S}}^{S_{c}}=\left\{\left(\bar{U}, U_{c}\right):(\bar{U}) \cap P_{m+2, c} \subseteq U_{c}\right\} \subseteq V_{\bar{S}} \times V_{S_{c}, c}
$$

which fit into the diagram


Proposition 4.11. (1) The image of $X_{S}^{S_{c}}$ under the vertical morphism $V_{S} \times V_{S_{c}, c} \rightarrow V_{\bar{S}} \times V_{S_{c}, c}$ is exactly $Y_{\bar{S}}^{S_{c}}$.
(2) The induced morphism $X_{S}^{S_{c}} \rightarrow Y_{\bar{S}}^{S_{c}}$ is a trivial fibration over its image with affine fibers of dimension $\left|S_{c}\right|-|S|$.
(3) The horizontal morphism $X_{S}^{S_{c}} \rightarrow V_{S} \times V_{\bar{S}_{c}, c}$ is injective.

Case 4.2. 3 Let $m \equiv 1 \bmod n-2$, and $c_{1}$ and $c_{2}$ the labels of the divided block immediately above the block at position $(0, m)$. Let $S_{1} \subseteq B_{m+1, c_{1}}, S_{2} \subseteq B_{m+1, c_{2}}$ be nonempty subsets at least one of which is maximal. Let moreover, $S \subseteq B_{m, j}$ be a maximal subset which is allowed by $S_{1}$ and $S_{2}$. In this case, this means the following: for each block $b$ in $S$, there is a divided block of with labels $\left(c_{1}, c_{2}\right)$ in the pattern either immediately above or to the right of $b$. In the first case, we require that at least one of these half-blocks is in $S_{1} \cup S_{2}$. In the second case, we require that both are contained in $S_{1} \cup S_{2}$.

Given this data, we define

$$
X_{S}^{S_{1}, S_{2}}=\left\{\left(U, U_{1}, U_{2}\right):(U) \cap P_{m+1, c_{1}} \subseteq U_{1},(U) \cap P_{m+1, c_{2}} \subseteq U_{2}\right\} \subseteq V_{S} \times V_{S_{1}, c_{1}} \times V_{S_{2}, c_{2}}
$$

and

$$
Y_{\bar{S}}^{S_{1}, S_{2}}=\left\{\left(\bar{U}, U_{1}, U_{2}\right):(\bar{U}) \cap P_{m+1, c_{1}} \subseteq U_{1},(\bar{U}) \cap P_{m+1, c_{2}} \subseteq U_{2}\right\} \subseteq V_{\bar{S}} \times V_{S_{1}, c_{1}} \times V_{S_{2}, c_{2}}
$$

We now have the following diagram:


Proposition 4.12. (1) The image of $X_{S}^{S_{1}, S_{2}}$ under the vertical morphism $V_{S} \times V_{S_{1}, c_{1}} \times$ $V_{S_{2}, c_{2}} \rightarrow V_{\bar{S}} \times V_{S_{1}, c_{1}} \times V_{S_{2}, c_{2}}$ is exactly $Y_{\bar{S}}^{S_{1}, S_{2}}$.
(2) The induced morphism $X_{S}^{S_{1}, S_{2}} \rightarrow Y_{\bar{S}}^{S_{1}, S_{2}^{S}}$ is a trivial fibration with fibers isomorphic to affine spaces of dimension $\left|S_{1}\right|+\left|S_{2}\right|-|S|$.
We remark that the analogue of (3) of Propositions 4.10 and 4.11 is not true in this case. What happens to $X_{S}^{S_{1}, S_{2}}$ when we project $V_{S_{1}, c_{1}} \times V_{S_{2}, c_{2}}$ to $V_{\bar{S}_{1}, c_{1}} \times V_{\bar{S}_{2}, c_{2}}$ will be the subject of 5 below. Case 4.2, 4 Finally, assume $m \not \equiv 0,1 \bmod n-2$ with a full block in position $(0, m)$. Let $c$ be the label of the full block immediately above this position, and $S_{c} \subseteq B_{m+1, c}$ a nonempty maximal subset. Let moreover $S \subseteq B_{m, j}$ be a maximal subset which is allowed by $S_{c}$; in this case, this means that above every block of $S$ there is a block in $S_{c}$. Consider the incidence varieties

$$
X_{S}^{S_{c}}=\left\{\left(U, U_{c}\right):(U) \cap P_{m+1, c} \subseteq U_{c}\right\} \subseteq V_{S} \times V_{S_{c}, c}
$$

and

$$
Y_{\bar{S}}^{S_{c}}=\left\{\left(\bar{U}, U_{c}\right):(\bar{U}) \cap P_{m+1, c} \subseteq U_{c}\right\} \subseteq V_{\bar{S}} \times V_{S_{c}, c} .
$$

There is the following diagram:


Proposition 4.13. (1) The image of $X_{S}^{S_{c}}$ under the vertical morphism $V_{S} \times V_{S_{c}, c} \rightarrow V_{\bar{S}} \times V_{S_{c}, c}$ is exactly $Y_{\bar{S}}^{S_{c}}$.
(2) The induced morphism $X_{S}^{S_{c}} \rightarrow Y_{\bar{S}}^{S_{c}}$ is a trivial fibration over its image with affine fibers of dimension $\left|S_{c}\right|-|S|$.
(3) The horizontal morphism $X_{S}^{S_{c}} \rightarrow V_{S} \times V_{\bar{S}_{c}, c}$ is injective.
4.3. Proof of Theorem 4.3, In this section we prove Theorem 4.3, thus completing the proof of Theorem 4.1, using the constructions and results stated in 4.2. Given a Young wall $Y \in \mathcal{Z}_{\Delta}$, we need to show that the corresponding stratum $Z_{Y}$ is an affine space. The following is the key combinatorial definition which underlies much of the rest of the paper.
Definition 4.14. Consider the Young wall $Y$, as usual in the transformed pattern. The salient blocks of $Y$ are those blocks in the complement of $Y$, whose absence from $Y$ does not follow from the shape of the rows below it, and which are at the leftmost positions in their rows with this property. In particular, these are

- missing half blocks under which there is a block in $Y$;
- missing undivided full blocks under which there is a block in $Y$;
- missing divided full blocks immediately to the right of the boundary of $Y$;
- the leftmost missing block(s) in the bottom row.

Given an ideal $I \in Z_{Y}$, it is easy to see that $I$ is necessarily generated by elements lying in cells corresponding to the salient blocks of $Y$. In most cases it is also true that all cells corresponding to salient blocks must contain a generator, but not always; we have already seen Example 4.5, where the divided missing blocks at position $(1,3)$ are salient blocks of $Y_{3}$, since they lie immediately to the right of the boundary of $Y_{3}$, but the corresponding cells do not necessarily contain generators of an ideal $I \in Z_{Y_{3}}$.

We start our analysis by defining maps from the strata $Z_{Y}$ to the Grassmannian cells defined in 3.5. Consider an arbitrary block or half-block of label $j$ at position $(k, l)$ in the Young wall pattern. Let $S(k, l, j) \subseteq B_{k+l, j}$ be the set of blocks of label $j$ from the $(k+l)$-th antidiagonal which are not in $Y$ and are above and including the position $(k, l) . S(k, l, j)$ is called the index set of $(k, l)$ in $Y$. If the block of label $j$ at position $(k, l)$ is not contained in $Y$, then the index set $S=S(k, l, j)$ contains ( $k, l$ ) as well. By the correspondence discussed at the end of 3.5 between the maximal Schubert cells of the relevant Grassmannian for $V_{k, l, j}$ and subsets of $B_{k+l, j}(k)$ which do contain $(k, l)$, there is a maximal Schubert cell $V_{S, j}$ corresponding to $S$. The affine cell $V_{S, j} \subset G_{k+l, j}^{r}$ for $r=|S|$ parameterizes certain affine subspaces of $V_{k, l, j}$, or, equivalently, projective subspaces of $P_{k+l, j}$, the projectivization of the space of degree $(k+l)$ homogeneous polynomials which transform in the representation $\rho_{j}$ with respect to $G_{\Delta}$. The correspondence is by projective closure of the corresponding affine subspace of $V_{k, l, j}$.

Lemma 4.15. For any block or half-block at position ( $k, l$ ) which is not contained in $Y$ and has index set $S=S(k, l, j)$, there is a morphism

$$
\begin{aligned}
Z_{Y} & \rightarrow V_{S, j} \\
I & \mapsto I \cap V_{k, l, j} .
\end{aligned}
$$

Proof. Assume first for simplicity that $(k, l)=(0, m)$. Let $\mathcal{U}_{Y} \triangleleft\left(\mathcal{O}_{Z_{Y}} \otimes \mathbb{C}[x, y]\right)$ be the universal family of homogeneous ideals over $Z_{Y}$ introduced in Remark 4.2. Consider the subfamily $\mathcal{V}=$ $\mathcal{U}_{Y} \cap\left(\mathcal{O}_{Z_{Y}} \otimes S_{m}\left[\rho_{j}\right]\right)$. This is a family of subspaces of dimension $r \operatorname{dim} \rho_{j}$ in $S_{m}\left[\rho_{j}\right]$ parameterized by $Z_{Y}$, where $r=|S|$. It is known, that the morphisms from $Z_{Y}$ to the equivariant Grassmannian $G_{m, j}^{r}$ of $S_{m}\left[\rho_{j}\right]$ are in one-to-one correspondence with $G_{\Delta}$ invariant subbundles of $\mathcal{O}_{Z_{Y}} \otimes S_{m}\left[\rho_{j}\right]$ of
rank $r \operatorname{dim} \rho_{j}\left[9\right.$, p. 88, Theorem 2.4]. Hence, there is a classifying morphism $Z_{Y} \rightarrow G_{m, j}^{r}$ inducing $\mathcal{V}$.

By the definition of $Z_{Y}$, the multigraded Hilbert polynomial of $\mathcal{U}_{Y}$ is constant. The Hilbert polynomial encodes the dimensions of the intersections with the cells of $P_{m, j}$. Therefore, over closed points of $Z_{Y}$ the elements of the family $\mathcal{V}$, when considered as projective subspaces of $P_{m, j}$, intersect exactly the cells $V_{k_{i}, l_{i}, j}$ for each element $\left(k_{i}, l_{i}\right) \in S$. As in 3.5, these projective subspaces are represented by points in the Schubert cell $V_{S, j} \subset G_{m, j}^{r}$. Hence, the image of the classifying morphism of $\mathcal{V}$ is inside the cell $V_{S, j}$. By the construction, the classifying morphism is just the same as taking the intersection of $I_{m, j}$ with $V_{0, m, j}$. This is denoted as $I \cap V_{0, m, j}$ above.

The general case follows similarly.
We will prove Theorem 4.1 by induction on the number of nonempty rows of $Y$. Consider an arbitrary Young wall $Y$ consisting of $l>0$ rows. Let $\bar{Y}$ denote the Young wall obtained from $Y$ by deleting its bottom row; we will call this the truncation of $Y$. Of course the labels of the half blocks are exchanged by $\kappa$, but we will supress this in the notations. The following result will be key to our induction.

Lemma 4.16. There exists a morphism of schemes

$$
\begin{aligned}
T: Z_{Y} & \rightarrow Z_{\bar{Y}} \\
I & \mapsto L_{1}^{-1}\left(I \cap L_{1} \mathbb{C}[x, y]\right) .
\end{aligned}
$$

Proof. Let $\mathcal{U}_{Y} \triangleleft\left(\mathcal{O}_{Z_{Y}} \otimes \mathbb{C}[x, y]\right)$ be the universal family of homogeneous ideals over $Z_{Y}$. Consider $\mathcal{I}=L_{1}^{-1}\left(\mathcal{U}_{Y} \cap L_{1}\left(\mathcal{O}_{Z_{Y}} \otimes \mathbb{C}[x, y]\right)\right)$. It is straightforward to check locally that this is still a sheaf of ideals in $\mathcal{O}_{Z_{Y}} \otimes \mathbb{C}[x, y]$. On closed points of $Z_{Y}$, it is also clear that the restriction has Young wall $\bar{Y}$. As mentioned above, the multigraded Hilbert polynomial of $\mathcal{U}_{Y}$ is constant. As $Z_{Y}$ is reduced, it then follows from [22, Ch III. Thm. 9.9] that $\mathcal{I}$ is a flat family of homogeneous ideals with Young wall $\bar{Y}$ over $Z_{Y}$. By Remark 4.2 there is a classifying morphism $Z_{Y} \rightarrow Z_{\bar{Y}}$ for this family, which is exactly the morphism $T$.

Next, we continue with an investigation of the combinatorics of the bottom two rows of our Young wall $Y$. The boundary of $Y$ in the transformed pattern is divided by the blocks into horizontal, vertical and diagonal straight line segments. The first two lines in the bottom can be connected in the following six possible ways:


Here a diagonal straight line borders a half block of the Young wall, which can be either a lower or an upper triangle. In the A cases the salient block in the bottom row is a full block, while in the B cases it is a half block.

Let the salient block of $Y$ in its bottom row be at position $(0, m)$. It can be either a divided or undivided full block, or a half block. In the first case, we have a type A corner at the bottom of $Y$, while in the second case there is a type B corner. As in 4.2 we need to consider four cases. In each case, we are going to define morphisms $Z_{Y} \rightarrow \mathcal{X}_{Y}$ and $Z_{\bar{Y}} \rightarrow \mathcal{Y}_{\bar{Y}}$ to incidence varieties defined in 4.2 .

Case 4.3. 1 Assume $m \equiv 0 \bmod n-2$ and we have vertex types A1 or A2 (A3 is not possible in this case). We are in the context of Case 4.2 1: the divided block at position $(0, m)$ has labels $\left(c_{1}, c_{2}\right)$, and index sets $S_{1}, S_{2}$; the block at position $(1, m)$ has label $c$, and index set $S_{c}$. By the Young wall rules for $Y, S_{1}, S_{2}$ is allowed by $S_{c}$. Lemma 4.15 implies that there is a morphism

$$
\begin{array}{ccc}
Z_{Y} & \rightarrow & V_{S_{1}} \times V_{S_{2}} \times V_{S_{c}} \\
I & \mapsto & \left(I \cap V_{0, m, c_{1}}, I \cap V_{0, m, c_{2}}, I \cap V_{1, m, c}\right) .
\end{array}
$$

By construction, the image of this morphism is contained in the incidence variety $X_{S_{1}, S_{2}}^{S_{c}} \subset V_{S_{1}} \times$ $V_{S_{2}} \times V_{S_{c}}$ from Case 4.2.1. Denote $\mathcal{X}_{Y}=X_{S_{1}, S_{2}}^{S_{c}} \subseteq V_{S_{1}} \times V_{S_{2}} \times V_{S_{c}}$; we thus obtain an induced morphism $Z_{Y} \rightarrow \mathcal{X}_{Y}$.

There is also a morphism

$$
\begin{array}{ccc}
Z_{\bar{Y}} & \rightarrow & V_{\bar{S}_{1}} \times V_{\bar{S}_{2}} \times V_{S_{c}} \\
I & \mapsto & \left(L_{1} I \cap V_{k_{1}, l_{1}, c_{1}}, L_{1} I \cap V_{k_{2}, l_{2}, c_{2}}, L_{1} I \cap V_{1, m, c}\right),
\end{array}
$$

where $\left(k_{i}, l_{i}\right)$ is the lowest block in $\bar{S}_{i}$ for $i=1,2$. We obtain an induced morphism $Z_{\bar{Y}} \rightarrow \mathcal{Y}_{\bar{Y}}$, where $\mathcal{Y}_{\bar{Y}}=Y_{\bar{S}_{1}, \bar{S}_{2}}^{S_{c}}$.
Case 4.3,2 Assume $m \equiv 0 \bmod n-2$ with vertex types B1 to B3. This is Case 4.2 2: the block at position $(0, m)$ has label $\kappa(c)$, and index set $S=S_{\kappa(c)}$; the block at position $(1, m+1)$ has label $c$, and index set $S_{c}$. We consider the morphisms

$$
\begin{array}{ccc}
Z_{Y} & \rightarrow & V_{S} \times V_{S_{c}} \\
I & \mapsto & \left(I \cap V_{0, m, \kappa(c)}, I \cap V_{1, m+1, c}\right),
\end{array}
$$

and

$$
\begin{array}{ccc}
Z_{\bar{Y}} & \rightarrow & V_{\bar{S}} \times V_{S_{c}} \\
I & \mapsto & \left(L_{1} I \cap V_{k, l, \kappa(c)}, L_{1} I \cap V_{1, m+1, c}\right),
\end{array}
$$

where again $(k, l)$ is the lowest block in $\bar{S}$. In this case we let $\mathcal{X}_{Y}=X_{S}^{S_{c}}$ and $\mathcal{Y}_{\bar{Y}}=Y_{\bar{S}}^{S_{c}}$. The images of the morphisms above are contained in these.
Case 4.3. 3 Assume $m \equiv 1 \bmod n-2$ with vertex types A1, A2, A3. This is Case 4.2. 3: $c_{1}$ and $c_{2}$ are the labels of the divided block immediately above the block at position $(0, m), S_{1} \subseteq B_{m+1, c_{1}}$, $S_{2} \subseteq B_{m+1, c_{2}}$ are their index sets, both (in cases A1 and A2) or one (in case A3) of which is maximal; $S$ is the index set of the block at position $(0, m)$ with label $j$. We get a morphism

$$
\begin{array}{ccc}
Z_{Y} & \rightarrow & V_{S} \times V_{S_{1}, c_{1}} \times V_{S_{2}, c_{2}} \\
I & \mapsto & \left(I \cap V_{0, m, j}, I \cap V_{1, m, c_{1}}, I \cap V_{1, m, c_{2}}\right)
\end{array}
$$

whose image is contained in $\mathcal{X}_{Y}=X_{S}^{S_{1}, S_{2}}$. In this way we obtain an induced morphism $Z_{Y} \rightarrow \mathcal{X}_{Y, \bar{Y}}$. Similarly, consider the morphism

$$
\begin{array}{ccc}
Z_{\bar{Y}} & \rightarrow & V_{\bar{S}} \times V_{S_{1}, c_{1}} \times V_{S_{2}, c_{2}} \\
I & \mapsto & \left(L_{1} I \cap V_{k, l, j}, L_{1} I \cap V_{1, m, c_{1}}, L_{1} I \cap V_{1, m, c_{2}}\right),
\end{array}
$$

where $(k, l)$ is the lowest block in $\bar{S}$. We obtain an induced morphism $Z_{\bar{Y}} \rightarrow \mathcal{Y}_{\bar{Y}}$, where $\mathcal{Y}_{\bar{Y}}=$ $Y_{\bar{S}}^{S_{1}, S_{2}}$.
Case 4.3.4 Assume finally that $m \not \equiv 1 \bmod n-2$ with vertex types A1 or A2. This is Case 4.2, 4: the full block at position $(0, m)$ has label $j$ and index set $S$ which is maximal; the full block at position $(1, m)$ has label $c$ and index set $S_{c} ; S$ is allowed by $S_{c}$. We get a morphism

$$
\begin{array}{ccc}
Z_{Y} & \rightarrow & V_{S} \times V_{S_{c}} \\
I & \mapsto & \left(I \cap V_{0, m, j}, I \cap V_{1, m, c}\right)
\end{array}
$$

whose image is contained in $X_{S}^{S_{c}} \subseteq V_{S} \times V_{S_{c}}$, an incidence variety we denote by $\mathcal{X}_{Y}$ to obtain an induced morphism $Z_{Y} \rightarrow \mathcal{X}_{Y, \bar{Y}}$.

Second, let

$$
\begin{array}{ccc}
Z_{\bar{Y}} & \rightarrow & V_{\bar{S}} \times V_{S_{c}} \\
I & \mapsto & \left(L_{1} I \cap V_{k, l, j}, L_{1} I \cap V_{1, m, c}\right),
\end{array}
$$

where $(k, l)$ is the lowest block in $\bar{S}$. By letting $\mathcal{Y}_{\bar{Y}}=Y_{\bar{S}}^{S_{c}}$ we obtain an induced morphism $Z_{\bar{Y}} \rightarrow \mathcal{Y}_{Y, \bar{Y}}$.

The last key step in our inductive proof is the following result, valid in all four cases above.
Proposition 4.17. The following is a scheme-theoretic fiber product diagram, with the right hand vertical map in each case given by the map induced by statement (1) of Propositions 4.10, 4.11
4.12 or 4.13 as appropriate.


Proof. It is immediate from the definitions in the different cases that the diagram is commutative. We thus need to show that it is a fibre product. Let $B$ be an arbitrary base scheme and let $f: B \rightarrow Z_{\bar{Y}}$ and $g: B \rightarrow \mathcal{X}_{Y}$ be morphisms; we need to show that these induce a unique morphism $B \rightarrow Z_{Y}$. We consider Case 4.31 ; the proof in the other cases is analogous. The map $f$ corresponds to a flat family of ideals $\mathcal{I}_{f} \triangleleft \mathcal{O}_{B} \otimes \mathbb{C}[x, y]$ with Young wall $\bar{Y}$. The map $g$ corresponds to a 3 -tuple of families $\mathcal{U}_{1, g}, \mathcal{U}_{2, g}, \mathcal{U}_{c, g}$ of subspaces of $\mathbb{C}[x, y]$ over $B$. Given this data, consider the family of ideals

$$
\mathcal{I}_{f, g}=\left(L_{1} \mathcal{I}_{f}, \mathcal{U}_{1, g}, \mathcal{U}_{2, g}\right) \triangleleft \mathcal{O}_{B} \otimes \mathbb{C}[x, y]
$$

over $B$. Here the parentheses mean the generated ideal as explained in 4.4. By the compatibility of $(f, g)$ it is immediate that the Young wall of the corresponding ideals is $Y$. The classifying map of this family is the unique possible extension of $(f, g)$ to a morphism $B \rightarrow Z_{Y}$.

Conclusion of the Proof of Theorem 4.3. Assume that we have shown for any Young wall $Y$ having less then $l$ rows that the corresponding stratum $Z_{Y}$ is affine, the $l=1$ case being obvious. Consider an arbitrary Young wall $Y$ consisting of $l$ rows. Let $\bar{Y}$ denote its truncation, as defined above. By the induction assumption, the space $Z_{\bar{Y}}$ is affine. Also, by Propositions 4.10, 4.11, 4.12 or 4.13 respectively, the map $\mathcal{X}_{Y} \rightarrow \mathcal{Y}_{\bar{Y}}$ of Proposition 4.17 is a trivial affine fibration in all cases. By Proposition 4.17, the map $Z_{Y} \rightarrow Z_{\bar{Y}}$ is a pullback of a trivial affine fibration and thus itself a trivial affine fibration. Using the induction hypothesis, $Z_{Y}$ is thus an affine space. The proof of Theorem 4.3 is complete.
Remark 4.18. One can deduce from the above proof that one can in fact canonically choose generators of a homogeneous ideal $I \in Z_{Y}$, which are in the cells of the some of the salient blocks of $Y$; as discussed before, not all salient cells necessarily contain a generator. For describing the coarse Hilbert scheme we have to keep track of these generators, but we will do this only implicitly.
Example 4.19. Returning to Examples 4.644.7, we see that for $Y_{3}$ the triangle of side $5, \overline{Y_{3}}=Y_{2}$, the triangle of size 4 . The map $T: Z_{Y_{3}} \cong \mathbb{A}^{2} \rightarrow Z_{\bar{Y}_{3}} \cong \mathbb{A}^{1}$ is the map identified at the end of the discussion of Example 4.7
4.4. Preparation for the proof of the incidence propositions. To prepare the ground for the proof of the propositions announced in 4.2, consider the operators defined in 3.4 We use these operators to describe projective coordinates on some of the Grassmannians $P_{m, j}$. We first record the following equalities, computing the isotypical summands of the homogeneous pieces of the ring $S=\mathbb{C}[x, y]$.

Lemma 4.20. We have

$$
\begin{gathered}
S_{2 k(n-2)}\left[\rho_{\kappa(0, k)}\right]=\left(L_{n-1}+L_{n}\right)^{2 k}\left[\rho_{\kappa(0, k)}\right]=\left(L_{n-1}^{2}+L_{n}^{2}\right)^{k} ; \\
S_{2 k(n-2)}\left[\rho_{\kappa(1, k)}\right]=\left(L_{n-1}+L_{n}\right)^{2 k}\left[\rho_{\kappa(1, k)}\right]=L_{n-1} L_{n}\left(L_{n-1}^{2}+L_{n}^{2}\right)^{k-1} ; \\
S_{(2 k+1)(n-2)}\left[\rho_{\kappa(n-1, k)}\right]=\left(L_{n-1}+L_{n}\right)^{2 k+1}\left[\rho_{\kappa(n-1, k)}\right]=L_{n-1}\left(L_{n-1}^{2}+L_{n}^{2}\right)^{k} ; \\
S_{(2 k+1)(n-2)}\left[\rho_{\kappa(n, k)}\right]=\left(L_{n-1}+L_{n}\right)^{2 k+1}\left[\rho_{\kappa(n, k)}\right]=L_{n}\left(L_{n-1}^{2}+L_{n}^{2}\right)^{k} .
\end{gathered}
$$

Proof. For each equality on the right, use (9) and an easy induction argument. For the equalities on the left, $\supseteq$ is always clear; then use dimension counting.

Thus, given an element $v \in P_{2 k(n-2), \kappa(0, k)}$, we can write it uniquely in the form

$$
v=\sum_{i=0}^{k} \alpha_{i} L_{n-1}^{2 i} L_{n}^{2(k-i)}
$$

for certain projective coordinates $\left[\alpha_{0}: \cdots: \alpha_{k}\right.$ ]. Analogous coordinates exist on $P_{2 k(n-2), \kappa(1, k)}$ and $P_{(2 k+1)(n-2), j}$ for $j=n-1, n$.

For the two-dimensional representations we similarly have

Lemma 4.21. For $1<j<n-1$,

$$
\begin{aligned}
S_{2 k(n-2)+j-1}\left[\rho_{j}\right] & =L_{j}\left(\left(L_{n-1}+L_{n}\right)^{2 k}\right) \\
& =L_{j}\left(\left(L_{n-1}+L_{n}\right)^{2 k}\left[\rho_{0}\right]\right) \oplus L_{j}\left(\left(L_{n-1}+L_{n}\right)^{2 k}\left[\rho_{1}\right]\right) ; \\
S_{(2 k+2)(n-2)-j+1}\left[\rho_{j}\right] & =L_{n-j}\left(\left(L_{n-1}+L_{n}\right)^{2 k+1}\right) \\
& =L_{n-j}\left(\left(L_{n-1}+L_{n}\right)^{2 k+1}\left[\rho_{n-1}\right]\right) \oplus L_{n-j}\left(\left(L_{n-1}+L_{n}\right)^{2 k+1}\left[\rho_{n}\right]\right) .
\end{aligned}
$$

Proof. Analogous.
For such $1<j<n-1$, the space $P_{2 k(n-2)+j-1, j}$ of two-dimensional $G_{\Delta}$-equivariant subspaces of $S_{2 k(n-2)+j-1}\left[\rho_{j}\right]$ is a projective space as noted above. Using Lemma 4.21, we get a collection of distinguished two-dimensional $G_{\Delta}$-equivariant subspaces $L_{j} L_{n-1}^{i} L_{n}^{2 k-i}$ in $S_{2 k(n-2)+j-1}\left[\rho_{j}\right]$; an arbitrary element $v \in P_{2 k(n-2)+j-1, j}$ can be uniquely written as

$$
v=\sum_{i=0}^{2 k} \alpha_{i} L_{j} L_{n-1}^{i} L_{n}^{2 k-i}
$$

for certain projective coordinates $\left[\alpha_{0}: \ldots: \alpha_{2 k}\right]$. Analogous coordinates also exist on the space $P_{(2 k+2)(n-2)-j+1, j}$.

For subspaces $U_{1}, \ldots, U_{i} \in \operatorname{Gr}(S)$ denote by $\left(U_{1}, \ldots, U_{i}\right)$ the $G$-invariant ideal of $S$ generated by the corresponding subspaces. In particular, the ideal generated by (the subspaces represented by) points $v_{1}, \ldots, v_{i} \in \mathbb{P} S$ is denoted by $\left(v_{1}, \ldots, v_{i}\right)$. Then $\left(U_{1}, \ldots, U_{i}\right)_{m, j}$ is represented by a projective linear subspace of $P_{m, j}$ (cf. Lemma 3.6). Thus, we can talk about its intersection with the cells of $P_{m, j}$. For simplicity we will use the notation $\left(U_{1}, \ldots, U_{i}\right) \cap V_{k, l, j}=\left(U_{1}, \ldots, U_{i}\right)_{k+l, j} \cap V_{k, l, j}$ for the intersection with $V_{k, l, j}$.

We need to study indidence relations between ideals generated by subspaces from the various strata defined above. First of all, let $v_{j} \in V_{0, m, j}$ for some $m$ and $j$, corresponding to a full or half block. Denote by $C$ the set of labels of full or half blocks immediately above or on the right of this block, i.e. in the positions $(1, m)$ or $(0, m+1)$. Then the definition of the McKay quiver clearly implies that we have $\left(v_{j}\right) \cap S_{m+1, c}=\emptyset$ whenever $c \notin C$. The following long statement discusses all the remaining cases when $c \in C$, split into the different possibilities.
Proposition 4.22. (1) For $j=0,1$, fix $v_{j} \in V_{0,2 k(n-2), j}$.
(a) We have $\left(v_{j}\right) \cap\left(\bigcup_{l \geq 1} V_{l, 2 k(n-2)-l+1,2}\right)=\emptyset$. Hence the unique point $\left(v_{j}\right) \cap P_{2 k(n-2)+1}$ necessarily lies in $V_{0,2 k(n-2)+1,2}$. This provides an injection

$$
V_{0,2 k(n-2), j} \rightarrow V_{0,2 k(n-2)+1,2}
$$

(b) $\left(v_{0}, v_{1}\right) \cap\left(\bigcup_{l>1} V_{l, 2 k(n-2)-l+1,2}\right)=\emptyset$. In particular, the projective line $\left(v_{0}, v_{1}\right) \cap$ $P_{2 k(n-2)+1,2}$ necessarily intersects $V_{1,2 k(n-2), 2}$. Let the intersection point be $L_{1} v_{2}$ for $a$ certain $v_{2} \in V_{0,2 k(n-2)-1,2}$. Then $v_{2}$ is the unique point of $V_{0,2 k(n-2)-1,2}$ such that $v_{0}, v_{1} \in\left(v_{2}\right)$. As a consequence, for any projective subspace $U_{2} \subseteq P_{2 k(n-2)-1,2}$, the intersection $\left(v_{0}, v_{1}\right) \cap\left(\bigcup_{l>0} V_{l, 2 k(n-2)-l+1,2}\right)$ is contained in $L_{1} U_{2}$ if and only if $v_{0}, v_{1} \in\left(U_{2}\right)$.
(2) Let $v_{2} \in V_{0,2 k(n-2)+1,2}$. For $j=0$, 1 , if $\left(v_{2}\right) \cap\left(\bigcup_{l>0} V_{l, 2 k(n-2)-l+2, j}\right)$ is not empty, then it is necessarily one-dimensional, and it equals $L_{1} v_{j}$ for a certain $v_{j} \in V_{0,2 k(n-2), j}$. Exactly one of the following three cases happens.

- $\left(v_{2}\right) \cap\left(\bigcup_{l>0} V_{l, 2 k(n-2)-l+2,1}\right)=L_{1} v_{0}$ and $\left(v_{2}\right) \cap\left(\bigcup_{l>0} V_{l, 2 k(n-2)-l+2,0}\right)=\emptyset$. This happens if and only if $v_{2} \in\left(v_{0}\right)$. In this case, $v_{0} \in V_{0,2 k(n-2), 0}$, and $\left(v_{2}\right) \cap S_{2 k(n-2)+2}$ has two (resp. three if $n=4$ ) irreducible components: $L_{1} v_{0}$ and $\left(v_{2}\right) \cap V_{0,2 k(n-2)+2,3}$ (resp. $\left(v_{2}\right) \cap V_{0,2 k(n-2)+2,3}$ and $\left.\left(v_{2}\right) \cap V_{0,2 k(n-2)+2,4}\right)$.
- $\left(v_{2}\right) \cap\left(\bigcup_{l>0} V_{l, 2 k(n-2)-l+2,1}\right)=\emptyset$ and $\left(v_{2}\right) \cap\left(\bigcup_{l>0} V_{l, 2 k(n-2)-l+2,0}\right)=L_{1} v_{1}$ with symmetrical statements as in the previous case.
- $\left(v_{2}\right) \cap\left(\bigcup_{l>0} V_{l, 2 k(n-2)-l+2,1}\right)=L_{1} v_{0}$ and $\left(v_{2}\right) \cap\left(\bigcup_{l>0} V_{l, 2 k(n-2)-l+2,0}\right)=L_{1} v_{1}$. This happens if and only if $v_{2} \in\left(v_{0}, v_{1}\right)$ but $v_{2} \notin\left(v_{0}\right) \cup\left(v_{1}\right)$. In this case at least one of the inclusions $v_{0} \in V_{0,2 k(n-2), 0}, v_{1} \in V_{0,2 k(n-2), 1}$ is satisfied but not necessarily both. Furthermore, $\left(v_{2}\right) \cap S_{2 k(n-2)+2}$ has three (resp. four if $n=4$ ) irreducible
components: $L_{1} v_{0}, L_{1} v_{1}$ and $\left(v_{2}\right) \cap V_{0,2 k(n-2)+2,3} \quad$ (resp. $\left(v_{2}\right) \cap V_{0,2 k(n-2)+2,3}$ and $\left.\left(v_{2}\right) \cap V_{0,2 k(n-2)+2,4}\right)$.
Thus for $n>4$, we obtain an isomorphism

$$
\begin{aligned}
V_{0,2 k(n-2)+1,2} & \rightarrow V_{0,2 k(n-2)+2,3} \\
v_{2} & \mapsto\left(v_{2}\right) \cap V_{0,2 k(n-2)+2,3}
\end{aligned}
$$

whereas for $n=4$, we get an isomorphism

$$
\begin{aligned}
V_{0,2 k(n-2)+1,2} & \rightarrow V_{0,2 k(n-2)+2,3} \times V_{0,2 k(n-2)+2,4} \\
v_{2} & \mapsto\left(\left(v_{2}\right) \cap V_{0,2 k(n-2)+2,3},\left(v_{2}\right) \cap V_{0,2 k(n-2)+2,4}\right) .
\end{aligned}
$$

Finally, for projective subspaces $U_{0} \subseteq P_{2 k(n-2), 0}, U_{1} \subseteq P_{2 k(n-2), 1}$,

- the conditions $\left(v_{2}\right) \cap\left(\bigcup_{l>0} V_{l, 2 k(n-2)-l+2,1}\right) \subseteq L_{1} U_{0}$ and $\quad\left(v_{2}\right) \cap$ $\left(\bigcup_{l>0} V_{l, 2 k(n-2)-l+2,0}\right)=\emptyset$ are satisfied if and only if $v_{2} \in\left(U_{0}\right)$;
- the conditions $\left(v_{2}\right) \cap\left(\bigcup_{l>0} V_{l, 2 k(n-2)-l+2,1}\right)=\emptyset$ and $\left(v_{2}\right) \cap\left(\bigcup_{l>0} V_{l, 2 k(n-2)-l+2,0}\right) \subseteq$ $L_{1} U_{1}$ are satisfied if and only if $v_{2} \in\left(U_{1}\right)$;
- the conditions $\left(v_{2}\right) \cap\left(\bigcup_{l>0} V_{l, 2 k(n-2)-l+2,1}\right) \subseteq L_{1} U_{0}$ and $\left(v_{2}\right) \cap$ $\left(\bigcup_{l>0} V_{l, 2 k(n-2)-l+2,0}\right) \subseteq L_{1} U_{1}$ are satisfied if and only if $v_{2} \in\left(U_{0}, U_{1}\right)$ but $v_{2} \notin\left(U_{0}\right) \cup\left(U_{1}\right)$.
(3) Assume that $3 \leq j \leq n-3$ (resp. $j=n-2$ ) and set $v_{j} \in V_{0,2 k(n-2)+j-1, j}$. Then $\left(v_{j}\right) \cap\left(\bigcup_{l>1} V_{l, 2 k(n-2)-l+j, j-1}\right)=\emptyset$. Furthermore, $\left(v_{j}\right) \cap S_{2 k(n-2)+j}$ has two (resp. three) irreducible components: a point $\left(v_{j}\right) \cap V_{1,2 k(n-2)+j-1, j-1}$ of the form $L_{1} v_{j-1}$ for some $v_{j-1} \in V_{0,2 k(n-2)+j-2, j-1}$, which is the unique element with $v_{j} \in\left(v_{j-1}\right)$, and another point $\left(v_{j}\right) \cap V_{0,2 k(n-2)+j, j+1}$ (resp. two other points $\left.\left(v_{j}\right) \cap V_{0,2 k(n-2)+j, n-1},\left(v_{j}\right) \cap V_{0,2 k(n-2)+j, n}\right)$. We obtain an isomorphism

$$
V_{0,2 k(n-2)+j-1, j} \rightarrow V_{0,2 k(n-2)+j, j+1}
$$

for $j \leq n-3$ and an isomorphism

$$
V_{0,2 k(n-2)+j-1, j} \rightarrow V_{0,2 k(n-2)+j, n-1} \times V_{0,2 k(n-2)+j, n}
$$

for $j=n-2$. Also, for a projective subspace $U_{j-1} \subseteq P_{2 k(n-2)+j-2, j-1}$, the intersection $\left(v_{j}\right) \cap\left(\bigcup_{l>0} V_{l, 2 k(n-2)-l+j, j-1}\right)$ is contained in $L_{1} U_{j-1}$ if and only if $v_{j} \in\left(U_{j-1}\right)$.
(4) For $j=n-1, n$, fix $v_{j} \in V_{0,(2 k+1)(n-2), j}$.
(a) We have $\left(v_{j}\right) \cap\left(\bigcup_{l \geq 1} V_{l,(2 k+1)(n-2)-l+1, n-2}\right)=\emptyset$. Hence, the point $\left(v_{j}\right) \cap$ $P_{(2 k+1)(n-2)+1}$ is necessarily in $V_{0,(2 k+1)(n-2)+1, n-2}$. This provides an injection

$$
V_{0,(2 k+1)(n-2), j} \rightarrow V_{0,(2 k+1)(n-2)+1, n-2}
$$

(b) $\left(v_{n-1}, v_{n}\right) \cap\left(\bigcup_{l>1} V_{l,(2 k+1)(n-2)-l+1,2}\right)=\emptyset$. In particular, the projective line $\left(v_{n-1}, v_{n}\right) \cap P_{(2 k+1)(n-2)+1, n-2}$ necessarily intersects $V_{1,(2 k+1)(n-2), n-2}$. Let the intersection point be $L_{1} v_{n-2}$ for a certain $v_{n-2} \in V_{0,(2 k+1)(n-2)-1, n-2}$. Then $v_{n-2}$ is the unique point of $V_{0,(2 k+1)(n-2)-1, n-2}$ such that $v_{n-1}, v_{n} \in\left(v_{n-2}\right)$. As a consequence, for any projective subspace $U_{n-2} \subseteq P_{(2 k+1)(n-2)-1, n-2}$, the intersection $\left(v_{n-1}, v_{n}\right) \cap\left(\bigcup_{l>0} V_{l,(2 k+1)(n-2)-l+1, n-2}\right)$ is contained in $L_{1} U_{n-2}$ if and only if $v_{n-1}, v_{n} \in\left(U_{n-2}\right)$.
(5) Let $v_{n-2} \in V_{0,(2 k+1)(n-2)+1, n-2}$. For $j=n-1, n$ if $\left(v_{n-1}\right) \cap\left(\bigcup_{l>0} V_{l,(2 k+1)(n-2)-l+2, j}\right) \neq \emptyset$ then this invariant subspace is one-dimensional, and it is of the form $L_{1} v_{j}$ for certain $v_{j} \in V_{0,(2 k+1)(n-2), j}$. Exactly one of the following three possibilities happens.

- $\left(v_{n-2}\right) \cap\left(\bigcup_{l>0} V_{l,(2 k+1)(n-2)-l+2, n}\right) \quad=\quad L_{1} v_{n-1}$ and $\left(v_{n-2}\right) \cap$ $\left(\bigcup_{l>0} V_{l,(2 k+1)(n-2)-l+2, n-1}\right)=\emptyset$. This happens if and only if $v_{n-2} \in\left(v_{n-1}\right)$. In this case $v_{n-1} \in V_{0,(2 k+1)(n-2), n-1}$, and $\left(v_{2}\right) \cap S_{(2 k+1)(n-2)+2}$ has two (resp. three if $n=4$ ) irreducible components: $L_{1} v_{n-1}$ and $\left(v_{n-2}\right) \cap V_{0,(2 k+1)(n-2)+2, n-3}$ (resp. $\left(v_{2}\right) \cap V_{0,(2 k+1)(n-2)+2,0}$ and $\left.\left(v_{2}\right) \cap V_{0,(2 k+1)(n-2)+2,1}\right)$.
- $\left(v_{n-2}\right) \cap\left(\bigcup_{l>0} V_{l,(2 k+1)(n-2)-l+2, n}\right)=\emptyset$ and $\left(v_{n-2}\right) \cap\left(\bigcup_{l>0} V_{l, 2 k(n-2)-l+2, n-1}\right)=L_{1} v_{n}$ with symmetrical statements as in the previous case.
- $\left(v_{n-2}\right) \cap\left(\bigcup_{l>0} V_{l,(2 k+1)(n-2)-l+2, n}\right) \quad=\quad L_{1} v_{n-1}$ and $\quad\left(v_{n-2}\right) \cap$ $\left(\bigcup_{l>0} V_{l,(2 k+1)(n-2)-l+2, n-1}\right)=L_{1} v_{n}$. This happens if and only if $v_{n-2} \in\left(v_{n-1}, v_{n}\right)$ but $v_{n-2} \notin\left(v_{n-1}\right) \cup\left(v_{n}\right)$. In this case at least one of the inclusions
$v_{n-1} \in V_{0,(2 k+1)(n-2), n-1}, \quad v_{n} \in V_{0,(2 k+1)(n-2), n}$ is satisfied but not necessarily both. Furthermore, $\left(v_{n-2}\right) \cap S_{(2 k+1)(n-2)+2}$ has three (resp. four if $n=4)$ irreducible components: $L_{1} v_{n-1}, L_{1} v_{n}$ and $\left(v_{n-2}\right) \cap V_{0,(2 k+1)(n-2)+2, n-3}$ (resp. $\left(v_{n-2}\right) \cap V_{0,(2 k+1)(n-2)+2,0}$ and $\left.\left(v_{n-2}\right) \cap V_{0,(2 k+1)(n-2)+2,1}\right)$.
For $n>4$, we obtain isomorphisms

$$
\begin{aligned}
V_{0,(2 k+1)(n-2)+1, n-2} & \rightarrow V_{0,(2 k+1)(n-2)+2, n-3} \\
v_{n-2} & \mapsto\left(v_{n-2}\right) \cap V_{0,(2 k+1)(n-2)+2, n-3}
\end{aligned},
$$

whereas for $n=4$ we obtain an isomorphism

$$
\begin{aligned}
V_{0,(2 k+1)(n-2)+1, n-2} & \rightarrow V_{0,2 k(n-2)+2,0} \times V_{0,(2 k+1)(n-2)+2,1} \\
v_{n-2} & \mapsto\left(\left(v_{n-2}\right) \cap V_{0,(2 k+1)(n-2)+2,0},\left(v_{n-2}\right) \cap V_{0,(2 k+1)(n-2)+2,1}\right) .
\end{aligned}
$$

Moreover, for projective subspaces $U_{n-1} \subseteq P_{(2 k+1)(n-2), n-1}, U_{n} \subseteq P_{(2 k+1)(n-2), n}$ the conditions

- $\left(v_{n-2}\right) \cap\left(\bigcup_{l>0} V_{l,(2 k+1)(n-2)-l+2, n}\right) \quad \subseteq \quad L_{1} U_{n-1} \quad$ and $\quad\left(v_{n-2}\right) \quad \cap$ $\left(\bigcup_{l>0} V_{l,(2 k+1)(n-2)-l+2, n-1}\right)=\emptyset$ are satisfied if and only if $v_{n-2} \in\left(U_{n-1}\right)$;
- $\left(v_{n-2}\right) \cap\left(\bigcup_{l>0} V_{l,(2 k+1)(n-2)-l+2, n}\right)=\emptyset$ and $\left(v_{n-2}\right) \cap\left(\bigcup_{l>0} V_{l,(2 k+1)(n-2)-l+2, n-1}\right) \subseteq$ $L_{1} U_{n}$ are satisfied if and only if $v_{n-2} \in\left(U_{n}\right)$;
- $\left(v_{n-2}\right) \cap\left(\bigcup_{l>0} V_{l,(2 k+1)(n-2)-l+2, n}\right) \quad \subseteq \quad L_{1} U_{n-1}$ and $\quad\left(v_{n-2}\right) \cap$ $\left(\bigcup_{l>0} V_{l,(2 k+1)(n-2)-l+2, n-1}\right) \subseteq L_{1} U_{n}$ are satisfied if and only if $v_{n-2} \in\left(U_{n-1}, U_{n}\right)$ but $v_{2} \notin\left(U_{n-1}\right) \cup\left(U_{n}\right)$.
(6) Assume that $3 \leq j \leq n-3$ (resp. $j=2$ ) and set $v_{j} \in V_{0,(2 k+2)(n-2)-j+1, j}$. Then $\left(v_{j}\right) \cap\left(\bigcup_{l>1} V_{l,(2 k+2)(n-2)-l-j+2, j+1}\right)=\emptyset$. Furthermore, $\left(v_{j}\right) \cap S_{(2 k+2)(n-2)-j+2}$ has two (resp. three) irreducible components: a point $\left(v_{j}\right) \cap V_{1,(2 k+2)(n-2)-j+1, j+1}$ of the form $L_{1} v_{j+1}$ for some $v_{j+1} \in V_{0,(2 k+2)(n-2)-j, j+1}$, which is the unique element with $v_{j} \in\left(v_{j+1}\right)$, and another point $\left(v_{j}\right) \cap V_{0,(2 k+2)(n-2)-j+2, j-1}$ (resp. two other points $\left.\left(v_{j}\right) \cap V_{0,(2 k+2)(n-2)-j+2,0,},\left(v_{j}\right) \cap V_{0,(2 k+2)(n-2)-j+2,1)}\right)$ providing an isomorphism $V_{0,(2 k+2)(n-2)-j+1, j} \rightarrow V_{0,2 k(n-2)-j+2, j-1} \quad$ resp. $V_{0,(2 k+2)(n-2)-j+1, j} \rightarrow V_{0,2 k(n-2)-j+2,0} \times$ $\left.V_{0,2 k(n-2)-j+2,1}\right)$. As a consequence, for a projective subspace $U_{j+1} \subseteq P_{(2 k+2)(n-2)-j, j+1}$, the intersection $\left(v_{j}\right) \cap\left(\bigcup_{l>0} V_{l,(2 k+2)(n-2)-l-j+2, j+1}\right)$ is contained in $L_{1} U_{j+1}$ if and only if $v_{j} \in\left(U_{j+1}\right)$.

Proof. For ease of notation in the proof, we will assume that $n$ is even. If $n$ is odd, then the argument works in the same way, except that $\kappa$ should be applied to the indices when appropriate.
(1) Statement (a) follows immediately from the definition of the cells. For the first part of (b), write $v_{0}=\sum_{i=0}^{k} \alpha_{i} L_{n-1}^{2 i} L_{n}^{2(k-i)}$ and $v_{1}=\sum_{i=0}^{k-1} \beta_{i} L_{n-1}^{2 i+1} L_{n}^{2(k-i)-1}$. The assumptions guarantee that $\sum \alpha_{i} \neq 0$ and $\sum \beta_{i} \neq 0$. The image $\left(v_{0}, v_{1}\right) \cap P_{2 k(n-2)+1,2}$, as subset of $G r\left(2, S_{2 k(n-2)+1}\left[\rho_{2}\right]\right)$, is a projective line and is spanned by

$$
L_{2} v_{0}=\sum_{i=0}^{k} \alpha_{i} L_{2} L_{n-1}^{2 i} L_{n}^{2(k-i)} \in L_{2}\left(L_{n-1}^{2}+L_{n}^{2}\right)^{k}
$$

and

$$
L_{2} v_{1}=\sum_{i=0}^{k-1} \beta_{i} L_{2} L_{n-1}^{2 i+1} L_{n}^{2(k-i)-1} \in L_{2} L_{n-1} L_{n}\left(L_{n-1}^{2}+L_{n}^{2}\right)^{k-1}
$$

In particular, if $\left(v_{0}, v_{1}\right) \cap V_{1,2 k(n-2)}=\left\{L_{1} v_{2}\right\}$, then there exist vectors $v_{x}, v_{y}$ in the two-dimensional vector space $v_{2}$ satisfying $v_{y}=\tau\left(v_{x}\right)$, as well as $a_{x}, b_{x}, a_{y}, b_{y} \in \mathbb{C}$, so that

$$
\begin{align*}
& L_{1} v_{x}=a_{x} L_{2, x} v_{0}+b_{x} L_{2, x} v_{1} \\
& L_{1} v_{y}=a_{y} L_{2, y} v_{0}+b_{y} L_{2, y} v_{1} \tag{11}
\end{align*}
$$

In $v_{0}$ the highest powers of $x$ and $y$ are $x^{2 k(n-2)}$ and $y^{2 k(n-2)}$, both with coefficient $\sum_{i} \alpha_{i}$. In $v_{1}$ the highest powers of $x$ and $y$ are $x^{2 k(n-2)}$ and $y^{2 k(n-2)}$, the first has coefficient $\sum_{i} \beta_{i}$, the second has coefficient $-\sum_{i} \beta_{i}$. We apply $L_{2, x}$ to these. In order for the sum to avoid the cell $V_{0,2 k(n-2)+1}$ the coefficients must satisfy $\left[a_{x}: b_{x}\right]=\left[\sum_{i} \beta_{i}:-\sum_{i} \alpha_{i}\right]$, since in this case the coefficient of $x^{2 k(n-2)+1}$ is 0 . Then the linear combination is necessarily in $V_{1,2 k(n-2)}$. Similarly, when applying $L_{2, y}$, the required coefficients are $\left[a_{y}: b_{y}\right]=\left[\sum_{i} \beta_{i}: \sum_{i} \alpha_{i}\right]$, so we have $a_{x} b_{y}=-a_{y} b_{x}$. Therefore, $b_{y} L_{2, y} L_{1} v_{x}-b_{x} L_{2, x} L_{1} v_{y}=b_{y} a_{x} L_{2, y} L_{2, x} v_{0}-b_{x} a_{y} L_{2, x} L_{2, y} v_{0}=L_{1} v_{0}$, where we have used that
$L_{2, x} L_{2, y}=L_{1}$. As a consequence, $b_{y} L_{2, y} v_{x}+b_{x} L_{2, x} v_{y}=v_{0}$ and similarly $-a_{y} L_{2, y} v_{x}+a_{x} L_{2, x} v_{y}=$ $v_{1}$, i.e. $v_{0}, v_{1} \in\left(v_{2}\right)$. This proves the first part of (b). The second part of (b), concerning projective subspaces, follows immediately from the first part.
(2) For the first part of the statement let $v_{2}=\sum_{i=0}^{2 k} \varepsilon_{i} L_{2} L_{n-1}^{i} L_{n}^{2 k-i}=$ $L_{2}\left(\sum_{i=0, i \text { even }}^{2 k} \varepsilon_{i} L_{n-1}^{i} L_{n}^{2 k-i}+\sum_{i=0, i \text { odd }}^{2 k} \varepsilon_{i} L_{n-1}^{i} L_{n}^{2 k-i}\right)=L_{2}\left(v_{\text {even }}+v_{\text {odd }}\right)$. If $n \neq 4$, then by applying $L_{2}$ again we get $L_{2}^{2}\left(v_{\text {even }}+v_{\text {odd }}\right)=\left(L_{1}+L_{3}\right)\left(v_{\text {even }}+v_{\text {odd }}\right)$, where the first sum is operator sum, and the second is vector sum. So $L_{2} v_{2}\left[\rho_{0}\right]=\left\{L_{1} v_{\text {even }}\right\}$ and $L_{2} v_{2}\left[\rho_{1}\right]=\left\{L_{1} v_{\text {odd }}\right\}$. If $\left(v_{2}\right) \cap\left(\bigcup_{l>1} V_{l, 2 k(n-2)-l+2,1}\right)=L_{1} v_{0}$ and $\left(v_{2}\right) \cap\left(\bigcup_{l>1} V_{l, 2 k(n-2)-l+2,0}\right)=\emptyset$, then $v_{\text {even }}=v_{0}$ and $v_{\text {odd }}=0$. The second case is just the opposite, and when $\left(v_{2}\right) \cap\left(\bigcup_{l>1} V_{l, 2 k(n-2)-l+2,1}\right)=L_{1} v_{0}$ and $\left(v_{2}\right) \cap\left(\bigcup_{l>1} V_{l, 2 k(n-2)-l+2,0}\right)=L_{1} v_{1}$, then $v_{\text {even }}=a v_{0}$ and $v_{\text {odd }}=b v_{1}$ for some coefficients $a, b \in \mathbb{C}$. If $n=4$, then $L_{2}^{2}\left(v_{\text {even }}+v_{\text {odd }}\right)=\left(L_{1}+L_{3}+L_{4}\right)\left(v_{\text {even }}+v_{\text {odd }}\right)$, and the rest is very similar to the first case.

For the second part of the statement observe that the results so far imply that $V_{0,2 k(n-2)+1,2}$ stratifies into disjoint, locally closed subspaces

$$
\begin{aligned}
V_{0,2 k(n-2)+1,2}= & \bigsqcup_{\left(v_{c_{1}}, v_{c_{2}}\right) \in P_{2 k(n-2), c_{1}} \times P_{2 k(n-2), c_{2}}}\left(\left(v_{c_{1}}, v_{c_{2}}\right) \backslash\left(\left(v_{c_{1}}\right) \cup\left(v_{c_{2}}\right)\right) \cap V_{0,2 k(n-2)+1,2}\right) \\
& \bigsqcup_{v_{c_{2}} \in V_{0,2 k(n-2), c_{2}}} \bigsqcup_{v_{c_{1}} \in V_{0,2 k(n-2), c_{1}}}\left(\left(v_{c_{2}}\right) \cap V_{0,2 k(n-2)+1,2}\right) \\
& \left.\left.\bigsqcup_{c_{1}}\right) \cap V_{0,2 k(n-2)+1,2}\right)
\end{aligned}
$$

The subset $\cup_{\left(v_{c_{1}}, v_{c_{2}}\right) \in V_{0,2 k(n-2), c_{1}} \times V_{0,2 k(n-2), c_{2}}}\left(\left(v_{c_{1}}, v_{c_{2}}\right) \backslash\left(\left(v_{c_{1}}\right) \cup\left(v_{c_{2}}\right)\right) \cap V_{0,2 k(n-2)+1,2}\right)$ is dense in the third stratum, since $V_{0,2 k(n-2), c_{1}} \times V_{0,2 k(n-2), c_{2}}$ is dense in $P_{2 k(n-2), c_{1}} \times P_{2 k(n-2), c_{2}}$. The first statement of (2) implies that the second statement is valid if $v_{2}$ is in this subset of $V_{0,2 k(n-2)+1,2}$. Similarly, the first statement implies the second statement on the loci $\sqcup_{v_{c_{1}} \in V_{0,2 k(n-2), c_{1}}}\left(\left(v_{c_{1}}\right) \cap V_{0,2 k(n-2)+1,2}\right)$ and $\sqcup_{v_{c_{2}} \in V_{0,2 k(n-2), c_{2}}}\left(\left(v_{c_{2}}\right) \cap V_{0,2 k(n-2)+1,2}\right)$. For $v_{2}$ in the closed complement of the union of these loci, the second statement follows from the linearity (and thus continuity) of the solution of (11), since the $U_{i}$ are projective.
(3) The statements in this case follow similarly to (2) by observing that $L_{2} L_{j}=L_{1} L_{j-1}+L_{j+1}$ $\left(\right.$ resp. $\left.L_{2} L_{n-2}=L_{1} L_{n-3}+L_{n-1}+L_{n}\right)$.

The cases (4), (5) and (6) are analogous.
Consider a full block in position $(0, m)$ with $m=k(n-2)+1$, with label $j$ which is 2 or $n-2$. In positions $(0, k(n-2))$ and $(1, m)$, above and to the left of this full block, are divided blocks with labels $\left(c_{1}, c_{2}\right)$, either $(0,1)$ or $(n-1, n)$. The next lemma gives $P_{m, j}$ the structure of a join (see B in the Appendix) of two projective subspaces.
Lemma 4.23. (1) The morphisms $V_{0, k(n-2), c_{i}} \rightarrow V_{0, m, j}$ constructed in Proposition 4.22 extend to morphisms

$$
\begin{array}{cc}
\phi_{i}: P_{k(n-2), c_{i}} & \rightarrow P_{m, j} \\
v & \mapsto(v) \cap P_{m, j} .
\end{array}
$$

The morphism $\phi_{i}$ is injective with image $N_{c_{i}}^{0}:=\operatorname{im}\left(\phi_{i}\right) \subset P_{m, j}$ such that $N_{c_{1}}^{0}, N_{c_{2}}^{0}$ are disjoint projective linear subspaces of $P_{m, j}$.
(2) The join of the disjoint linear subspaces $N_{c_{1}}^{0}, N_{c_{1}}^{0} \subset P_{m, j}$ is $P_{m, j}$ itself. Thus given $\left(v_{1}, v_{2}\right) \in P_{k(n-2), c_{1}} \times P_{k(n-2), c_{2}}$, there is a projective line $\mathbb{P}^{1} \cong v_{1} v_{2} \subset P_{m, j}$ containing both $\phi_{i}\left(v_{i}\right)$, namely, the line defined by $v_{1}, v_{2}$ with endpoints $\phi_{i}\left(v_{i}\right)$. The lines $v_{1} v_{2}$ cover $P_{m, j}$. For all $\left(v_{1}, v_{2}\right),\left(v_{1}^{\prime}, v_{2}^{\prime}\right) \in P_{k(n-2), c_{1}} \times P_{k(n-2), c_{2}}$, the intersection $v_{1} v_{2} \cap v_{1}^{\prime} v_{2}^{\prime}$ can only be at a common endpoint.
(3) For all $\left(v_{1}, v_{2}\right) \in P_{k(n-2), c_{1}} \times P_{k(n-2), c_{2}}$, the intersection $v_{1} v_{2} \cap V_{0, m, j}$ is

- either empty, exactly when $v_{1} \notin V_{0, k(n-2), c_{1}}$ and $v_{2} \notin V_{0, k(n-2), c_{2}}$;
- or an affine line otherwise.

Proof. (1) is immediate. (2) then follows from $\operatorname{dim} P_{k(n-2), c_{1}}+\operatorname{dim} P_{k(n-2), c_{i}}+1=\operatorname{dim} P_{m, j}$ and Lemma B.1. (3) is again immediate.

As we did already in the statement above, we will sometimes omit the inclusion maps $\phi_{i}$; thus, for subspaces $U_{1} \subseteq P_{k(n-2), c_{1}}$ and $U_{2} \subseteq P_{k(n-2), c_{2}}$, we will denote by $J\left(U_{1}, U_{2}\right) \subseteq P_{m, j}$ the join of $\phi_{1}\left(U_{1}\right)$ and $\phi_{1}\left(U_{1}\right)$ in $P_{m, j}$.

Let

$$
M_{c_{i}}^{0}:=\phi_{i}\left(V_{0, k(n-2), c_{i}}\right) \subset V_{0, m, j}
$$

these are disjoint affine linear subspaces of the affine space $V_{0, m, j}$. Also consider

$$
N_{c_{i}}=J\left(P_{k(n-2), c_{i}}, \bar{P}_{k(n-2), c_{3-i}}\right) \subset P_{m, j}
$$

This is the locus of points in $P_{m, j}$ covered by lines $v_{1} v_{2}$ one of whose endpoints is at a point "at infinity", in $\bar{P}_{k(n-2), c_{3-i}}=P_{k(n-2), c_{3-i}} \backslash V_{0, k(n-2), c_{3-i}}$. Let

$$
M_{c_{i}}=N_{c_{i}} \cap V_{0, m, j} \subset V_{0, m, j}
$$

be the intersection with the large affine cell of $P_{m, j}$.
Lemma 4.24. There exists morphisms $\psi_{i}: M_{c_{i}} \rightarrow M_{c_{i}}^{0}$, given by associating to a point

$$
v \in M_{c_{i}} \subset V_{0, m, j}
$$

the "non-infinity" endpoint of the (mostly unique) line $v_{1} v_{2}$ passing through it. The maps $\psi_{i}$ are trivial vector bundles over affine spaces.

Proof. See Lemma B. 2 ,
Corollary 4.25. (1) For $i=1,2$ the decomposition $P_{k(n-2), c_{3-i}}=$ $\bigsqcup_{\left(k^{\prime}, l^{\prime}\right) \in B_{k(n-2), c_{3-i}}} V_{k^{\prime}, l^{\prime}, c_{3-i}}$ induces a decomposition into locally closed subspaces

$$
\begin{equation*}
M_{c_{i}} \backslash M_{c_{i}}^{0}=\bigsqcup_{\left(k^{\prime}, l^{\prime}\right) \in B_{k(n-2), c_{3-i}} \backslash(1, m)}\left(\left(J\left(V_{0, k(n-2), c_{i}}, V_{k^{\prime}, l^{\prime}, c_{3-i}}\right) \cap V_{0, m, j}\right) \backslash M_{c_{i}}^{0}\right) \tag{12}
\end{equation*}
$$

(2) Taking into account the bijections $B_{k(n-2), c_{i}} \cong B_{k(n-2)+2, c_{3-i}}$ and the decomposition (12), the space $V_{0, m, j}$ decomposes into locally closed subspaces as

$$
\begin{aligned}
V_{0, m, j}= & V_{0, m, j}(1, m, 1, m) \\
& \bigsqcup\left(\begin{array}{c}
\hbar_{\left(k_{1}, l_{1}\right) \in B_{k(n-2)+2, c_{1}}^{\prime}} V_{0, m, j}\left(k_{1}, l_{1}, 1, m\right) \\
\\
\\
\\
\vdots\left(k_{2}, l_{2}\right) \in B_{k(n-2)+2, c_{2}}^{\prime}
\end{array} V_{0, m, j}\left(1, m, k_{2}, l_{2}\right)\right),
\end{aligned}
$$

where we introduced the notations

- $B_{k(n-2)+2, c_{i}}^{\prime}=\left(B_{k(n-2)+2, c_{i}} \cup\{\emptyset\}\right) \backslash\{(1, m)\} ;$
- $V_{0, m, j}(\emptyset, 1, m)=M_{c_{1}}^{0}$;
- $V_{0, m, j}\left(k_{1}, l_{1}, 1, m\right)=\left(J\left(V_{0, k(n-2), c_{1}}, V_{k_{1}, l_{1}, c_{2}}\right) \cap V_{0, m, j}\right) \backslash M_{c_{1}}^{0}$;
- $V_{0, m, j}(1, m, \emptyset)=M_{c_{2}}^{0}$;
- $V_{0, m, j}\left(1, m, k_{2}, l_{2}\right)=\left(J\left(V_{0, k(n-2), c_{2}}, V_{k_{2}, l_{2}, c_{1}}\right) \cap V_{0, m, j}\right) \backslash M_{c_{2}}^{0}$;
- $V_{0, m, j}(1, m, 1, m)=V_{0, m, j} \backslash\left(M_{c_{1}} \cup M_{c_{2}}\right)$.

The meaning of the notations is that $V_{0, m, j}\left(k_{1}, l_{1}, k_{2}, l_{2}\right)$ consists of exactly those points $v \in$ $V_{0, m, j}$ such that $(v) \cap P_{k(n-2)+2, c_{i}} \in V_{k_{i}, l_{i}, c_{i}}$ for $i=1,2$. The symbol $\emptyset$ at an argument replacing a pair ( $k_{i}, l_{i}$ ) means that there is no such intersection.
4.5. Proofs of propositions about incidence vareties. Here we prove the propositions announced in 4.2, The arguments for Propositions 4.10, 4.11 and 4.13 are very similar, so we will spell out the proof for one of these. The discussion will also prepare the ground for the proof of Proposition 4.12, which is substantially more complicated.

Consider first the situation of 4.13 Namely, $m \not \equiv 0,1 \bmod n-2$ is a positive integer, $V_{0, m, j}$ the cell of a full block, $c$ is the label of the full block immediately above the position $(0, m), S_{c} \subseteq B_{m+1, c}$ is a nonempty maximal subset, and $S \subseteq B_{m, j}$ is a maximal subset which is allowed by $S_{c}$.

Proof of Proposition 4.13. By Proposition4.22(3) and (6), for an arbitrary $U \in V_{S},\left(U, U_{c}\right) \in X_{S}^{S_{c}}$ if and only if $U \subseteq\left(L_{1}^{-1} U_{c}\right) \cap V_{0, m, j}$. Moreover, the composition of $L_{1}^{-1}$ and the isomorphism $V_{0, m-1, c} \rightarrow V_{0, m, j}$ gives an isomorphism $V_{1, m, c} \rightarrow V_{0, m, j}$. Hence, for a pair $\left(\bar{U}, U_{c}\right) \in Y_{\bar{S}}^{S_{c}}$, we have

$$
\left\{\left.U \in V_{S}\right|_{\bar{U}}:\left(U, U_{c}\right) \in X_{S}^{S_{c}}\right\}=\left(\left(L_{1}^{-1} U_{c}\right) \cap V_{0, m, j}\right) / \bar{U}
$$

and there is a canonical quotient map $V_{0, m, j} / \bar{U} \rightarrow V_{1, m, c} / U_{c}$.
Let us define two families parameterized by $V_{\bar{S}} \times V_{S_{c}, c}$. The family $\mathcal{F}_{c}$ is defined to have the fiber $\left(L_{1}^{-1} U_{c}\right) \cap V_{0, m, j}$ over a pair $\left(\bar{U}, U_{c}\right) \in V_{\bar{S}} \times V_{S_{c}, c}$. This is a family of affine subspaces of $V_{0, m, j}$. The family $\overline{\mathcal{F}}$ is defined to have the fiber $\bar{U} \subset P_{m, j}$ over a pair $\left(\bar{U}, U_{c}\right) \in V_{\bar{S}} \times V_{S_{c}, c}$. This is a family of projective subspaces contained in $\bar{P}_{m, j}$. Since the tautological bundle over any Schubert cell in any Grassmannian is trivial, the two families are trivial with affine, respectively projective space fibres. Consider these families over the subset $Y_{\bar{S}}^{S_{c}} \subset V_{\bar{S}} \times V_{S_{c}, c}$. By construction, over each point of $Y_{\bar{S}}^{S_{c}}$, the fibre of $\overline{\mathcal{F}}$ is a subspace of the projective closure of the fiber of $\mathcal{F}_{c}$ over the same point. In particular, we can take quotients fiberwise. By the considerations above,

$$
X_{S}^{S_{c}}=\mathcal{F}_{c} / \overline{\mathcal{F}}
$$

Moreover, the morphism $\omega \times \mathrm{Id}: X_{S}^{S_{c}} \rightarrow Y_{\bar{S}}^{S_{c}}$ over a pair $\left(\bar{U}, U_{c}\right)$ is given by the quotient morphism $V_{0, m, j} / \bar{U} \rightarrow V_{1, m, c} / U_{c}$ times the identity. This shows (1) and (2).

The injectivity statement (3) follows again from the isomorphism $V_{1, m, c} \cong V_{0, m, j}$ given by $L_{1}$, since for every pair $\left(U, U_{c}\right) \in X_{S}^{S_{1}, S_{2}}$ one has $U_{c}=\left(U, \bar{U}_{c}\right) \cap V_{1, m, c}$.

Consider now the situation of Proposition 4.12, thus $m \equiv 1 \bmod n-2$ is a positive integer, $c_{1}$ and $c_{2}$ are the labels of the divided block immediately above the block at position $(m, j), S_{1} \subseteq B_{m+1, c_{1}}$, $S_{2} \subseteq B_{m+1, c_{2}}$ are nonempty subsets at least one of which is maximal, and $S \subseteq B_{m, j}$ is a maximal subset which is allowed by $S_{1}$ and $S_{2}$.
Lemma 4.26. For $i=1,2$ fix $U_{i} \in V_{S_{i}, c_{i}}$.
(a) For an arbitrary $U \in V_{S},\left(U, U_{1}, U_{2}\right) \in X_{S}^{S_{1}, S_{2}}$ if and only if $U \subseteq$ $J\left(\phi_{1}\left(L_{1}^{-1} U_{1}\right), \phi_{2}\left(L_{1}^{-1} U_{2}\right)\right) \cap V_{0, m, j}$.
(b) If $\left(\bar{U}, U_{1}, U_{2}\right) \in Y_{\bar{S}}^{S_{1}, S_{2}}$, then $\bar{U} \subseteq J\left(\phi_{1}\left(L_{1}^{-1} U_{1}\right), \phi_{2}\left(L_{1}^{-1} U_{2}\right)\right) \cap V_{0, m, j}$.
(c) If $\left(\bar{U}, U_{1}, U_{2}\right) \in Y_{\bar{S}}^{S_{1}, S_{2}}$, then

$$
\left\{\left.U \in V_{S}\right|_{\bar{U}}:\left(U, U_{1}, U_{2}\right) \in X_{S}^{S_{1}, S_{2}}\right\}=\left(J\left(\phi_{1}\left(L_{1}^{-1} U_{1}\right), \phi_{2}\left(L_{1}^{-1} U_{2}\right)\right) \cap V_{0, m, j}\right) / \bar{U}
$$

Proof. (a) By Proposition 4.22 for any pair of vectors $\left(v_{1}, v_{2}\right) \in U_{1} \times U_{2}$ those points of $P_{m, j}$ for which $\left(v_{1}, v_{2}\right) \cap P_{m+1, c_{i}}$ is either $v_{i}$ or empty are exactly those which are on $J\left(\phi_{1}\left(L_{1}^{-1} v_{1}\right), \phi_{2}\left(L_{1}^{-1} v_{2}\right)\right)$. Hence, to satisfy the conditions $U$ has to be a subset of

$$
\bigcup_{\left.v_{2}\right) \in U_{1} \times U_{2}} J\left(\phi_{1}\left(L_{1}^{-1} v_{1}\right), \phi_{2}\left(L_{1}^{-1} v_{2}\right)\right) \cap V_{0, m, j}=J\left(\phi_{1}\left(L_{1}^{-1} U_{1}\right), \phi_{2}\left(L_{1}^{-1} U_{2}\right)\right) \cap V_{0, m, j}
$$

(b) If $\left(\bar{U}, U_{1}, U_{2}\right) \in Y$, then $(\bar{U}) \cap P_{m+1, c_{i}} \subseteq U_{i}$. Hence, $\phi_{i}\left(L_{1}^{-1}\left((\bar{U}) \cap P_{m+1, c_{i}}\right)\right) \subseteq \phi_{i}\left(L_{1}^{-1} U_{i}\right)$, and

$$
J\left(\phi_{1}\left(L_{1}^{-1}\left((\bar{U}) \cap P_{m+1, c_{1}}\right)\right), \phi_{2}\left(L_{1}^{-1}\left((\bar{U}) \cap P_{m+1, c_{2}}\right)\right)\right) \subseteq J\left(\phi_{1}\left(L_{1}^{-1} U_{1}\right), \phi_{2}\left(L_{1}^{-1} U_{2}\right)\right)
$$

By Proposition 4.22 there is an isomorphism $V_{1, m-1, j} \cong V_{0, m-1, c_{1}} \times V_{0, m-1, c_{1}}$ in such a way that

$$
\bar{U} \cap V_{1, m-1, j} \subseteq J\left(\phi_{1}\left(L_{1}^{-1}\left((\bar{U}) \cap P_{m+1, c_{1}}\right)\right), \phi_{2}\left(L_{1}^{-1}\left((\bar{U}) \cap P_{m+1, c_{2}}\right)\right)\right) \cap V_{1, m-1, j}
$$

Similarly, on each cell $V_{k, l, j}$ such that $k+l=m$ and $k \geq 1$, the affine subspace $\bar{U} \cap V_{k, l, j}$ is a subvariety of $J\left(\phi_{1}\left(L_{1}^{-1}\left((\bar{U}) \cap P_{m+1, c_{1}}\right)\right), \phi_{2}\left(L_{1}^{-1}\left((\bar{U}) \cap V_{m+1, c_{2}}\right)\right)\right) \cap V_{k, l, j}$. All these mean that $\bar{U} \subseteq J\left(\phi_{1}\left(L_{1}^{-1} U_{1}\right), \phi_{2}\left(L_{1}^{-1} U_{2}\right)\right)$.
(c) Recall, that $\bar{U}$ also represents a subspace at infinity for $V_{0, m, j}$, and $\left.V_{S}\right|_{\bar{U}}=V_{0, m, j} / \bar{U}$. In fact, we can take the quotient of an arbitrary subspace of $V_{0, m, j}$, whose closure in $P_{m, j}$ contains $\bar{U}$ with respect to (an arbitrary affine subspace representing) $\bar{U}$. Then the statement follows from (a) and (b).

Proof of Proposition 4.12. It follows from the definitions that $(\omega \times \mathrm{Id} \times \mathrm{Id})\left(X_{S}^{S_{1}, S_{2}}\right) \subseteq Y_{\bar{S}}^{S_{1}, S_{2}}$. The surjectivity will follow from the calculation of the fibers.

We will define three families of subspaces in $P_{m, j}$ over $V_{\bar{S}} \times V_{S_{1}, c_{1}} \times V_{S_{2}, c_{2}}$. For $i=1,2$ the family $\mathcal{F}_{i}$ is defined to have the fiber $\phi_{i}\left(L_{1}^{-1} U_{i}\right) \subseteq P_{m, j}$ over a three-tuple $\left(\bar{U}, U_{1}, U_{2}\right) \in V_{\bar{S}} \times V_{S_{1}, c_{1}} \times V_{S_{2}, c_{2}}$. Let the third family $\mathcal{F}$ has the fiber $\bar{U} \subseteq P_{m, j}$ over the same element. This is of course empty, if $|S|=1$. It is important to note, that in all cases the fibers are always projective subspaces of $P_{m, j}$.

By Lemma 4.23, there is an embedding $\phi_{i} \circ L_{1}^{-1}: P_{m+1, c_{3-i}} \rightarrow N_{c_{i}}^{0} \subset P_{m, j}$. Apply this embedding on the fibers of the projectivization of the tautological bundle over the Schubert cell $V_{S_{3-i}}$. Then multiply the base with $V_{\bar{S}} \times V_{S_{i}}$, and extend the family into this direction as a constant. This gives the bundle $\mathcal{F}_{i}$. Again, by the fact that the tautological bundle over any Schubert cell is trivial it follows that the $\mathcal{F}_{i}$ 's are also trivial, that is, $\mathcal{F}_{i} \cong \mathbb{P}^{\left|S_{i}\right|-1} \times V_{\bar{S}} \times V_{S_{1}} \times V_{S_{2}}$. Similarly, $\mathcal{F} \cong \mathbb{P}^{|S|-2} \times V_{\bar{S}} \times V_{S_{1}} \times V_{S_{2}}$.

By Lemma B.4, the join of trivial families over a common base is a trivial family of the joins of the fibers:

$$
\begin{aligned}
J\left(\mathcal{F}_{1}, \mathcal{F}_{2}\right) & =J\left(\mathbb{P}^{\left|S_{1}\right|-1} \times V_{\bar{S}} \times V_{S_{1}} \times V_{S_{2}}, \mathbb{P}^{\left|S_{2}\right|-1} \times V_{\bar{S}} \times V_{S_{1}} \times V_{S_{2}}\right) \\
& \cong J\left(\mathbb{P}^{\left|S_{1}\right|-1}, \mathbb{P}^{\left|S_{2}\right|-1}\right) \times V_{\bar{S}} \times V_{S_{1}} \times V_{S_{2}} \\
& \cong \mathbb{P}^{\left|S_{1}\right|+\left|S_{2}\right|-1} \times V_{\bar{S}} \times V_{S_{1}} \times V_{S_{2}} \subseteq P_{m, j} \times V_{\bar{S}} \times V_{S_{1}} \times V_{S_{2}}
\end{aligned}
$$

Therefore, $J\left(\mathcal{F}_{1}, \mathcal{F}_{2}\right) \cap\left(V_{0, m, j} \times V_{\bar{S}} \times V_{S_{1}} \times V_{S_{2}}\right)$ is a trivial family of affine subspaces of $V_{0, m, j}$ over $V_{\bar{S}} \times V_{S_{1}} \times V_{S_{2}}$.

By Lemma 4.26 (b) $\mathcal{F}$ is a (trivial) subfamily of $J\left(\mathcal{F}_{1}, \mathcal{F}_{2}\right)$ over $Y_{\bar{S}}^{S_{1}, S_{2}}$. By Lemma 4.26 (c) $X_{S}^{S_{1}, S_{2}}$ can be constructed as

$$
X_{S}^{S_{1}, S_{2}}=\left(J\left(\mathcal{F}_{1}, \mathcal{F}_{2}\right) \cap\left(V_{0, m, j} \times V_{\bar{S}} \times V_{S_{1}} \times V_{S_{2}}\right)\right) /\left.\mathcal{F}\right|_{Y_{\bar{S}}^{S_{1}, S_{2}}}
$$

Hence, $X_{S}^{S_{1}, S_{2}}$ is a trivial family of affine spaces of dimension $\left|S_{1}\right|+\left|S_{2}\right|-|S|$, since it is the quotient of a trivial affine family of fibre dimension $\left|S_{1}\right|+\left|S_{2}\right|-1$ by another trivial affine family of fibre dimension $|S|-1$.

## 5. Type $D_{n}$ : special loci

5.1. Support blocks. In this section, we analyze the cases when a cell corresponding to a salient block (see Def. 4.14) of a Young wall $Y$ fails to contain a generator of a corresponding ideal $I \in Z_{Y}$. As an example, recall once again Example 4.5, where the divided missing blocks at position $(1,3)$ are salient blocks of $Y_{3}$, but the corresponding cells do not necessarily contain generators of an ideal $I \in Z_{Y_{3}}$. That this phenomenon can happen at all is one of the main sources of difficulty in our analysis of the strata of the singular Hilbert scheme. We introduce the notion of a support block for a salient block. Intuitively, the intersection of an ideal $I$ with the cell in the support block can generate the intersection in the salient block (such as the support block at position $(0,3)$ for $Y_{3}$ ), and thus the salient block contains no new generator of $I$. We will make this statement more precise in the rest of this section.

We start with some combinatorial preliminaries. Recall the setup of Proposition 4.12, $m \equiv$ $1 \bmod n-2 ; c_{1}$ and $c_{2}$ are the labels of the divided block immediately above the block of label $j$ at position $(m, 0) ; S_{1} \subseteq B_{m+1, c_{1}}$ and $S_{2} \subseteq B_{m+1, c_{2}}$ are nonempty subsets at least one of which is maximal; $S \subseteq B_{m, j}$ is a maximal subset which is allowed by $S_{1}$ and $S_{2}$.

For a half-block $b$ of $S_{i}$, consider the following two conditions.
(1) The blocks below or to the left of $b$ are not contained in $S$.
(2) The block below $b$ is contained in $S$, the complementary half-block $b^{\prime}$ is contained in $S_{3-i}$, and the block to the left of their position in not contained in $S$.
For $i, j=1,2$, let us denote by $S_{i}^{s, j} \subset S_{i}$ the subset of half-blocks of label $c_{i}$ satisfying condition $(j)$. Let moreover $S_{i}^{s}=S_{i}^{s, 1} \cup S_{i}^{s, 2}$.

The next lemma, whose proof is immediate, connects the global Definition 4.14 with the local conditions (1)-(2) where we consider only the $m$-th and $m+1$-st diagonals for a particular $m$, and index sets $S, S_{1}$ and $S_{2}$ as above.

Lemma 5.1. Given a Young wall $Y \in \mathcal{Z}_{\Delta}$, let $S$, respectively $S_{1}$ and $S_{2}$ denote the set of missing blocks, respectively half-blocks of $Y$ on the $m$-th and $(m+1)$-st diagonals. The blocks $S_{i}^{s} \subset S_{i}$ are exactly the salient blocks of $Y$ of label $c_{i}$ on the $(m+1)$-st diagonal.

Thus we can legitimately call the blocks in $S_{i}^{s}$ salient blocks in this local situation.
Let us introduce the following subsets of $S$.

- $S^{l}$ consists of blocks $b \in S$ that are directly to the left of a divided block with labels $\left(c_{1}, c_{2}\right)$.
- $S^{b, 0}$ consists of blocks $b \in S$, so that $b$ is immediately below a divided block with labels ( $c_{1}, c_{2}$ ), the block immediately up-left of $b$ is not in $S$, and both of the divided blocks above $b$ are in $S_{1} \cup S_{2}$.
- $S^{b, c_{i}}$ consists of blocks $b \in S$ that are immediately below a divided block, so that the block immediately up-left of $b$ is not in $S$, and the block of label $c_{3-i}$ above $b$ is in $S_{3-i}$.
- $S^{b, c_{1} \cup c_{2}}$ consists of blocks $b \in S$ such that $b$ lies immediately below a divided block, and the block immediately up-left of $b$ is contained in $S$.
- $S^{b}=S^{b, 0} \cup S^{b, c_{1}} \cup S^{b, c_{2}} \cup S^{b, c_{1} \cup c_{2}}$.

Note that by the Young wall rules, we necessarily have $S^{b}=S \backslash S^{l}$. We will call the blocks in the set $S^{c_{i}}=S^{b, 0} \cup S^{b, c_{i}} \cup S^{b, c_{1} \cup c_{2}}$ support blocks for label $c_{i}$. We will define a support relation from $S^{c_{i}}$ to $S_{i}^{S}$ in 5.2 below.
5.2. Special loci in orbifold strata and the supporting rules. Let $Y \in \mathcal{Z}_{\Delta}$ be a Young wall with a salient block in its bottom row in position $(0, m)$ with $m \equiv 1 \bmod n-2$, a full block immediately below a divided block with labels $\left(c_{1}, c_{2}\right)$. As before, let $S$, respectively $S_{1}$ and $S_{2}$ denote the set of missing blocks, respectively half-blocks of $Y$ on the $m$-th and ( $m+1$ )-st diagonals with the corresponding labels.

We introduce index sets depending on $S, S_{1}$ and $S_{2}$. We consider two cases.
If $\bar{S}$ is not maximal, then let

$$
I\left(S, S_{1}, S_{2}\right)=\left\{\left(k_{1}, l_{1}, k_{2}, l_{2}\right): \quad \begin{array}{cc}
\left(k_{i}, l_{i}\right) \in S_{i}^{s} \cup\{\emptyset\} \cup\left(\{(1, m)\} \cap S_{i}\right) \text { for } i=1,2, \\
\text { and at least one }\left(k_{i}, l_{i}\right)=(1, m)
\end{array}\right\} .
$$

We partition this index set into the following (possibly empty) disjoint subsets:

- $I\left(S, S_{1}, S_{2}\right)_{0}=\left\{\left(k_{1}, l_{1}, k_{2}, l_{2}\right) \in I\left(S, S_{1}, S_{2}\right):\left(k_{i}, l_{i}\right) \notin\{\emptyset,(1, m)\}\right.$ for some $\left.i=1,2\right\}$;
- $I\left(S, S_{1}, S_{2}\right)_{1}=\{(1, m, \emptyset),(\emptyset, 1, m)\} \cap I\left(S, S_{1}, S_{2}\right)$;
- $I\left(S, S_{1}, S_{2}\right)_{-1}=\{(1, m, 1, m)\} \cap I\left(S, S_{1}, S_{2}\right)$.

If $\bar{S}$ is maximal, then let

$$
I\left(S, S_{1}, S_{2}\right)=\left\{\left(k_{1}, l_{1}, k_{2}, l_{2}\right):\left(k_{i}, l_{i}\right) \in S_{i}^{s} \cup\{\emptyset\} \text { for } i=1,2\right\}
$$

We remark that in this case $(1, m) \notin S_{i}^{s}$ for both $i=1,2$. The index set $I\left(S, S_{1}, S_{2}\right)$ in this case can be partitioned into the following subsets:

- $I\left(S, S_{1}, S_{2}\right)_{0}=\left\{\left(k_{1}, l_{1}, k_{2}, l_{2}\right) \in I\left(S, S_{1}, S_{2}\right):\left(k_{i}, l_{i}\right) \neq \emptyset\right.$ for some $\left.i=1,2\right\} ;$
- $I\left(S, S_{1}, S_{2}\right)_{1}=\{(\emptyset, \emptyset)\}$.

As before, $\emptyset$ is used as a symbol replacing a pair in these defintions.
For projective subspaces $P_{1} \subseteq P_{2} \subseteq P_{m+1, c}$ we introduce the following notation. $\left(P_{2} \backslash P_{1}\right) \dashv V_{k, l, c}$ if and only if $\left(P_{2} \backslash P_{1}\right) \cap V_{k, l, c} \neq \emptyset$ and $k$ is maximal with this property. This is the smallest cell whose intersection with $P_{2}$ is larger than that with $P_{1}$.

Recall the truncated Young wall $\bar{Y}$ and the morphism $T: Z_{Y} \rightarrow Z_{\bar{Y}}$ from 4.3. The following statement will be proved below in 5.4

Theorem 5.2. There is a decomposition into locally closed subspaces

$$
Z_{Y}=\bigsqcup_{\left(k_{1}, l_{1}, k_{2}, l_{2}\right) \in I\left(S, S_{1}, S_{2}\right)} Z_{Y}\left(k_{1}, l_{1}, k_{2}, l_{2}\right),
$$

where

$$
Z_{Y}\left(k_{1}, l_{1}, k_{2}, l_{2}\right)=\left\{I \in Z_{Y}:\left(\left(I \cap P_{m, j}\right) \backslash\left(I \cap \bar{P}_{m, j}\right)\right) \cap P_{m+1, c_{i}} \dashv V_{k_{i}, l_{i}, c_{i}} \text { for } i=1,2\right\} .
$$

The symbol $\emptyset=\left(k_{i}, l_{i}\right)$ means that there is no intersection with $P_{m+1, c_{i}}$. Moreover, if $\left(k_{1}, l_{1}, k_{2}, l_{2}\right) \in I\left(S, S_{1}, S_{2}\right)_{e}$, then the nonempty fibers of $T: Z_{Y}\left(k_{1}, l_{1}, k_{2}, l_{2}\right) \rightarrow Z_{\bar{Y}}$ have Euler characterestic e.

The space $Z_{Y}\left(k_{1}, l_{1}, k_{2}, l_{2}\right)$ should be thought as the space of those ideals, where the generator in the cell of support block at position $(0, m)$ has an image on the $(m+1)$-st diagonal at the cells $V_{k_{i}, l_{i}, c_{i}}$ for $i=1,2$ which does not come from the rows above. We remark that the mentioned support block in the bottom row is also salient since it is the first missing block in its row. But if one attaches further rows to the bottom of $Y$, then it may become non-salient.

Recall from 5.1 the set $S^{c_{i}}=S^{b, 0} \cup S^{b, c_{i}} \cup S^{b, c_{1} \cup c_{2}}$ of support blocks for label $c_{i}$. In the light of the definition of the sets $I\left(S, S_{1}, S_{2}\right)$ and Theorem 5.2 we will say that the blocks in $S^{c_{i}}$ can support the salient blocks of label $c_{i}$ in $S_{i}^{s}$ in the following sense.

Each support block $b \in S$ for label $c_{i}$ can support at most one salient block of label $c_{i}$ above and to the left of $b$ on the $(m+1)$-st antidiagonal. More precisely, this supporting relationship has to respect the following supporting rules.

- Each block in $S^{b, 0}$ supports precisely one or two salient blocks, at most one from each label, and at least one of these has to be immediately above it;
- each block in $S^{b, c_{i}}$ can support at most one salient block of label $c_{i}$ which is not immediately above it;
- each block in $S^{b, c_{1} \cup c_{2}}$ can support none, one or two salient blocks, at most one from each label, and neither of these is immediately above it.
In this way we define a correspondence from a subset of $S^{c_{1}}$ to $S_{1}^{s}$ and one from a subset of $S^{c_{2}}$ to $S_{2}^{s}$ but these two have to satisfy the restrictions mentioned above on the intersection $S^{c_{1}} \cap S^{c_{2}}=$ $S^{b, 0} \cup S^{b, c_{1} \cup c_{2}}$. Neither correspondence has to be surjective or be defined on the whole domain, but, where they are defined, they should be injective.

We shall call a salient block $b^{\prime} \in S_{i}$ of label $c_{i}$ supported, if the number of support blocks for label $c_{i}$ in $S$ which are below $b$ is at least as much as the total number of salient blocks of label $c_{i}$ in $S_{i}$ counted from the top left, including $b^{\prime}$ itself. A supported salient block $b^{\prime}$ satisfying condition (2) above, so that there is a support block in $S$ immediately below $b^{\prime}$, will be called directly supported. The others will be called non-directly supported. The supporting relationship will be globalized for the whole diagram in the notion of closing datum, to be defined in 6.1 below.

Recall that during the inductive process in the proof of Theorem4.1, at each step a new generator appears in the cell corresponding to the salient block in the bottom row. Assume that for $I \in Z_{Y}$, $\left(I \cap P_{m, j}\right) \cap P_{m+1, c_{i}}=I \cap P_{m+1, c_{i}}$ for $i=1,2$. In this case we will say that there is no generator of label $c_{i}$ on the $(m+1)$-st antidiagonal. Let $S, S_{1}$ and $S_{2}$ be the index sets for $V_{0, m, j}, V_{1, m, c_{1}}$ and $V_{1, m, c_{2}}$ respectively. Then using inductively Theorem 5.2 for each block $b \in S^{c_{i}}$ we get that there is at most one block $b_{i} \in S_{i}^{s}$ such that when the row of $b$ is added to the Young wall, the new generator in the cell of $b$ has nontrivial image in the cell of $b_{i}$. Conversely, for each block $b_{i} \in S_{i}^{s}$ there corresponds a support block $b \in S^{c_{i}}$ determined by $I$. In particular, this implies
Corollary 5.3. Assume that for $I \in Z_{Y},\left(I \cap P_{m, j}\right) \cap P_{m+1, c_{i}}=I \cap P_{m+1, c_{i}}$ for some $i=1,2$. Let $S, S_{1}$ and $S_{2}$ be the index sets for $V_{0, m, j}, V_{1, m, c_{1}}$ and $V_{1, m, c_{2}}$ respectively.
(1) $\left|S_{i}^{u, 1}\right| \leq\left|S^{b, c_{i}}\right|+\left|S^{b, c_{1} \cup c_{2}}\right|$;
(2) every salient block of label $c_{i}$ is supported;
(3) to each salient block of label $c_{i}$ there corresponds a unique support block for label $c_{i}$ in the way described above.
5.3. Special loci in Grassmannians. We prepare the ground for the proof of Theorem 5.2 by analyzing the incidence varieties of 4.2 in the case $m \equiv 1 \bmod n-2$. Once again, we use the notations of Proposition 4.12, The composition with the projection from $V_{S} \times V_{S_{1}, c_{1}} \times V_{S_{2}, c_{2}}$ to its first factor, followed by the affine linear fibration $\omega: V_{S} \rightarrow V_{\bar{S}}$, defines a projection map $p_{V_{\bar{S}}}: V_{S} \times V_{S_{1}, c_{1}} \times V_{S_{2}, c_{2}} \rightarrow V_{\bar{S}}$.

For $i=1,2$ let $S_{i}(\bar{U})=\left\{\left(k_{i}, l_{i}\right) \in B_{m+1, c_{i}}:(\bar{U}) \cap V_{k_{i}, l_{i}, c_{i}}=\emptyset\right\}$ be the blocks in the partial profile of $(\bar{U}) \cap P_{m+1, c_{i}}$ on the $(m+1)$-st diagonal. Then the index sets $I\left(S, S_{1}(\bar{U}), S_{2}(\bar{U})\right)$ and $I\left(S, S_{1}(\bar{U}), S_{2}(\bar{U})\right)_{e}$ introduced above make sense. The following lemma stratifies the fibers of the affine linear fibration $\omega: V_{S} \rightarrow V_{\bar{S}}$.

Lemma 5.4. For any $\bar{U} \in V_{\bar{S}}$, there is a stratification

$$
V_{0, m, j} / \bar{U}=\bigsqcup_{\left(k_{1}, l_{1}, k_{2}, l_{2}\right) \in I\left(S, S_{1}(\bar{U}), S_{2}(\bar{U})\right)} V_{0, m, j}\left(k_{1}, l_{1}, k_{2}, l_{2}\right) / \bar{U},
$$

where

$$
V_{0, m, j}\left(k_{1}, l_{1}, k_{2}, l_{2}\right) / \bar{U}=\left\{U \in V_{0, m, j} / \bar{U}:((U) \backslash(\bar{U})) \cap P_{m+1, c_{i}} \dashv V_{k_{i}, l_{i}, c_{i}} \text { for } i=1,2\right\} .
$$

Moreover, if $\left(k_{1}, l_{1}, k_{2}, l_{2}\right) \in I\left(S, S_{1}(\bar{U}), S_{2}(\bar{U})\right)_{e}$, then the space $V_{0, m, j}\left(k_{1}, l_{1}, k_{2}, l_{2}\right) / \bar{U}$ is of Euler characterestic e.

Proof. We have to distinguish the cases when $\bar{S}$ is maximal or not. The latter case is significantly simpler, so we start with that.

If $\bar{U} \in V_{\bar{S}}$ with $\bar{S}$ not maximal, then $\left(M_{c_{1}}+\bar{U}\right) \cap\left(M_{c_{2}}+\bar{U}\right)=\emptyset$ since $M_{c_{1}}$ and $M_{c_{2}}$ are distinct parallel hyperplanes in $V_{0, m, j}$, and there are affine subspaces $U_{i}$ representing $\bar{U}$ such that $U_{i} \subseteq M_{c_{i}}$. Recall from Corollary 4.25 the stratification of $V_{0, m, j}$ which basically comes from the join structure on its closure $P_{m, j}$. This induces a decomposition of $V_{0, m, j} / \bar{U}$ into non-empty, locally closed, but not necessarily disjoint spaces

$$
\begin{aligned}
V_{0, m, j} / \bar{U}= & \left(\left(V_{0, m, j} \backslash\left(M_{c_{1}} \cap M_{c_{2}}\right)\right) / \bar{U}\right) \cup\left(M_{c_{1}} / \bar{U}\right) \cup\left(M_{c_{1}} / \bar{U}\right) \\
& \bigcup\left(\bigcup_{\left(k_{1}, l_{1}\right) \in B_{m+1, c_{1}}}\left(\left(J\left(V_{0, m-1, c_{1}}, V_{k_{1}, l_{1}, c_{2}}\right) \cap V_{0, m, j}\right) \backslash M_{c_{1}}^{0}+\bar{U}\right) / \bar{U}\right) \\
& \bigcup\left(\bigcup_{\left(k_{2}, l_{2}\right) \in B_{m+1, c_{1}}}\left(\left(J\left(V_{0, m-1, c_{2}}, V_{k_{2}, l_{2}, c_{1}}\right) \cap V_{0, m, j}\right) \backslash M_{c_{2}}^{0}+\bar{U}\right) / \bar{U}\right) .
\end{aligned}
$$

Consider a block $\left(k_{i}, l_{i}\right) \in B_{m+1, c_{i}} \backslash S_{i}(\bar{U})$. Then the intersection $(\bar{U}) \cap V_{k_{i}, l_{i}, c_{i}} \neq \emptyset$. Assume that there is an $U \in V_{0, m, j} / \bar{U}$ such that $((U) \backslash(\bar{U})) \cap V_{k_{i}, l_{i}, c_{i}} \neq \emptyset$. Then $\operatorname{dim}\left((U) \cap V_{k_{i}, l_{i}, c_{i}}\right)>$ $\operatorname{dim}\left((\bar{U}) \cap V_{k_{i}, l_{i}, c_{i}}\right)$ so by Lemma3.6 there is at least one other block in a row above $k_{i}$ which has a trivial intersection with $(\bar{U})$ but a nontrivial one with $(U)$. Hence, for any $\left(k_{i}, l_{i}\right) \in B_{m+1, c_{i}} \backslash S_{i}(\bar{U})$ we have

$$
\left\{U \in V_{0, m, j} / \bar{U}:((U) \backslash(\bar{U})) \cap P_{m+1, c_{i}} \dashv V_{k_{i}, l_{i}, c_{i}}\right\}=\emptyset .
$$

On the other hand, if $\left(k_{i}, l_{i}\right) \in S_{i}(\bar{U}) \cup\{\emptyset\}$ then

$$
\begin{gathered}
\left(\left(J\left(V_{0, m-1, c_{i}}, V_{k_{i}, l_{i}, c_{3-i}}\right) \cap V_{0, m, j}\right) \backslash M_{c_{i}}^{0}+\bar{U}\right) / \bar{U}= \\
\left\{U \in V_{0, m, j} / \bar{U}:((U) \backslash(\bar{U})) \cap P_{m+1, c_{i}} \dashv V_{k_{i}, l_{i}, c_{i}} \text { and }((U) \backslash(\bar{U})) \cap P_{m+1, c_{3-i}} \dashv V_{1, m, c_{3-i}}\right\} .
\end{gathered}
$$

By dimension constrains these spaces are disjoint and it is easy to see that together with ( $V_{0, m, j} \backslash$ $\left.\left(M_{c_{1}} \cup M_{c_{2}}\right)\right) / \bar{U}, M_{c_{1}}^{0} / \bar{U}$, and $M_{c_{2}}^{0} / \bar{U}$ they cover $V_{0, m, j} / \bar{U}$. Thus we get a stratification

$$
\begin{aligned}
& V_{0, m, j} / \bar{U}=V_{0, m, j}(1, m, 1, m) / \bar{U} \\
& \bigsqcup\left(\bigsqcup_{\left.\left(k_{1}, l_{1}\right) \in\left(S_{1} \overline{( } \bar{U}\right) \cup\{\emptyset\}\right) \backslash\{(1, m)\}} V_{0, m, j}\left(k_{1}, l_{1}, 1, m\right) / \bar{U}\right) \\
& \bigsqcup\left(\underset{\left(k_{2}, l_{2}\right) \in\left(S_{2}(\bar{U}) \cup\{\emptyset\}\right) \backslash\{(1, m)\}}{\bigsqcup_{0}} V_{0, m, j}\left(1, m, k_{2}, l_{2}\right) / \bar{U}\right) .
\end{aligned}
$$

In particular, there is a stratification

$$
\begin{equation*}
M_{c_{3-i}} / \bar{U}=\bigsqcup_{\left(k_{i}, l_{i}\right) \in\left(S_{i}(\bar{U}) \cup\{\emptyset\}\right) \backslash\{(1, m)\}} V_{0, m, j}\left(k_{i}, l_{i}, 1, m\right) / \bar{U} \tag{13}
\end{equation*}
$$

Being an affine space, $M_{c_{i}}^{0} / \bar{U}$ has Euler characteristic 1 for $i=1,2$. By Lemma B. 3 the spaces $\left(J\left(V_{0, k(n-2), c_{i}}, V_{k_{i}, l_{i}, c_{3-i}}\right) \cap V_{0, m, j}\right) \backslash M_{c_{i}}^{0}$ have Euler characteristic 0, and the same is true for $\left(\left(J\left(V_{0, k(n-2), c_{i}}, V_{k_{i}, l_{i}, c_{3-i}}\right) \cap V_{0, m, j}\right) \backslash M_{c_{i}}^{0}+\bar{U}\right) / \bar{U}$. This last step follows from the fact that the subspace $\bar{U} \subset \bar{P}_{m, j}$ avoids both the image of $V_{k_{i}, l_{i}, c_{3-i}}$ and $M_{c_{i}}^{0}$.

If $\bar{U} \in V_{\bar{S}}$ such that $\bar{S}$ is maximal, then $\left(M_{c_{1}}+\bar{U}\right)=\left(M_{c_{2}}+\bar{U}\right)=V_{0, m, j}$ since $\bar{U}$ is transversal to $M_{c_{1}}$ and $M_{c_{2}}$. Therefore, there are two stratifications for $V_{0, m, j} / \bar{U}$ with $i=1,2$ as in (13). The claimed stratification is the largest common refinement of these two. In particular, there are three types of strata. First, if $U \in\left(\left(J\left(V_{0, k(n-2), c_{1}}, V_{k_{1}, l_{1}, c_{2}}\right) \cap V_{0, m, j}\right) \backslash M_{c_{1}}^{0}+\bar{U}\right) \cap$ $\left(\left(J\left(V_{0, k(n-2), c_{2}}, V_{k_{2}, l_{2}, c_{1}}\right) \cap V_{0, m, j}\right) \backslash M_{c_{2}}^{0}+\bar{U}\right) / \bar{U}$ for arbitrary $\left(k_{1}, l_{1}\right) \in S_{1}(\bar{U})$ and $\left(k_{2}, l_{2}\right) \in S_{2}(\bar{U})$,
then $((U) \backslash(\bar{U})) \cap P_{m+1, c_{i}} \dashv V_{k_{i}, l_{i}, c_{i}}$ for $i=1,2$. Second, if $U \in\left(\left(\left(J\left(V_{0, k(n-2), c_{1}}, V_{k_{i}, l_{i}, c_{3-i}}\right) \cap\right.\right.\right.$ $\left.\left.\left.V_{0, m, j}\right)+\bar{U}\right) \backslash \bigcup_{\left(k_{3-i}, l_{3-i}\right) \in S_{3-i}(\bar{U})}\left(J\left(V_{0, k(n-2), c_{2}}, V_{k_{3-i}, l_{3-i}, c_{i}}\right) \cap V_{0, m, j}\right)+\bar{U}\right) / \bar{U}$, then $((U) \backslash(\bar{U})) \cap$ $P_{m+1, c_{i}} \dashv V_{k_{i}, l_{i}, c_{i}}$ but $((U) \backslash(\bar{U})) \cap P_{m+1, c_{3-i}}=\emptyset$. Third, if $U \in\left(\left(M_{c_{1}}^{0}+\bar{U}\right) \cap\left(M_{c_{2}}^{0}+\bar{U}\right)\right) / \bar{U}$, then $((U) \backslash(\bar{U})) \cap P_{m+1, c_{i}}=\emptyset$ for $i=1,2$. To sum it up, we get a stratification into locally closed spaces

$$
V_{0, m, j} / \bar{U}=\bigsqcup_{\substack{\left(k_{1}, l_{1}\right) \in S_{1}(\bar{U}) \cup\{\emptyset\} \\\left(k_{2}, l_{2}\right) \in S_{2}(\bar{U}) \cup\{\emptyset\}}} V_{0, m, j}\left(k_{1}, l_{1}, k_{2}, l_{2}\right) / \bar{U}
$$

The Euler characteristic of the stratum $V_{0, m, j}(\emptyset, \emptyset) / \bar{U}=\left(\left(M_{c_{1}}+\bar{U}\right) \cap\left(M_{c_{2}}+\bar{U}\right)\right) / \bar{U}$ is 1 . It is left to the reader that the others have Euler characteristic 0 .

Let $\mathcal{I}_{S}=\left\{\left(S_{1}(\bar{U}), S_{2}(\bar{U})\right): \bar{U} \in V_{\bar{S}}\right\}$. Actually $\mathcal{I}_{S}$ only depends on $\bar{S}$. For each $\left(S_{1}^{\prime}, S_{2}^{\prime}\right) \in \mathcal{I}_{S}$, let

$$
V_{\bar{S}}\left(S_{1}^{\prime}, S_{2}^{\prime}\right)=\left\{\bar{U} \in V_{\bar{S}}:\left(S_{1}(\bar{U}), S_{2}(\bar{U})\right)=\left(S_{1}^{\prime}, S_{2}^{\prime}\right)\right\}
$$

Corollary 5.5. For a fixed $\left(S_{1}^{\prime}, S_{2}^{\prime}\right) \in \mathcal{I}(S)$ and $\left(k_{1}, l_{1}, k_{2}, l_{2}\right) \in I\left(S, S_{1}^{\prime}, S_{2}^{\prime}\right)$ the spaces $V_{0, m, j}\left(k_{1}, l_{1}, k_{2}, l_{2}\right) / \bar{U}$ are isomorphic for every $\bar{U} \in V_{\bar{S}}\left(S_{1}^{\prime}, S_{2}^{\prime}\right)$. Moreover, they fit together into a locally closed subvariety $V_{S}\left(k_{1}, l_{1}, k_{2}, l_{2}\right) \subseteq V_{S}$ which is a trivial family over $V_{\bar{S}}\left(S_{1}^{\prime}, S_{2}^{\prime}\right)$. Using induction and the fact that the fiber product of locally closed spaces is locally closed, we get that there is a stratification

$$
V_{S}=\bigsqcup_{\left(S_{1}^{\prime}, S_{2}^{\prime}\right) \in \mathcal{I}_{S}}\left(\bigsqcup_{\left(k_{1}, l_{1}, k_{2}, l_{2}\right) \in I\left(S, S_{1}^{\prime}, S_{2}^{\prime}\right)} V_{S}\left(k_{1}, l_{1}, k_{2}, l_{2}\right)\right)
$$

into locally closed subvarieties. Furthermore, if $\left(k_{1}, l_{1}, k_{2}, l_{2}\right) \in I\left(S, S_{1}^{\prime}, S_{2}^{\prime}\right)_{e}$, then the fiber of $\omega: V_{S}\left(k_{1}, l_{1}, k_{2}, l_{2}\right) \rightarrow V_{\bar{S}}$ has Euler characteristic e.
Proof. The triviality of the family $V_{S}\left(k_{1}, l_{1}, k_{2}, l_{2}\right) \rightarrow V_{\bar{S}}\left(S_{1}^{\prime}, S_{2}^{\prime}\right)$ follows from Lemma B. 4 and the fact that $V_{0, m, j}\left(k_{1}, l_{1}, k_{2}, l_{2}\right) / \bar{U}$ is constructed using (union, intersection and difference of) joins in $P_{m, j}$. The rest of the statement is obvious.
5.4. Proof of Theorem5.2. As before, we fix $S, S_{1}$ and $S_{2}$. Recall that the fiber of the morphism $\omega \times \mathrm{Id} \times \mathrm{Id}: X_{S}^{S_{1}, S_{2}} \rightarrow Y_{\bar{S}}^{S_{1}, S_{2}}$ over an element $\left(\bar{U}, U_{1}, U_{2}\right) \in Y_{\bar{S}}^{S_{1}, S_{2}}$ is $J\left(L_{1}^{-1} U_{1}, L_{1}^{-1} U_{2}\right) / \bar{U}$.

For $i=1,2$ let $S_{i}(\bar{U})=\left\{\left(k_{i}, l_{i}\right) \in S_{i}:(\bar{U}) \cap V_{k_{i}, l_{i}, c_{i}}=\emptyset\right\}$ be the blocks in the partial profile of $(\bar{U}) \cap P_{m+1, c_{i}}$ on the $(m+1)$-st diagonal. Then the index sets $I\left(S, S_{1}(\bar{U}), S_{2}(\bar{U})\right)$ and $I\left(S, S_{1}(\bar{U}), S_{2}(\bar{U})\right)_{e}$ are defined. The following lemma, whose proof is the same as that of Lemma 5.4] stratifies the fibers of the affine linear fibration $\omega \times \operatorname{Id} \times \operatorname{Id}: X_{S}^{S_{1}, S_{2}} \rightarrow Y_{\bar{S}}^{S_{1}, S_{2}}$.

Lemma 5.6. For any $U_{1} \in V_{S_{1}, c_{1}}, U_{2} \in V_{S_{2}, c_{2}}$ the stratification of Lemma 5.4 restricts to $a$ stratification

$$
J\left(L_{1}^{-1} U_{1}, L_{1}^{-1} U_{2}\right) / \bar{U}=\bigsqcup_{\left(k_{1}, l_{1}, k_{2}, l_{2}\right) \in I\left(S, S_{1}(\bar{U}), S_{2}(\bar{U})\right)} J\left(L_{1}^{-1} U_{1}, L_{1}^{-1} U_{2}\right)\left(k_{1}, l_{1}, k_{2}, l_{2}\right) / \bar{U}
$$

where

$$
\begin{gathered}
J\left(L_{1}^{-1} U_{1}, L_{1}^{-1} U_{2}\right)\left(k_{1}, l_{1}, k_{2}, l_{2}\right) / \bar{U}= \\
\left\{U \in J\left(L_{1}^{-1} U_{1}, L_{1}^{-1} U_{2}\right) / \bar{U}:((U) \backslash(\bar{U})) \cap P_{m+1, c_{i}} \dashv V_{k_{i}, l_{i}, c_{i}} \text { for } i=1,2\right\} .
\end{gathered}
$$

Moreover, if $\left(k_{1}, l_{1}, k_{2}, l_{2}\right) \in I\left(S, S_{1}(\bar{U}), S_{2}(\bar{U})\right)_{e}$, then the space $J\left(L_{1}^{-1} U_{1}, L_{1}^{-1} U_{2}\right)\left(k_{1}, l_{1}, k_{2}, l_{2}\right) / \bar{U}$ is of Euler characterestic e.

Let $\mathcal{I}_{S}^{S_{1}, S_{2}}=\left\{\left(S_{1}(\bar{U}), S_{2}(\bar{U})\right):\left(\bar{U}, U_{1}, U_{2}\right) \in Y_{\bar{S}}^{S_{1}, S_{2}}\right\}$. Actually $\mathcal{I}_{S}^{S_{1}, S_{2}}$ only depends on $\bar{S}, S_{1}$ and $S_{2}$. For each $\left(S_{1}^{\prime}, S_{2}^{\prime}\right) \in \mathcal{I}_{S}^{S_{1}, S_{2}}$, let

$$
Y_{\bar{S}}^{S_{1}, S_{2}}\left(S_{1}^{\prime}, S_{2}^{\prime}\right)=\left\{\left(\bar{U}, U_{1}, U_{2}\right) \in Y_{\bar{S}}^{S_{1}, S_{2}}:\left(S_{1}(\bar{U}), S_{2}(\bar{U})\right)=\left(S_{1}^{\prime}, S_{2}^{\prime}\right)\right\}
$$

Corollary 5.7. For fixed $\left(k_{1}, l_{1}, k_{2}, l_{2}\right) \in I\left(S, S_{1}^{\prime}, S_{2}^{\prime}\right)$ the spaces $J\left(L_{1}^{-1} U_{1}, L_{1}^{-1} U_{2}\right)\left(k_{1}, l_{1}, k_{2}, l_{2}\right) / \bar{U}$ are isomorphic for every $\bar{U} \in Y_{\bar{S}}^{S_{1}, S_{2}}\left(S_{1}^{\prime}, S_{2}^{\prime}\right)$. Moreover, they fit together into a locally closed
subvariety $X_{S}^{S_{1}, S_{2}}\left(k_{1}, l_{1}, k_{2}, l_{2}\right) \subseteq X_{S}^{S_{1}, S_{2}}$. Using induction and the fact that the fiber product of locally closed spaces is locally closed, we get that there is a stratification

$$
X_{S}^{S_{1}, S_{2}}=\bigsqcup_{\left(S_{1}^{\prime}, S_{2}^{\prime}\right) \in \mathcal{I}_{S}^{S_{1}, S_{2}}}\left(\bigsqcup_{\left(k_{1}, l_{1}, k_{2}, l_{2}\right) \in I\left(S, S_{1}^{\prime}, S_{2}^{\prime}\right)} X_{S}^{S_{1}, S_{2}}\left(k_{1}, l_{1}, k_{2}, l_{2}\right)\right)
$$

into a locally closed subvarieties. Furthermore, if $\left(k_{1}, l_{1}, k_{2}, l_{2}\right) \in I\left(S, S_{1}^{\prime}, S_{2}^{\prime}\right)_{e}$, then the fiber of $\omega \times \mathrm{Id} \times \mathrm{Id}: X_{S}^{S_{1}, S_{2}}\left(k_{1}, l_{1}, k_{2}, l_{2}\right) \rightarrow Y_{\bar{S}}^{S_{1}, S_{2}}\left(S_{1}^{\prime}, S_{2}^{\prime}\right)$ has Euler characteristic e.
Proof of Theorem 5.2. With all these preparations the proof itself is very easy. We just observe that $Z_{Y}\left(k_{1}, l_{1}, k_{2}, l_{2}\right)$ consist of those points in $Z_{Y}$, which map in (10) to $X_{S}^{S_{1}, S_{2}}\left(k_{1}, l_{1}, k_{2}, l_{2}\right)$ for some $\left(S_{1}^{\prime}, S_{2}^{\prime}\right) \in \mathcal{I}_{S}^{S_{1}, S_{2}}$ such that $\left(k_{1}, l_{1}, k_{2}, l_{2}\right) \in I\left(S, S_{1}^{\prime}, S_{2}^{\prime}\right)$. The result then follows from Corollary 5.7

## 6. Type $D_{n}$ : Decomposition of the coarse Hilbert scheme

6.1. Distinguished 0-generated Young walls. In this section, we describe some distinguished subsets of the set of Young walls $\mathcal{Z}_{\Delta}$ of type $D_{n}$. They will consist of Young walls which are the analogues of the 0-generated partitions from 2.2, as always, in the type $D$ case there are substantial extra complications. For a Young wall $Y \in \mathcal{Z}_{\Delta}$, denote by $\mathrm{wt}_{0}(Y)$ the 0 -weight of $Y$, the number of half-blocks labelled 0 in $Y$.

Recall from 4.3, respectively 5.1 the notions of a salient block and a support block for a given label $c \in\{0,1, n-1, n\}$; we will use also all other notation introduced in the latter section. We call a pair of salient half-blocks $\left(b, b^{\prime}\right)$ sharing the same position a salient block-pair.

Consider the following conditions for a Young wall $Y \in \mathcal{Z}_{\Delta}$.
(A1) All salient blocks of $Y$ are labelled $0,1, n-1$ or $n$.
(A2) Every salient block of $Y$ labelled $c \in\{1, n-1, n\}$ is supported.
Let $\mathcal{Z}_{\Delta}^{\prime} \subset \mathcal{Z}_{\Delta}$ be the set of Young walls $Y$ which satisfy conditions (A1)-(A2). We will prove in Theorem 6.4 that $\mathcal{Z}_{\Delta}^{\prime}$ is the set of Young walls $Y \in \mathcal{Z}_{\Delta}$ for which $Z_{Y} \cap \operatorname{im}\left(i^{*}\right) \neq \emptyset$, where $i^{*}$ is the pullback map defined in 1.1.

Given a Young wall $Y \in \mathcal{Z}_{\Delta}$, a closing datum for $Y$ is a function $d$ from the set of the salient blocks of $Y$ of label $c \in\{1, n-1, n\}$, and some subset of the salient blocks of $Y$ with label 0 , to the set of support blocks of $Y$, such that

- for each salient block $b$ of label $c$ for which $d$ is defined, the associated support block $d(b)$ is a support block for label $c$, and lies on the previous antidiagonal and in a lower row than that of $b$;
- for each fixed $c \in\{0,1, n-1, n\}$ the different salient blocks of label $c$ are mapped to different support blocks;
- each support block for label $c$ can support at most one salient block of label $c$ according to the supporting rules spelled out at the end of 5.1.
By condition (A2), for every $Y \in \mathcal{Z}_{\Delta}^{\prime}$ with a nonempty set of salient blocks the set $\operatorname{cd}(\mathrm{Y})$ of closing data for $Y$ is nonempty. If all salient blocks of $Y$ of label $1, n-1$ or $n$ are directly supported, then a closing datum $d \in \operatorname{cd}(\mathrm{Y})$ is called contributing, if to every salient block of label on which $d$ is defined, it associates the support block immediately below it.

We define two subsets of $\mathcal{Z}_{\Delta}^{\prime}$. Consider the following conditions for a Young wall $Y \in \mathcal{Z}_{\Delta}$.
(R1) The salient blocks of $Y$ of label $n-1$ or $n$ are all part of a directly supported salient block-pair.
(R2) $Y$ has no salient block with label in the set $\{1, \ldots, n-2\}$.
(R3) Any consecutive series of rows of $Y$ having equal length and ending in half $n-1 / n$-blocks is longer than $n-2$, or $n-1$ if the length of the rows is $n-1$, and the last one starts with a block labelled 1 (see Example 6.15 below for the latter condition being broken).
Young walls satisfying (R1)-(R2) will be called 0 -generated. Let $\mathcal{Z}_{\Delta}^{1} \subset \mathcal{Z}_{\Delta}$ denote the set of 0 generated Young walls. They automatically satisfy (A1)-(A2), so indeed $\mathcal{Z}_{\Delta}^{1} \subset \mathcal{Z}_{\Delta}^{\prime}$. Let further $\mathcal{Z}_{\Delta}^{0} \subset \mathcal{Z}_{\Delta}^{1}$ be the set of those Young walls which in addition satisfy (R3) also. These will be called distinguished.
Lemma 6.1. Let $Y \in \mathcal{Z}_{\Delta}$ be an arbitrary Young wall.
(1) The positions of the 0 blocks of $Y$ determine uniquely a Young wall $Y_{0} \in \mathcal{Z}_{\Delta}^{0}$ such that $\mathrm{wt}_{0}(Y)=\mathrm{wt}_{0}\left(Y_{0}\right)$, and the 0 blocks in $Y_{0}$ are exactly the 0 blocks in $Y$. $Y_{0}$ does not necessarily contain $Y$.
(2) There is a unique $Y_{1} \in \mathcal{Z}_{\Delta}^{1}$ that contains $Y$ and is minimal with this property with respect to containment.

We will prove Lemma 6.1 at the end of the section.

Lemma 6.2. There is a combinatorial reduction map red: $\mathcal{Z}_{\Delta}^{\prime} \rightarrow \mathcal{Z}_{\Delta}^{0}$, restricting to the identity on $\mathcal{Z}_{\Delta}^{0} \subset \mathcal{Z}_{\Delta}^{\prime}$, which associates to each $Y \in \mathcal{Z}_{\Delta}^{\prime}$ a distinguished 0 -generated Young wall $\operatorname{red}(Y) \in \mathcal{Z}_{\Delta}^{0}$ of the same 0 -weight.

Proof. Starting with a Young wall $Y \in \mathcal{Z}_{\Delta}^{\prime}$, we construct red( $Y$ ) by enforcing (R1)-(R2) and (R3) in turn, making sure in the second step that (R1)-(R2) remain fulfilled.

First, if a Young wall $Y \in \mathcal{Z}_{\Delta}^{\prime}$ violates (R1) or (R2) at a salient block of label $1, n-1$ or $n$, find the lowest row where this happens, and extend $Y$ by adding as many extra blocks as possible to this row without modifying the 0 -weight. Thus, the extension stops either just before any full block below which there is a missing full block, or just before the next 0 block in the row, whichever comes earlier. Then in the row above this, one or two blocks may become salient. If at least one of these new salient blocks is not of label 0 , then we repeat the same procedure. Following this procedure all the way to the top of $Y$ gives a new Young wall which satisfies (R1) and (R2). These moves may increase the number of places where the Young wall violates (R3).

Second, assume a Young wall $Y$ satisfies (R1) and (R2) but violates (R3): there is a consecutive series of rows having equal length and ending in half $n-1 / n$-blocks, but the length of this series is $m \leq n-2$. Remove the half block from the end of the lowest row of such a series. Then a new supported salient block-pair appears. If the block $b$ immediately above this block-pair is contained in $Y$, then we remove $b$, as well as the blocks to the right of it in order to obtain a valid Young wall. Any full block above $b$ also cannot be present in a valid Young wall, so we remove that as well, together with all of the blocks to the right of it. Continue the removal process until there is a full block above the last removed block. This process terminates after $m$ steps, when it arrives at a row which is already short enough. In this way we decreased the number places where the Young wall violates (R3), but haven't increased the number of places where it violates (R1) or (R2). The 0 -weight of the Young wall remains unchanged, since the length of the series was at most $n-2$.

Combining these steps, we obtain a Young wall red $(Y)$ that satisfies (R3) as well as (R1) and (R2), and so lies in $\mathcal{Z}_{\Delta}^{0}$, and has the same 0 -weight.

Example 6.3. Let $n=4$ and let us consider the following six Young walls.


It can be checked that all these satisfy conditions (A1) and (A2), and hence lie in $\mathcal{Z}_{\Delta}^{\prime}$. On the other hand, the Young walls $Y_{1}$ and $Y_{3}$ violate (R1), and are extended to $Y_{2}$, resp. $Y_{4}$ in the first step of the reduction process. Both of these still violate (R3); in the second step, they become $\operatorname{red}\left(Y_{1}\right)=\operatorname{red}\left(Y_{3}\right)=Y_{6} \in \mathcal{Z}_{\Delta}^{0} . Y_{5}$ satisfies (R1), but violates (R2) and is extended to $Y_{6}$ in the first step with $\operatorname{red}\left(Y_{5}\right)=Y_{6} \in \mathcal{Z}_{\Delta}^{0}$ also. In each case the bullets indicate where exactly these violations occur. The blocks without numbers are not contained in the Young walls.
For a Young wall $Y \in \mathcal{Z}_{\Delta}^{0}$, let $\operatorname{Rel}(Y)=\operatorname{red}^{-1}(Y)$ denote the set of relatives of $Y$, the Young walls we can get from $Y$ by the inverses of the reduction steps above. This is a finite directed set, directed by the steps of the proof of Lemma 6.2 In Example 6.3 all $Y_{i}$ are relatives of $Y_{6}$.
Proof of Lemma 6.1] Fix $Y \in \mathcal{Z}_{\Delta}$. For part (1), consider the minimal Young wall $Y_{m}$ that contains the same 0 -blocks as $Y$. Similarly to the proof of Lemma 6.2 extend each row of $Y_{m}$ with as many blocks as possible taking into account the Young wall rules and the conditions (R1)-(R3), and without modifying the 0 -weight.

For part (2), consider the set $\operatorname{Rel}_{Y}\left(Y_{0}\right) \subset \mathcal{Z}_{\Delta}^{\prime}$ of those relatives of $Y_{0}$ which contain $Y$. Here $Y_{0}$ is the Young wall obtained in part (1). This set of relatives is nonempty, since we can always extend $Y_{0}$ with the inverse of the move (R3) in Lemma 6.2 until there are only label 0 salient blocks. There can be several of these since there is an ambiguity in the inverse of the move (R3), but there is no Young wall having the same 0 weight as $Y_{0}$ which is not contained in at least one of these extended Young walls.
Suppose that $\operatorname{Rel}_{Y}\left(Y_{0}\right)$ has two distinct minimal elements $Y_{2}, Y_{3}$ with respect to containment. Then there is at least one row ending in a half block, where one of $Y_{2}, Y_{3}$ has a left triangle, and the other has a right triangle, but otherwise the row has the same length. Then the length of the series of successive rows with the same length is the same in the two Young walls. If this length is
more than $n-2$, then they cannot both contain $Y$. If it is $n-2$ or less, then $Y_{2}, Y_{3}$ are the two results of the inverse of the move (R3) applied on a smaller Young wall. Since $Y$ was a Young wall, also this smaller Young wall contains $Y$. Hence, neither of $Y_{2}, Y_{3}$ could be minimal. The same reasoning applies to all places where there is the left triangle/right triangle ambiguity. Thus, there is a unique minimal element $Y_{1}$ in the set of relatives of $Y_{0}$ containing $Y$.

Since $Y_{1}$ can be obtained from $Y_{0}$ by inverses of the move (R3), it is an element in $\mathcal{Z}_{\Delta}^{1}$
6.2. The decomposition of the coarse Hilbert scheme. Let us turn to the Hilbert scheme of points on the quotient $\mathbb{C}^{2} / G_{\Delta}$, the coarse Hilbert scheme $\operatorname{Hilb}\left(\mathbb{C}^{2} / G_{\Delta}\right)=\sqcup_{n} \operatorname{Hilb}^{n}\left(\mathbb{C}^{2} / G_{\Delta}\right)$. Recall that the inclusion $\mathbb{C}[x, y]^{G_{\Delta}} \subset \mathbb{C}[x, y]$ defines a morphism

$$
p_{*}: \operatorname{Hilb}\left(\left[\mathbb{C}^{2} / G_{\Delta}\right]\right) \rightarrow \operatorname{Hilb}\left(\mathbb{C}^{2} / G_{\Delta}\right), \quad J \mapsto J^{G_{\Delta}}=J \cap \mathbb{C}[x, y]^{G_{\Delta}}
$$

and a map of sets

$$
i^{*}: \operatorname{Hilb}\left(\mathbb{C}^{2} / G_{\Delta}\right)(\mathbb{C}) \rightarrow \operatorname{Hilb}\left(\left[\mathbb{C}^{2} / G_{\Delta}\right]\right)(\mathbb{C}), \quad I \mapsto \mathbb{C}[x, y] . I
$$

between the coarse and the orbifold Hilbert schemes.
The purpose of this section is to prove the following result.
Theorem 6.4. The decomposition of the equivariant Hilbert scheme $\operatorname{Hilb}\left(\left[\mathbb{C}^{2} / G_{\Delta}\right]\right)$ from Theorem 4.1 induces a locally closed decomposition

$$
\operatorname{Hilb}\left(\mathbb{C}^{2} / G_{\Delta}\right)=\bigsqcup_{Y \in \mathcal{Z}_{\Delta}^{\prime}} \operatorname{Hilb}\left(\mathbb{C}^{2} / G_{\Delta}\right)_{Y}
$$

of the coarse Hilbert scheme $\operatorname{Hilb}\left(\mathbb{C}^{2} / G_{\Delta}\right)$ into strata indexed bijectively by the set $\mathcal{Z}_{\Delta}^{\prime}$ of Young walls of type $D_{n}$ satisfying (A1)-(A2) above. The stratum $\operatorname{Hilb}\left(\mathbb{C}^{2} / G_{\Delta}\right)_{Y}$ is contained in the $m$-th Hilbert scheme $\operatorname{Hilb}^{m}\left(\mathbb{C}^{2} / G_{\Delta}\right)$ for $m=\mathrm{wt}_{0}(Y)$.

Proof. We start with the universal ideal $\mathcal{J} \triangleleft \mathcal{O}_{\mathrm{Hilb}\left(\mathbb{C}^{2} / G_{\Delta}\right)} \otimes \mathbb{C}[x, y]^{G_{\Delta}}$, which exists since $\operatorname{Hilb}\left(\mathbb{C}^{2} / G_{\Delta}\right)$ is a fine moduli space. Using the relative pullback, we obtain an invariant ideal $\mathbb{C}[x, y] . \mathcal{J} \triangleleft \mathcal{O}_{\operatorname{Hilb}\left(\mathbb{C}^{2} / G_{\Delta}\right)} \otimes \mathbb{C}[x, y]$, which however is not a flat family of invariant ideals over $\operatorname{Hilb}\left(\mathbb{C}^{2} / G_{\Delta}\right)$. Take the flattening stratification of the base with index set $F$, to obtain a decomposition

$$
\begin{equation*}
\operatorname{Hilb}\left(\mathbb{C}^{2} / G_{\Delta}\right)=\sqcup_{f \in F} \operatorname{Hilb}\left(\mathbb{C}^{2} / G_{\Delta}\right)_{f} \tag{14}
\end{equation*}
$$

over which the restrictions $(\mathbb{C}[x, y] . \mathcal{J})_{f}$ are flat. These flat families of invariant ideals of $\mathbb{C}[x, y]$ define classifying maps

$$
i_{f}: \operatorname{Hilb}\left(\mathbb{C}^{2} / G_{\Delta}\right)_{f} \rightarrow \operatorname{Hilb}^{\rho}\left(\left[\mathbb{C}^{2} / G_{\Delta}\right]\right)
$$

from these strata to components of the equivariant Hilbert scheme. The latter smooth varieties are decomposed into locally closed strata by Theorem4.1 as

$$
\begin{equation*}
\operatorname{Hilb}\left(\left[\mathbb{C}^{2} / G_{\Delta}\right]\right)=\sqcup_{Y \in \in \mathcal{Z}_{\Delta}} \operatorname{Hilb}\left(\left[\mathbb{C}^{2} / G_{\Delta}\right]\right)_{Y} \tag{15}
\end{equation*}
$$

The stratification (15) gives a stratification on $\operatorname{im}\left(i_{f}\right)$ for each $f \in F$ since over each $\operatorname{Hilb}\left(\left[\mathbb{C}^{2} / G_{\Delta}\right]\right)_{Y}$ the classifying map is flat. Hence, we can pull it back to obtain a decomposition

$$
\operatorname{Hilb}\left(\mathbb{C}^{2} / G_{\Delta}\right)=\bigsqcup_{Y \in \mathcal{Z}_{\Delta}} \operatorname{Hilb}\left(\mathbb{C}^{2} / G_{\Delta}\right)_{Y}
$$

where we have, set-theoretically,

$$
\operatorname{Hilb}\left(\mathbb{C}^{2} / G_{\Delta}\right)_{Y}(\mathbb{C})=\left\{I \in \operatorname{Hilb}\left(\mathbb{C}^{2} / G_{\Delta}\right)(\mathbb{C}): i^{*}(I) \in \operatorname{Hilb}\left(\left[\mathbb{C}^{2} / G_{\Delta}\right]\right)_{Y}(\mathbb{C})\right\}
$$

The whole construction is compatible with the $T=\mathbb{C}^{*}$-action, so we can also decompose the $T$-fixed locus representing homogeneous ideals as

$$
\operatorname{Hilb}\left(\mathbb{C}^{2} / G_{\Delta}\right)^{T}=\bigsqcup_{Y \in \mathcal{Z}_{\Delta}} W_{Y}
$$

where

$$
W_{Y}(\mathbb{C})=\left\{I \in \operatorname{Hilb}\left(\mathbb{C}^{2} / G_{\Delta}\right)(\mathbb{C}): i^{*}(I) \in Z_{Y}(\mathbb{C})\right\}
$$

Notice also that by construction, the maps $i_{\rho}$ above are given by the pullback map $i^{*}$. In other words, when restricted to the strata $\operatorname{Hilb}\left(\mathbb{C}^{2} / G_{\Delta}\right)_{Y} \supset W_{Y}$, the map $i^{*}$ becomes a morphism of
schemes. On the other hand, it is also clear that, letting $\widetilde{W}_{Y}$ denote the image of $i^{*}$ inside $Z_{Y}, p_{*}$ and $i^{*}$ are two-sided inverses and so $W_{Y} \cong \widetilde{W}_{Y} \subset Z_{Y}$.

To conclude, we need to show that for a fixed $I \in \operatorname{Hilb}\left(\mathbb{C}^{2} / G_{\Delta}\right)$, the Young wall $Y$ associated with the pullback ideal $J=i^{*}(I)$ necessarily lies in $\mathcal{Z}_{\Delta}^{\prime}$. It is clearly enough to assume that $I$, and so $J$, are homogeneous. The ideal $J$, being a pullback, is of course generated by invariant polynomials. On the other hand, as we have seen during the proof of Theorem4.1, a homogeneous ideal is generated by polynomials lying in salient blocks. While not all salient blocks necessarily contain new a generator, it is clear that salient blocks labelled $j$ with $2 \leq j \leq n-2$ must contain a generator. Since such a generator is not allowed in an invariant ideal, $Y$ must satisfy condition (A1).

To discuss the other condition (A2), let us return to the inductive proof of Theorem 4.1 Corollary 5.3 (2) implies that if there is no generator on a given antidiagonal of $Y$, then the salient blocks of label $c \in\{1, n, n-1\}$ on this antidiagonal are supported. For an invariant ideal, this condition is required to be satisfied for all salient blocks of label $c \in\{1, n, n-1\}$. This concludes the proof.
6.3. Possibly and almost invariant ideals. We wish to study the Euler characteristics of the strata of the coarse Hilbert scheme obtained in Theorem 6.4, using the inductive approach used in 4.3 in our study of the orbifold Hilbert scheme. However, as things stand, the setup does not lend itself well to induction based on the removal of the bottom row from a Young wall, since the set of Young walls $\mathcal{Z}_{\Delta}^{\prime}$ is clearly not closed under the removal of the bottom row. To remedy this, we introduce two auxiliary constructions. From now on, except when noted, every ideal is supposed to be $T$-invariant.

First, call an ideal $I \triangleleft \mathbb{C}[x, y]$ possibly invariant, if it is generated by

- polynomials which transform under $G_{\Delta}$ according to $\rho_{0}, \rho_{1}, \rho_{n-1}$ or $\rho_{n}$,
- and at most one $\tau$-invariant pair of polynomials of the same degree, forming a twodimensional representation of $G_{\Delta}$, and not lying in the image of the operator $L_{1}$.
Equivalently, the second condition says that the corresponding two-dimensional subspace lies in the large stratum of the appropriate projective space parameterizing such subspaces.

Second, a possibly invariant ideal $I$ will be called almost invariant, if it is generated by

- $G_{\Delta}$-invariant elements,
- and at most a single polynomial, or pair of polynomials of the same degree, forming a oneor two-dimensional representation of $G_{\Delta}$, and not lying in the image of the operator $L_{1}$.
Let us denote by $\mathcal{Z}_{\Delta}^{P} \subset \mathcal{Z}_{\Delta}$ the subset of all Young walls which are characterized by the following condition:
(A1') all salient blocks of $Y$ are labelled $0,1, n-1$ or $n$, except possibly for a single salient block in the bottom row of a different label.
Moreover, let $\mathcal{Z}_{\Delta}^{A} \subset \mathcal{Z}_{\Delta}^{P}$ be the set of Young walls which satisfy condition (A1') as well as the following second condition:
(A2') every salient block of $Y$ labelled $c \in\{1, n-1, n\}$ is supported, except possibly the ones in the bottom row.
The following statement follows immediately from our setup.


## Proposition 6.5.

(1) Possibly invariant ideals correspond to points in the strata $Z_{Y} \subset \operatorname{Hilb}\left(\left[\mathbb{C}^{2} / G_{\Delta}\right]\right)$ where $Y \in \mathcal{Z}_{\Delta}^{P}$.
(2) Points parameterizing almost invariant ideals lie in constructible subsets $\widetilde{W}_{Y} \subset Z_{Y}$ of strata $Z_{Y} \subset \operatorname{Hilb}\left(\left[\mathbb{C}^{2} / G_{\Delta}\right]\right)$ for $Y \in \mathcal{Z}_{\Delta}^{A}$.
By definition, we have $\mathcal{Z}_{\Delta}^{\prime} \subset \mathcal{Z}_{\Delta}^{A}$. For $Y \in \mathcal{Z}_{\Delta}^{\prime}$, the constructible subset provided by Proposition 6.5(2) is exactly $\widetilde{W}_{Y}=i_{*}\left(W_{Y}\right)$. Therefore, in the sequel we will denote these strata as $\widetilde{W}_{Y}$ for arbitrary $Y \in \mathcal{Z}_{\Delta}^{A}$ as well. Further, if for $Y \in \mathcal{Z}_{\Delta}^{A}$, also $\bar{Y} \in \mathcal{Z}_{\Delta}^{A}$, where $\bar{Y}$ is the Young wall obtained by removing from $Y$ the bottom row as in 4.3, then the map $T: Z_{Y} \rightarrow Z_{\bar{Y}}$ of Lemma4.16 takes $\widetilde{W}_{Y}$ to $\widetilde{W}_{\bar{Y}}$.

Furthermore, let $\mathcal{Z}_{\Delta}^{0, A} \subset \mathcal{Z}_{\Delta}^{A}$ be the subset defined by the following conditions:
(R1') the salient blocks of label $n-1$ or $n$ not in the bottom row are all directly supported;
(R2') there is no salient block with label in $\{1, \ldots, n-2\}$ except possibly for the bottom row;
(R3') any consecutive series of rows, except the one starting in the bottom row, having equal length and ending in half $n-1 / n$-blocks is longer than $n-2$, or $n-1$ if the length of the rows is $n-1$ and the last one starts with a block labelled 1 .

Lemma 6.6. There is a combinatorial reduction map red: $\mathcal{Z}_{\Delta}^{P} \rightarrow \mathcal{Z}_{\Delta}^{0, A}$ associating to each $Y^{\prime} \in \mathcal{Z}_{\Delta}^{P}$ a unique $Y \in \mathcal{Z}_{\Delta}^{0, A}$.
Proof. The proof of Lemma 6.2 above goes through unchanged in this setting and the reduction process gives a well-defined element in $\mathcal{Z}_{\Delta}^{P}$. After the reduction there is no indirectly supported salient block. Hence, each salient block is directly supported, and in particular supported as required by condition (A2'). Therefore, the output of the reduction is an element in $\mathcal{Z}_{\Delta}^{A}$ with the properties (R1')-(R3').

Once again, for $Y \in \mathcal{Z}_{\Delta}^{0, A}$ the Young walls $\operatorname{Rel}(Y)=\operatorname{red}^{-1}(Y)$ will be called the relatives of $Y$. The following statement is clear from the definitions.
Lemma 6.7. The sets $\mathcal{Z}_{\Delta}^{P}$ and $\mathcal{Z}_{\Delta}^{0, A}$ are closed under the operation of bottom row removal.
The sets of Young walls introduced so far can be placed in the following commutative diagram:

$$
\begin{array}{ccc}
\mathcal{Z}_{\Delta}^{0} & \hookrightarrow \mathcal{Z}_{\Delta}^{\prime} & \\
\mathfrak{f} & & \downarrow \\
\mathcal{Z}_{\Delta}^{0, A} & \hookrightarrow & \\
\mathcal{Z}_{\Delta}^{A} & \hookrightarrow & \mathcal{Z}_{\Delta}^{P} .
\end{array}
$$

The relatives of a Young wall $Y \in \mathcal{Z}_{\Delta}^{0}$ are the same in $\mathcal{Z}_{\Delta}^{\prime}$ and in $\mathcal{Z}_{\Delta}^{A}$, but there may be some new relatives in $\mathcal{Z}_{\Delta}^{P}$. This will cause no problem; see the discussion above Corollary 6.13 Even though we are interested in strata of the Young walls in the upper row, it is easier to work in the lower row because of Lemma 6.7

The notion of a closing datum generalizes word by word for ideals in the stratum of $Y \in \mathcal{Z}_{\Delta}^{A}$. For ideals in the stratum of Young walls in $\mathcal{Z}_{\Delta}^{P}$ we have to relax it, since not all salient blocks of label $1, n-1$ or $n$ are supported. A partial closing datum for a Young wall in $\mathcal{Z}_{\Delta}^{P}$ is given by associating to some of its salient blocks of label $c \in\{0,1, n-1, n\}$ a support block for label $c$ in the previous antidiagonal and below the salient block, in such a way that to each support block for label $c$ at most one salient block of label $c$ is associated. We say that those salient blocks to which there is an associated support block are closed. The set of all partial closing data for $Y \in \mathcal{Z}_{\Delta}^{P}$ will be denoted as $\operatorname{pcd}(Y)$. Closing data are special partial closing data for Young walls in $\mathcal{Z}_{\Delta}^{A}$, in which all salient blocks of label $1, n-1$ or $n$ are closed, except possibly one on the bottom row.

Fix a Young wall $Y \in \mathcal{Z}_{\Delta}^{P}$ such that in the bottom row the salient block is a support block for label $c \in\{0,1, n-1, n\}$. Let $\bar{Y}$ be the truncation of $Y$. By Lemma 6.7, $\bar{Y} \in \mathcal{Z}_{\Delta}^{P}$. If $\bar{d}$ is a partial closing datum associated to some $\bar{I} \in Z_{\bar{Y}}$, then using the decomposition of Theorem 5.2 we can extend it to a partial closing datum for each ideal in the fiber over $\bar{I}$; in particular, if $I \in Z_{Y}\left(k_{1}, l_{1}, k_{2}, l_{2}\right)$ for some pairs $\left(k_{1}, l_{1}\right)$ and $\left(k_{2}, l_{2}\right)$, and either of these is a salient block of label $c$, then we associate to them the support block in the bottom row. By induction, we obtain

Corollary 6.8. Given $Y \in \mathcal{Z}_{\Delta}^{P}$, every $I \in Z_{Y}$ defines a unique partial closing datum $d(I) \in$ $\operatorname{pcd}(Y)$.

For $d \in \operatorname{pcd}(Y)$, let $Z_{Y}(d) \subseteq Z_{Y}$ be the subset of those ideals which have partial closing datum $d$. Then

$$
Z_{Y}=\sqcup_{d \in \operatorname{pcd}(Y)} Z_{Y}(d)
$$

is a locally closed decomposition of the stratum $Z_{Y}$. Similarly, for an element $Y \in \mathcal{Z}_{\Delta}^{A}$ let $\widetilde{W}_{Y}(d) \subseteq$ $\widetilde{W}_{Y}$ be the subset of those ideals which have closing datum $d$. Then

$$
\widetilde{W}_{Y}=\sqcup_{d \in \operatorname{cd}(Y)} \widetilde{W}_{Y}(d)
$$

is a locally closed decomposition.
6.4. Euler characteristics of strata and the coarse generating series. In this section, we derive information about the topological Euler characteristics of the strata of the coarse Hilbert scheme constructed above.

Fix a Young wall $Y \in \mathcal{Z}_{\Delta}^{P}$, and a partial closing datum $d \in \operatorname{pcd}(Y)$. We say that a support block for label $c$ is of type $e \in\{-1,0,1\}$ if, when its row is considered as the bottom row, the associated half blocks according to $d$ are in the set $I\left(S_{1}, S_{2}, S\right)_{e}$ in the notation of Theorem 5.2.
Lemma 6.9. Assume that $(Y, d)$ are such such that in the bottom row of $Y$, the salient block b of label $j \in\{2, n-2\}$ is a support block for label $c \in\{0,1, n-1, n\}$. Let $(\bar{Y}, \bar{d})$ be the truncation of $(Y, d)$. If the salient block $b$ is of type $e \in\{-1,0,1\}$, then

$$
\chi\left(Z_{Y}(d)\right)=e \cdot \chi\left(Z_{\bar{Y}}(\bar{d})\right)
$$

Proof. Using the notations of 5.2, let $\left(k_{1}, l_{1}, k_{2}, l_{2}\right)$ be the quadruple of the half blocks associated to the support block, and consider the diagram

$$
\begin{array}{rlrl}
Z_{Y}(d) & \subseteq Z_{Y} \longrightarrow \mathcal{X}_{Y} \\
& & \left.\right|_{T} & \downarrow_{\bar{Y}} \\
& & \\
Z_{\bar{Y}}(\bar{d}) & \subseteq Z_{\bar{Y}} \longrightarrow \mathcal{Y}_{\bar{Y}}
\end{array}
$$

Returning once again to the process proving Theorem 4.1 when we obtained $Z_{Y}$ from the truncation $Z_{\bar{Y}}$, we saw that those ideals in $Z_{Y}$ that does not have a generator in the strata of the missing half blocks at $\left(k_{1}, l_{1}\right)$ and at $\left(k_{2}, l_{2}\right)$ are necessarily in $Z_{Y}\left(k_{1}, l_{1}, k_{2}, l_{2}\right)$ and all ideals in $Z_{Y}\left(k_{1}, l_{1}, k_{2}, l_{2}\right)$ have this property. Formally, a point of $Z_{Y}$ over $Z_{\bar{Y}}(\bar{d})$ is in $Z_{Y}(d)$ if and only if

$$
\left(\left(I \cap P_{m, j}\right) \backslash\left(I \cap \bar{P}_{m, j}\right)\right) \cap P_{m+1, c_{i}} \dashv V_{k_{i}, l_{i}, c_{i}} \text { for } i=1,2 .
$$

Under the operation $T$ the space $Z_{Y}(d)$ necessarily maps onto $Z_{\bar{Y}}(\bar{d})$. Hence, we get that

$$
Z_{Y}(d)=Z_{\bar{Y}}(\bar{d}) \times_{Z_{\bar{Y}}} Z_{Y}\left(k_{1}, l_{1}, k_{2}, l_{2}\right) .
$$

By Theorem 5.2 the fibers of $T$ on $Z_{Y}\left(k_{1}, l_{1}, k_{2}, l_{2}\right)$ have Euler characteristic $e$. Thus

$$
\chi\left(Z_{Y}(d)\right)=e \cdot \chi\left(Z_{\bar{Y}}(\bar{d})\right) .
$$

For $e \in\{-1,0,1\}$ and $c \in\{0 / 1, n-1 / n\}$ let $s_{e}(d, c)$ be the number of support blocks for label $c$ of type $e$, and let $s_{e}(d)=s_{e}(d, 0 / 1)+s_{e}(d, n-1 / n)$. Applying Lemma 6.9 inductively, we get the following.
Proposition 6.10. For $Y \in \mathcal{Z}_{\Delta}^{P}$ and $d \in \operatorname{pcd}(Y)$,

$$
\chi\left(Z_{Y}(d)\right)=(-1)^{s_{-1}(d)} \cdot 0^{s_{0}(d)} \cdot 1^{s_{1}(d)}
$$

where we adopt the convention $0^{0}=1$.
Corollary 6.11. (1) Let $Y \in \mathcal{Z}_{\Delta}^{P}$ and $d \in \operatorname{pcd}(Y)$. If $Y$ contains a salient block of any label to which a support block not immediately below it is associated under d, then $\chi\left(Z_{Y}(d)\right)=0$.
(2) Let $Y \in \mathcal{Z}_{\Delta}^{A}$ and $d \in \operatorname{cd}(Y)$.
(a) If $Y$ contains a nondirectly supported salient block of label $1, n-1$ or $n$, then $\chi\left(\widetilde{W}_{Y}(d)\right)=0$.
(b) If $Y$ does not contain any nondirectly supported salient block of label 1, $n-1$ or $n$, but $d$ is not contributing, then $\chi\left(\widetilde{W}_{Y}(d)\right)=0$.
Proof. (1) follows from Proposition 6.10 and both parts of (2) follows from (1).
The main ingredient for calculating the coarse generating series is the following statement.
Proposition 6.12. For all $Y \in \mathcal{Z}_{\Delta}^{0, A}$,

$$
\sum_{Y^{\prime} \in \operatorname{Rel}(Y)} \chi\left(\widetilde{W}_{Y^{\prime}}\right)=1
$$

Recall from 6.3 that the relatives of a Young wall $Y \in \mathcal{Z}_{\Delta}^{0}$ are the same in $\mathcal{Z}_{\Delta}^{\prime}$ and in $\mathcal{Z}_{\Delta}^{A}$, but there may be some new relatives in $\mathcal{Z}_{\Delta}^{P}$ which are not in $\mathcal{Z}_{\Delta}^{A}$. On the other hand, since $\widetilde{W}_{Y}=\emptyset$ for $Y \in \mathcal{Z}_{\Delta}^{P} \backslash \mathcal{Z}_{\Delta}^{A}$, Proposition 6.12 implies

Corollary 6.13. For all $Y \in \mathcal{Z}_{\Delta}^{0}$,

$$
\sum_{Y^{\prime} \in \operatorname{Rel}(Y)} \chi\left(\widetilde{W}_{Y^{\prime}}\right)=1
$$

Proof. By Corollary 6.11 (2.a), the strata associated to those Young walls $Y^{\prime} \in \operatorname{Rel}(Y)$ that have at least one undirectly supported salient block of label $n-1$ or $n$ above the bottom row do not contribute to the sum. Therefore we can restrict our attention to the subset of Young walls in which the salient blocks not in the bottom row are

- directly supported salient block-pairs of label $0 / 1$ or $n-1 / n$, or
- arbitrary salient blocks of label 0 .

In particular, we can assume that all Young walls we are working with satisfy (R1'). Moreover, by Corollary 6.11 (2.b) we can assume that each closing datum is contributing. That is, for each salient block of label $1, n-1$ or $n$ not in the bottom row, the associated support block is immediately below it, and this holds also for those label-0 salient blocks which have an associated support block.

As in the proof of Theorem 4.3, we will prove the proposition by induction on the number of rows. Fix a Young wall $Y \in \mathcal{Z}_{\Delta}^{0, A}$. We re-write the statement into the form

$$
\sum_{Y^{\prime} \in \operatorname{Rel}(Y)} \sum_{d^{\prime} \in \operatorname{cd}\left(Y^{\prime}\right)} \chi\left(\widetilde{W}_{Y^{\prime}}\left(d^{\prime}\right)\right)=1
$$

Let $\bar{Y}$ be the truncation of $Y$; by Lemma 6.7] $\bar{Y} \in \mathcal{Z}_{\Delta}^{0, A}$ also.
Assume first that the block above the salient block in the bottom row of $Y$ is not divided. Then the relatives of $Y$ are exactly the extensions of those of $\bar{Y}$ with the bottom row of $Y$. A closing datum on any such Young wall extends uniquely to the extended Young wall, and the corresponding strata are isomorphic. In this case, the induction step is obvious.

Assume next that the block above the salient block in the bottom row of $Y$ is divided. Let $S$ be the index set of the block at $(0, m), S_{1}$ and $S_{2}$ the index sets for the blocks at $(1, m)$. The following cases can occur:
(1) $\bar{S}$ is maximal;
(2) $\bar{S}$ is not maximal, and exactly one of $S_{1}$ or $S_{2}$ is maximal;
(3) $\bar{S}$ is not maximal, and both $S_{1}$ and $S_{2}$ are maximal.

If $\bar{S}$ is maximal, then to each relative $\bar{Y}^{\prime} \in \operatorname{Rel}(\bar{Y})$ there corresponds a unique relative $Y^{\prime} \in$ $\operatorname{Rel}(Y)$ which satisfies (R1') and (R2'): we extend each relative of $\bar{Y}$ with the bottom row of $Y$. A closing datum $\bar{d}^{\prime}$ on one of these relatives $\bar{Y}^{\prime}$ can be extended to a closing datum $d^{\prime}$ on the corresponding relative of $Y$ by assigning $(\emptyset, \emptyset)$ to the new salient block appearing in the bottom row. By Theorem 5.2, $(\emptyset, \emptyset) \in I\left(S, S_{1}, S_{2}\right)_{1}$, and

$$
\sum_{Y^{\prime} \in \operatorname{Rel}(Y)} \sum_{d^{\prime} \in \operatorname{cd}\left(Y^{\prime}\right)} \chi\left(\widetilde{W}_{Y^{\prime}}\left(d^{\prime}\right)\right)=\sum_{\bar{Y}^{\prime} \in \operatorname{Rel}(\bar{Y})} \sum_{\bar{d}^{\prime} \in \operatorname{cd}\left(\bar{Y}^{\prime}\right)} \chi\left(\widetilde{W}_{\bar{Y}^{\prime}}\left(\bar{d}^{\prime}\right)\right)=1
$$

If $\bar{S}$ is not maximal, and $S_{i}$ is maximal while $S_{3-i}$ is not, then the missing block at $(1, m)$ is not salient, and the subspace in its stratum is necessarily in the image of the subspace at $(0, m)$. Again, to each relative of $\bar{Y}$ there corresponds a unique relative of $Y$, and

$$
\sum_{Y^{\prime} \in \operatorname{Rel}(Y)} \sum_{d^{\prime} \in \operatorname{cd}\left(Y^{\prime}\right)} \chi\left(\widetilde{W}_{Y^{\prime}}\left(d^{\prime}\right)\right)=\sum_{\bar{Y}^{\prime} \in \operatorname{Rel}(\bar{Y})} \sum_{\bar{d}^{\prime} \in \operatorname{cd}\left(\bar{Y}^{\prime}\right)} \chi\left(\widetilde{W}_{\bar{Y}^{\prime}}\left(\bar{d}^{\prime}\right)\right)=1 .
$$

Finally, suppose that $\bar{S}$ is not maximal, but $S_{1}$ and $S_{2}$ are maximal. In this case new relatives appear when we add the bottom row of $Y$ to $\bar{Y}$. We investigate two possible cases individually.

If the divided block above the new salient block is a full $n-1 / n$ block, then for each relative of $\bar{Y}$ there are two other relatives, which contain either the label $(n-1)$ half block or the label $n$ half block. If the relatives of $\bar{Y}$ are $\left\{\bar{Y}_{0}=\bar{Y}, \bar{Y}_{1}, \bar{Y}_{2}, \ldots\right\}$, then the relatives of $Y$ are

$$
\left\{Y_{0}=Y, Y_{1}, Y_{2}, \ldots\right\} \cup \bigcup_{c \in\{n-1, n\}}\left\{Y_{0, c}=Y_{c}, Y_{1, c}, Y_{2, c}, \ldots\right\}
$$

We obtain these by performing the inverse of the move (R3) in the algorithm of Lemma 6.2, The addition of the bottom row of $Y$ to an $\overline{Y_{i}}$ schematically looks like this:


Here the block of label $c$ is in the complement of the Young wall $Y_{i, c}$, while its pair is in $Y_{i, c}$. There is only one contributing closing datum on each of $Y_{i}, Y_{i, n-1}$ and $Y_{i, n}$ extending a contributing closing datum $\bar{d}_{i}$ on $\bar{Y}_{i}$. Namely, to the new support block of $Y_{i}$ we associate both the divided block above it, and every other part of $\bar{d}_{i}$ is kept constant. We denote these by $d_{i}, d_{i, n-1}$ and $d_{i, n}$, respectively.

We claim that for $c \in\{n-1, n\}$,

$$
\chi\left(\widetilde{W}_{Y_{i, c}}\left(d_{i, c}\right)\right)=\chi\left(\widetilde{W}_{\bar{Y}_{i}}\left(\bar{d}_{i}\right)\right)
$$

To show this, we define a morphism $\widetilde{W}_{\bar{Y}_{i}}\left(\bar{d}_{i}\right) \rightarrow \widetilde{W}_{\bar{Y}_{i, c}}\left(\bar{d}_{i, c}\right)$, where $\bar{Y}_{i, c}$ is the truncation of $Y_{i, c}$. This morphism is given by the restriction of an ideal $I \in \widetilde{W}_{\bar{Y}_{i}}\left(\bar{d}_{i}\right)$ to the union of those cells whose block is missing from $\bar{Y}_{i, c}$. The Young wall $\bar{Y}_{i, c}$ has the same salient blocks as $\bar{Y}_{i}$ except a half block in the bottom row. Hence, the result is necessarily an ideal, which has the same generators as $I$ except for the cell $V_{0, m, c}$ where it does not have any generator. Therefore, the image of this morphism is in $W_{\bar{Y}_{i, c}}\left(\bar{d}_{i, c}\right)$; here $(0, m)$ is the position of the salient block-pair in the bottom row of $\bar{Y}_{i}$, which is also the position of the salient half block of label $\kappa(c)$ in the bottom row of $\bar{Y}_{i, c}$.

Assume that there are two ideals $I, I^{\prime} \in \widetilde{W}_{\bar{Y}_{i}}\left(\bar{d}_{i}\right)$ which map to the same element of $\widetilde{W}_{\bar{Y}_{i, c}}\left(\bar{d}_{i, c}\right)$ under this morphism. Then they only differ in the function at $V_{0, m, c}$, or more precisely in the point of $V_{0, m, c} / \bar{U}$ which represents the subspace in $V_{0, m, c}$. Here $\bar{U}=I \cap \bar{P}_{m, c}=I^{\prime} \cap \bar{P}_{m, c}$. Then any ideal $I^{\prime \prime}$ which is their affine linear combination, i.e. whose corresponding subspace in $V_{0, m, c}$ is a linear combination of those of $I$ and $I^{\prime}$, is also an element of $\widetilde{W}_{\bar{Y}_{i}}\left(\bar{d}_{i}\right)$, mapped to the same ideal as $I$ and $I^{\prime}$ under the morphism. In particular, the fibers of the morphism are affine spaces. Taking into account that $\chi\left(\widetilde{W}_{Y_{i, c}}\left(d_{i, c}\right)\right)=\chi\left(\widetilde{W}_{\bar{Y}_{i, c}}\left(\bar{d}_{i, c}\right)\right)$, this proves the claim.

Thus,

$$
\begin{gathered}
\chi\left(\widetilde{W}_{Y_{i}}\left(d_{i}\right)\right)+\chi\left(\widetilde{W}_{Y_{i, n-1}}\left(d_{i, n-1}\right)\right)+\chi\left(\widetilde{W}_{Y_{i, n}}\left(d_{i, n}\right)\right) \\
=-\chi\left(\widetilde{W}_{\bar{Y}_{i}}\left(\bar{d}_{i}\right)\right)+\chi\left(\widetilde{W}_{\bar{Y}_{i}}\left(\bar{d}_{i}\right)\right)+\chi\left(\widetilde{W}_{\bar{Y}_{i}}\left(\bar{d}_{i}\right)\right)=\chi\left(\widetilde{W}_{\bar{Y}_{i}}\left(\bar{d}_{i}\right)\right),
\end{gathered}
$$

where in the first equality we used Lemma 6.9
Second, if the divided block above the new salient block is a full $0 / 1$ block, then for each relative of $\bar{Y}$ there is a new relative which contains the label 1 half block (and possibly some other blocks above it). Let us denote the relatives of $\bar{Y}$ as $\left\{\bar{Y}_{0}=\bar{Y}, \bar{Y}_{1}, \bar{Y}_{2}, \ldots\right\}$. Then the relatives of $Y$ are

$$
\left\{Y_{0}=Y, Y_{1}, Y_{2}, \ldots\right\} \cup\left\{Y_{0,1}, Y_{1,1}, Y_{2,1}, \ldots\right\}
$$

We obtain these by performing the inverse of the move (R2) in the algorithm of Lemma 6.2. Schematically:


In this case $Y_{i, 0}$ cannot appear as a relative, since that would change the 0 -weight. Given a contributing closing datum $\bar{d}_{i}$ on $\bar{Y}_{i}$, shifting it appropriately gives a unique contributing closing datum $d_{i, 1}$ on $Y_{i, 1}$. On $Y_{i}, \bar{d}_{i}$ induces two contributing closing data:

- $d_{i}$ is obtained by associating the support block in the bottom row to the salient half-block pair above it;
- $d_{i}^{\prime}$ is obtained by associating the support block in the bottom row to the salient half block of label 1 above it only.
Again, using Lemma 6.9 we get

$$
\begin{aligned}
\chi\left(\widetilde{W}_{Y_{i}}\left(d_{i}\right)\right)+\chi\left(\widetilde{W}_{Y_{i}}\left(d_{i}^{\prime}\right)\right)+\chi\left(\widetilde{W}_{Y_{i, 1}}\left(d_{i, 1}\right)\right) & =-\chi\left(\widetilde{W}_{\bar{Y}_{i}}\left(\bar{d}_{i}\right)\right)+\chi\left(\widetilde{W}_{\bar{Y}_{i}}\left(\bar{d}_{i}\right)\right)+\chi\left(\widetilde{W}_{\bar{Y}_{i}}\left(\bar{d}_{i}\right)\right) \\
& =\chi\left(\widetilde{W}_{\bar{Y}_{i}}\left(\bar{d}_{i}\right)\right) .
\end{aligned}
$$

Summing over these, we obtain in all cases

$$
\sum_{Y^{\prime} \in \operatorname{Rel}(Y)} \sum_{d^{\prime} \in \operatorname{cd}\left(Y^{\prime}\right)} \chi\left(\widetilde{W}_{Y^{\prime}}(d)\right)=\sum_{\bar{Y}^{\prime} \in \operatorname{Rel}(\bar{Y})} \sum_{\bar{d}^{\prime} \in \operatorname{cd}\left(\bar{Y}^{\prime}\right)} \chi\left(\widetilde{W}_{\bar{Y}^{\prime}}(\bar{d})\right),
$$

which proves the induction step.
Putting everything together, we obtain the following result, the analogue of Corollary 2.3 in type $D$.
Theorem 6.14. Let $\Delta$ be of type $D_{n}$. Then the generating function of Euler characteristics of the coarse Hilbert schemes of points of the corresponding singular surface $\mathbb{C}^{2} / G_{\Delta}$ is given by the combinatorial generating series

$$
\sum_{m=0}^{\infty} \chi\left(\operatorname{Hilb}^{m}\left(\mathbb{C}^{2} / G_{\Delta}\right)\right) q^{m}=\sum_{Y \in \mathcal{Z}_{\Delta}^{0}} q^{\mathrm{wt}_{0}(Y)}
$$

Proof. We use the decomposition of Theorem 6.4. For $Y \in \mathcal{Z}_{\Delta}^{\prime}$, we have $\operatorname{Hilb}\left(\mathbb{C}^{2} / G_{\Delta}\right)_{Y}^{T}=W_{Y}$ and thus $\chi\left(W_{Y}\right)=\chi\left(\operatorname{Hilb}\left(\mathbb{C}^{2} / G_{\Delta}\right)_{Y}\right)$. Now use Corollary 6.13 to sum the Euler characteristics of the strata $W_{Y}$ along the fibres of the combinatorial reduction map red of Lemma 6.2,

Example 6.15. Let $n=4$ and let $Y \in \mathcal{Z}_{\Delta}^{0}$ be the distinguished 0 -generated Young wall


The parameter space $Z_{Y}$ of this Young wall $Y$ is isomorphic to that of Example 4.6 which in turn is isomorphic to that of Example 4.5. In particular, $Z_{Y} \cong \mathbb{A}^{1}$. The difference is that in this case the salient blocks are the 0-labelled blocks at $(0,4)$ and $(1,5)$.

Denote by $Y_{3}$ and $Y_{4}$ the 0-generated non-distinguished Young walls


We have $Y_{3}, Y_{4} \in \mathcal{Z}_{\Delta}^{\prime}$ and are in fact 0-generated, but both violate condition (R3) at their fourth row, so they are not distinguished. In the first step of the reduction algorithm of Lemma 6.2 we remove the half block of label 3 (resp., 4) from $Y_{3}$ (resp, $Y_{4}$ ). Then we remove the blocks of label 2 and 4 (resp., 2 and 3 ) from the fifth row. Finally we remove the blocks of 1,2 and 3 (reps., 1,2 and 4) from the sixth row. This shows that both Young walls $Y_{3}, Y_{4}$ reduce to $Y$. It is can be seen that these are in fact all the relatives of $Y$. As explained in Example 4.6, when the two generators $\left(v_{0}, v_{0}^{\prime}\right) \in V_{0,4,0} \times V_{1,5,0}$ are in a special position such that $L_{1} v_{0}, v_{0}^{\prime} \in\left(L_{1} L_{3}\right)$, then the ideal $\left(v_{0}, v_{0}^{\prime}\right)$ has Young wall $Y_{3}$. Similarly, if $L_{1} v_{0}, v_{0}^{\prime} \in\left(L_{1} L_{4}\right)$, then $\left(v_{0}, v_{0}^{\prime}\right)$ has Young wall $Y_{4}$. In fact we can think of the corresponding strata as gluing to one stratum inside the invariant Hilbert scheme $\operatorname{Hilb}^{3}\left(\mathbb{C}^{2} / D_{4}\right)$, the strata of $Y_{3}$ and $Y_{4}$ "patching in" the gaps of the stratum of $Y$. At least on the level of Euler characteristics, this is what Proposition 6.12 shows in full generality.

## 7. Type $D_{n}$ : abacus Combinatorics

7.1. Young walls and abacus of type $D_{n}$. We continue to work with the root system $\Delta$ of type $D_{n}$, and introduce some associated combinatorics. From now on, we return to the untransformed representation of the type $D_{n}$ Young wall pattern introduced in 3.2,

Recalling the Young wall rules (YW1)-(YW4), it is clear that every $Y \in \mathcal{Z}_{\Delta}$ can be decomposed as $Y=Y_{1} \sqcup Y_{2}$, where $Y_{1} \in \mathcal{Z}_{\Delta}$ has full columns only, and $Y_{2} \in \mathcal{Z}_{\Delta}$ has all its columns ending in a half-block. These conditions define two subsets $\mathcal{Z}_{\Delta}^{f}, \mathcal{Z}_{\Delta}^{h} \subset \mathcal{Z}_{\Delta}$. Because of the Young wall rules, the pair $\left(Y_{1}, Y_{2}\right)$ uniquely reconstructs $Y$, so we get a bijection

$$
\begin{equation*}
\mathcal{Z}_{\Delta} \longleftrightarrow \mathcal{Z}_{\Delta}^{f} \times \mathcal{Z}_{\Delta}^{h} \tag{16}
\end{equation*}
$$

Given a Young wall $Y \in \mathcal{Z}_{\Delta}$ of type $D_{n}$, let $\lambda_{k}$ denote the number of blocks (full or half blocks both contributing 1) in the $k$-th vertical column. By the rules of Young walls, the resulting positive integers $\left\{\lambda_{1}, \ldots, \lambda_{r}\right\}$ form a partition $\lambda=\lambda(Y)$ of weight equal to the total weight $|Y|$, with the additional property that its parts $\lambda_{k}$ are distinct except when $\lambda_{k} \equiv 0 \bmod (n-1)$. Corresponding to the decomposition (16), we get a decomposition $\lambda(Y)=\mu(Y) \sqcup \nu(Y)$. In $\mu(Y)$, no part is congruent to 0 modulo $(n-1)$, and there are no repetitions; all parts in $\nu(Y)$ are congruent to 0 modulo $(n-1)$ and repetitions are allowed. Note that the pair $(\mu(Y), \nu(Y))$ does almost, but not quite, encode $Y$, because of the ambiguity in the labels of half-blocks on top of non-complete columns.

We now introduce the typ $6{ }^{5} D_{n}$ abacus, following [25, 27]. This abacus is the arrangement of positive integers, called positions, in the following pattern.

$$
\begin{array}{cccccccc}
1 & \ldots & n-2 & n-1 & n & \ldots & 2 n-3 & 2 n-2 \\
2 n-1 & \ldots & 3 n-4 & 3 n-3 & 3 n-2 & \ldots & 4 n-5 & 4 n-4 \\
\vdots & & \vdots & \vdots & \vdots & & \vdots & \vdots
\end{array}
$$

For any integer $1 \leq k \leq 2 n-2$, the set of positions in the $k$-th column of the abacus is the $k$-th runner, denoted $R_{k}$. As in type $A$, positions on the runners are occupied by beads. For $k \not \equiv 0 \bmod (n-1)$, the runners $R_{k}$ can only contain normal (uncolored) beads, with each position occupied by at most one bead. On the runners $R_{n-1}$ and $R_{2 n-2}$, the beads are colored white and black. An arbitrary number of white or black beads can be put on each such position, but each position can only contain beads of the same color.

[^4]Given a type $D_{n}$ Young wall $Y \in \mathcal{Z}_{\Delta}$, let $\lambda=\mu \sqcup \nu$ be the corresponding partition with its decomposition. For each nonzero part $\nu_{k}$ of $\nu$, set

$$
n_{k}=\#\left\{1 \leq j \leq l(\mu) \mid \mu_{j}<\nu_{k}\right\}
$$

to be the number of full columns shorter than a given non-full column. The abacus configuration of the Young wall $Y$ is defined to be the set of beads placed at positions $\lambda_{1}, \ldots, \lambda_{r}$. The beads at positions $\lambda_{k}=\mu_{j}$ are uncolored; the color of the bead at position $\lambda_{k}=\nu_{l}$ corresponding to a column $C$ of $Y$ is

$$
\begin{cases}\text { white, } & \text { if the block at the top of } C \text { is } \triangle \text { and } n_{l} \text { is even; } \\ \text { white, } & \text { if the block at the top of } C \text { is } \nabla \text { and } n_{l} \text { is odd; } \\ \text { black, } & \text { if the block at the top of } C \text { is } \nabla \text { and } n_{l} \text { is even; } \\ \text { black, } & \text { if the block at the top of } C \text { is } \Delta \text { and } n_{l} \text { is odd. }\end{cases}
$$

One can check that the abacus rules are satisfied, that all abacus configurations satisfying the above rules, with finitely many uncolored, black and white beads, can arise, and that the Young wall $Y$ is uniquely determined by its abacus configuration.

Example 7.1. The abacus configuration associated to the Young wall $Y_{6}$ of Example 6.3 is

| $R_{1}$ | $R_{2}$ | $R_{3}$ | $R_{4}$ | $R_{5}$ | $R_{6}$ |
| :---: | :---: | :---: | :---: | :---: | :---: |
| 1 | 2 | 3 | 4 | 5 | $(6$ |
| 7 | 8 | 9 | 10 | $(11$ | $(12)$ |
| 13 | 14 | 15 | 16 | 17 | 18 |
|  | $\vdots$ |  |  | $\vdots$ |  |

7.2. Core Young walls and their abacus representation. In parallel with the type $A$ story, we now introduce the combinatorics of core Young walls of type $D_{n}$, and the corresponding abacus moves. On the Young wall side, define a bar to be a connected set of blocks and half-blocks, with each half-block occuring once and each block occuring twice. A Young wall $Y \in \mathcal{Z}_{\Delta}$ will be called a core Young wall, if no bar can be removed from it without violating the Young wall rules. For an example of bar removal, see [25, Example 5.1(2)]. Let $\mathcal{C}_{\Delta} \subset \mathcal{Z}_{\Delta}$ denote the set of all core Young walls. The following statement is the type $D$ analogue of the discussion of 2.3,
Proposition 7.2. (25, 27]) Given a Young wall $Y \in \mathcal{Z}_{\Delta}$, any complete sequence of bar removals through Young walls results in the same core $\operatorname{core}(Y) \in \mathcal{C}_{\Delta}$, defining a map of sets

$$
\text { core: } \mathcal{Z}_{\Delta} \rightarrow \mathcal{C}_{\Delta}
$$

The process can be described on the abacus, respects the decomposition (16), and results in a bijection

$$
\begin{equation*}
\mathcal{Z}_{\Delta} \longleftrightarrow \mathcal{P}^{n+1} \times \mathcal{C}_{\Delta} \tag{17}
\end{equation*}
$$

where $\mathcal{P}$ is the set of ordinary partitions. Finally, there is also a bijection

$$
\begin{equation*}
\mathcal{C}_{\Delta} \longleftrightarrow \mathbb{Z}^{n} \tag{18}
\end{equation*}
$$

Proof. Decompose $Y$ into a pair of Young walls $\left(Y_{1}, Y_{2}\right)$ as above. Let us first consider $Y_{1}$. On the corresponding runners $R_{k}, k \not \equiv 0 \bmod (n-1)$, the following steps correspond to bar removals 25, Lemma 5.2].
(B1) If $b$ is a bead at position $s>2 n-2$, and there is no bead at position $(s-2 n+2)$, then move $b$ one position up and switch the color of the beads at positions $k$ with $k \equiv 0 \bmod (n-1)$, $s-2 n+2<k<s$.
(B2) If $b$ and $b^{\prime}$ are beads at position $s$ and $2 n-2-s(1 \leq s \leq n-2)$ respectively, then remove $b$ and $b^{\prime}$ and switch the color of the beads at positions $k \equiv 0 \bmod (n-1), s<k<2 n-2-s$.
Performing these steps as long as possible results in a configuration of beads on the runners $R_{k}$ with $k \not \equiv 0 \bmod (n-1)$ with no gaps from above, and for $1 \leq s \leq n-2$, beads on only one of $R_{s}, R_{2 n-2-s}$. This final configuration can be uniquely described by an ordered set of integers $\left\{z_{1}, \ldots, z_{n-2}\right\}, z_{s}$ being the number of beads on $R_{s}$ minus the number of beads on $R_{2 n-2-s}[27$, Remark $3.10(2)$ ]. In the correspondence (18) this gives $\mathbb{Z}^{n-2}$. It turns out that the reduction steps
in this part of the algorithm can be encoded by an $(n-2)$-tuple of ordinary partitions, with the summed weight of these partitions equal to the number of bars removed 25, Theorem 5.11(2)].

Let us turn to $Y_{2}$, represented on the runners $R_{k}, k \equiv 0 \bmod (n-1)$. On these runners, the following steps correspond to bar removals [27, Sections 3.2 and 3.3].
(B3) Let $b$ be a bead at position $s \geq 2 n-2$. If there is no bead at position $(s-n+1)$, and the beads at position $(s-2 n+2)$ are of the same color as $b$, then shift $b$ up to position $(s-2 n+2)$.
(B4) If $b$ and $b^{\prime}$ are beads at position $s \geq n-1$, then move them up to position $(s-n+1)$. If $s-n+1>0$ and this position already contains beads, then $b$ and $b^{\prime}$ take that same color. During these steps, there is a boundary condition: there is an imaginary position 0 in the rightmost column, which is considered to contain unvisible white beads; placing a bead there means that this bead disappears from the abacus. It turns out that the reduction steps in this part of the algorithm can be described by a triple of ordinary partitions, again with the summed weight of these partitions equal to the number of bars removed [27, Proposition 3.6]. On the other hand, the final result can be encoded by a pair of 2 -core partitions (see 2.3).

The different bar removal steps (B1)-(B4) construct the map $c$ algorithmically and uniquely. The stated facts about parameterizing the steps prove the existence of the bijection (17). To complete the proof of (18), we only need to remark further that the set of 2 -core partitions, in our language $A_{1}$-core partitions, is in bijection with the set of integers by bijection (7) in 2.3 (see also [27, Remark $3.10(2)]$ ). This gives the remaining $\mathbb{Z}^{2}$ factor in the bijection (18).

Example 7.3. The abacus configuration associated to the core of the Young wall $Y_{6}$ of Examples 6.3 and 7.1 is

| $R_{1}$ | $R_{2}$ | $R_{3}$ | $R_{4}$ | $R_{5}$ | $R_{6}$ |
| :---: | :---: | :---: | :---: | :---: | :---: |
| 1 | 2 | 3 | 4 | 5 | $(6)$ |
| 7 | 8 | 9 | 10 | 11 | $(12)$ |
|  | $\vdots$ |  | $\vdots$ |  |  |

We next determine the multi-weight of a Young wall $Y$ in terms of the bijections (17)-(18). The quotient part is easy: the multi-weight of each bar is $(1,1,2, \ldots, 2,1,1)$, so in complete analogy with the type $A$ situation, the contribution of the $(n+1)$-tuple of partitions to the multi-weight is easy to compute. Turning to cores, under the bijection $\mathcal{C}_{\Delta} \leftrightarrow \mathbb{Z}^{n}$, the total weight of a core Young wall $Y \in \mathcal{C}_{\Delta}$ corresponding to $\left(z_{1}, \ldots, z_{n}\right) \in \mathbb{Z}^{n}$ is calculated in [27, Remark 3.10]:

$$
\begin{equation*}
|Y|=\frac{1}{2} \sum_{i=1}^{n-2}\left((2 n-2) z_{i}^{2}-(2 n-2 i-2) z_{i}\right)+(n-1) \sum_{i=n-1}^{n}\left(2 z_{i}^{2}+z_{i}\right) . \tag{19}
\end{equation*}
$$

A refinement of this formula calculates the multi-weight of $Y$.
Theorem 7.4. Let $q=q_{0} q_{1} q_{2}^{2} \ldots q_{n-2}^{2} q_{n-1} q_{n}$, corresponding to a single bar.
(1) Composing the bijection (18) with an appropriate $\mathbb{Z}$-change of coordinates in the lattice $\mathbb{Z}^{n}$, the multi-weight of a core Young wall $Y \in \mathcal{C}_{\Delta}$ corresponding to an element $\bar{m}=\left(m_{1}, \ldots, m_{n}\right) \in \mathbb{Z}^{n}$ can be expressed as

$$
q_{1}^{m_{1}} \cdots \cdots q_{n}^{m_{n}}\left(q^{1 / 2}\right)^{\bar{m}^{\top} \cdot C \cdot \bar{m}}
$$

where $C$ is the Cartan matrix of type $D_{n}$.
(2) The multi-weight generating series

$$
Z_{\Delta}\left(q_{0}, \ldots, q_{n}\right)=\sum_{Y \in \mathcal{Z}_{\Delta}} \underline{q}^{\underline{w}(Y)}
$$

of Young walls for $\Delta$ of type $D_{n}$ can be written as

$$
Z_{\Delta}\left(q_{0}, \ldots, q_{n}\right)=\frac{\sum_{\bar{m}=\left(m_{1}, \ldots, m_{n}\right) \in \mathbb{Z}^{n}}^{\infty} q_{1}^{m_{1}} \cdots \cdot q_{n}^{m_{n}}\left(q^{1 / 2}\right)^{\bar{m}^{\top} \cdot C \cdot \bar{m}}}{\prod_{m=1}^{\infty}\left(1-q^{m}\right)^{n+1}}
$$

(3) The following identity is satisfied between the coordinates $\left(m_{1}, \ldots, m_{n}\right)$ and $\left(z_{1}, \ldots, z_{n}\right)$ on $\mathbb{Z}^{n}$ :

$$
\sum_{i=1}^{n} m_{i}=-\sum_{i=1}^{n-2}(n-1-i) z_{i}-(n-1) c\left(z_{n-1}+z_{n}\right)-(n-1) b
$$

Here $z_{1}+\cdots+z_{n-2}=2 a-b$ for integers $a \in \mathbb{Z}, b \in\{0,1\}$, and $c=2 b-1 \in\{-1,1\}$.
Proof. The coordinate change $\left(z_{1}, \ldots, z_{n}\right) \mapsto\left(m_{1}, \ldots, m_{n}\right)$ and the multiweight formula of (1), as well as (3), follow from somewhat involved but routine calculations which we omit; details can be found in the thesis 20] of the first named author. (2) clearly follows from (1) and the preceeding discussion.
7.3. 0-generated Young walls and their abacus representations. In this section, we characterize the abacus configurations corresponding to Young walls in the sets $\mathcal{Z}_{\Delta}^{0}$ and $\mathcal{Z}_{\Delta}^{1}$ defined in 6.1

Recall conditions (R1)-(R3) on Young walls $Y \in \mathcal{Z}_{\Delta}$ from 6.1. Recall also that a Young wall corresponds uniquely to an abacus configuration where the beads are placed at the positions $\lambda_{1}, \ldots, \lambda_{r}$. Finally recall that the quantity $n_{k}$ denotes the number of full columns shorter than a given non-full column of height $\lambda_{k}$.

Lemma 7.5. Conditions (R1)-(R2) on Young walls are equivalent to the following conditions for an abacus configuration.
(D1) In each row, the rightmost bead is always on the $(2 n-2)-n d$ ruler, and either all the beads of the row are at this position, or the number of beads in this position is odd.
(D2) If a row ends with a white (resp., black) bead corresponding to a column of height $\lambda_{k}$, then $k+n_{k}$ must be odd (resp. even). If several beads are placed on this position, which is allowed since it is on the ruler $R_{2 n-2}$, then this condition refers to the smallest possible $k$.
(D3) The total number of beads in the whole abacus is even, or the total number of beads on the rulers $R_{1}, \ldots, R_{n-1}$ in the first row is $n-2$.
(D4) The beads on the rulers $R_{n} \ldots, R_{2 n-3}$ are pushed to the right as much as possible in each row. In any given row, the positions on the rulers $R_{1}, \ldots, R_{n-1}$ are empty unless all the positions on the rulers $R_{n}, \ldots, R_{2 n-2}$ are filled.
(D5) The beads on the rulers $R_{1} \ldots R_{n-1}$ in any given row are either all on the ruler $R_{n-1}$, or on the rules $R_{1} \ldots R_{n-2}$, and pushed to the right as much as possible.
Condition (R3) is equivalent to the following condition.
(D6) Let $s$ be the total number of beads on the rulers $R_{1}, \ldots, R_{n-1}$ in any given row.
(a) If $s>n-2$, then all these beads are on $R_{n-1}$.
(b) If $s \leq n-2$, then all these beads are on the rulers $R_{1}, \ldots, R_{n-2}$, pushed to the right.

Thus 0 -generated Young walls $Y \in \mathcal{Z}_{\Delta}^{1}$ correspond to abacus configurations satisfying (D1)-(D5), whereas distinguished 0 -generated Young walls $Y \in \mathcal{Z}_{\Delta}^{0}$ correspond to those satisfying (D6) also.
Proof. Two kinds of salient blocks can appear in a Young wall that satisfies (R1)-(R2):

- label 0 half-blocks,
- and salient block-pairs of label $n-1 / n$.

As in the type $A$ case a salient block corresponds to the first bead in a consecutive series of beads in the abacus. More precisely, if there are several columns of height $\lambda_{k}$, or equivalently, if there are several beads placed on the position $\lambda_{k}$, then the salient block corresponds to that column of height $\lambda_{k}$ which has the smallest possible index $k$ among these.

The label 0 blocks always correspond to positions which are on the ruler $R_{2 n-2}$. In the odd columns of the type $D$ pattern they are of the shape $\nabla$ while in the even columns they are of the shape $\Delta$. Condition (D2) encodes the fact, that the salient blocks of label 0 are upper triangles in odd columns and lower triangles in even columns, and that the color of the beads corresponding to them is also affected by the parity of the appropriate $n_{k}$.

If there is a salient block of label 0 , then some columns of the same height, let's say, $\lambda_{k}$, can follow it. If the first column after them has height $\lambda_{k}-1$ then on its top there is again a salient block of label 0 . This block can only have the opposite orientation than the aforementioned salient
block, hence the number of columns of height $\lambda_{k}$ in this case can only be odd. This gives condition (D1).

Condition (D3) follows again from the absence of label 1 salient blocks. To see this recall that in the bottom row of the type $D$ pattern there are half blocks which have label 0 in the odd columns and have label 1 in the even columns. Since there are no salient blocks of label 1 in $Y$, this total number of columns can only be odd if the last column has height 1 , the column to the left of it has height 2 , and so on. This is can only happen when in the bottom row $s=n-2$.

The fact that there is no salient block of label $2, \ldots, n-2$ implies that for each bead on the rulers $R_{1}, \ldots, R_{2 n-1}$ there has to be a block placed on its right. The only exception is the ruler $R_{n-2}$. There cannot be any bead on this ruler, except when there is a salient block pair of label $n-1 / n$ which corresponds to a hole on $R_{n-1}$. These considerations imply conditions (D4) and (D5).

Condition (D6) corresponds directly to property (R3).
Given $Y \in \mathcal{Z}_{\Delta}^{0}$, let $t_{i}$ denote the total number of beads in the $i$-th row of its abacus representation, and $l_{i}$ the number of beads in the rightmost position of the $i$-th row. We obtain a sequence of pairs $\left(t_{i}, l_{i}\right)_{i \in \mathbb{Z}_{+}}$, only finitely many of which do not equal ( 0,0 ).
Corollary 7.6. Given $Y \in \mathcal{Z}_{\Delta}^{0}$, the associated sequence of pairs $\left(t_{i}, l_{i}\right)_{i \in \mathbb{Z}_{+}}$satisfies the following conditions.
(1) For all $i, 0 \leq l_{i} \leq t_{i}$.
(2) For all $i$, if $t_{i}>0$, then either $l_{i}=t_{i}$ is even, or $l_{i}$ is odd.
(3) Either $\sum_{i} t_{i}$ is even, or $t_{1}-l_{1}=2 n-4$.

Conversely, any sequence $\left(t_{i}, l_{i}\right)_{i \in \mathbb{Z}_{+}}$satisfying these conditions arises as a sequence associated to at least one Young wall $Y \in \mathcal{Z}_{\Delta}^{0}$. More precisely, the number of different Young walls $Y \in \mathcal{Z}_{\Delta}^{0}$ corresponding to any given sequence is $2^{m}$, where $m$ is the number of indices $i$ such that $t_{i}-l_{i}>$ $2 n-2$. All Young walls $Y$ corresponding to a single sequence have the same multi-weight, when the weights for labels $n-1$ and $n$ are counted together.

Proof. Condition (1) is clear from the definition of $\left(t_{i}, l_{i}\right)$. Condition (2) follows from (D1) above. Condition (3) is equivalent to condition (D3).

Conversely, given a sequence $\left(t_{i}, l_{i}\right)_{i \in \mathbb{Z}_{+}}$satisfying conditions (1)-(3), we can reconstruct a corresponding $Y \in \mathcal{Z}_{\Delta}^{0}$ in its abacus representation as follows. On the $i$-th row, $l_{i}$ beads have to be put on the last position; (D1) is satisfied because of (1). They are white if $\sum_{j<i} t_{j}+\sum_{j>i, t_{j} \neq 0 \bmod n-1} t_{j} \equiv$ $1 \bmod 2$, and black otherwise; this is just a reformulation of (D2). (D3) is satisfied because of (3). At most one bead can be put on each ruler between $R_{n}$ and $R_{2 n-3}$, pushed to the right as much as possible; this is (D4). If $t_{i}-l_{i} \leq 2 n-2$, then the rest of the beads fill up the rulers between $R_{1}$ and $R_{n-2}$, pushed to the right. If $t_{i}-l_{i}>2 n-2$, then there are no beads in this row on the rulers between $R_{1}$ and $R_{n-2}$; instead, the remaining beads are all placed on the ( $n-1$ )-st ruler, and they can be either white or black. These rules reconstruct a configuration satisfying (D5)-(D6), and give the stated ambiguity in the reconstruction.
7.4. The generating series of distinguished 0-generated walls. In light of Theorem6.14, in order to complete the proof of our main Theorem 1.7 for type $D$, we need the following combinatorial result, the precise analogue of Proposition 2.7 in type $A$.
Theorem 7.7. Let $\Delta$ be of type $D_{n}$, and let $\xi$ be a primitive $(2 n-1)$-st root of unity. Then the generating series of the set $\mathcal{Z}_{\Delta}^{0}$ of distinguished 0-generated Young walls is given in terms of the generating function of all Young walls by the following substitution:

$$
\sum_{Y \in \mathcal{Z}_{\Delta}^{0}} q^{\mathrm{wt}_{0}(Y)}=\left.Z_{\Delta}\left(q_{0}, \ldots, q_{n}\right)\right|_{q_{0}=\xi^{2} q, q_{1}=\cdots=q_{n}=\xi}
$$

Recall that by Lemma 6.1 (1) for each $Y \in \mathcal{Z}_{\Delta}$ there is a unique Young wall $Y_{0} \in \mathcal{Z}_{\Delta}^{0}$ which has exactly the same set of label 0 blocks as $Y$. This induces a combinatorial map

$$
p: \mathcal{Z}_{\Delta} \rightarrow \mathcal{Z}_{\Delta}^{0}
$$

which, in analogy once again with the type $A$ proof, will be the key construction in our proof of Theorem 7.7

Proposition 7.8. On the abacus representation of Young walls, the map $p: \mathcal{Z}_{\Delta} \rightarrow \mathcal{Z}_{\Delta}^{0}$ corresponds to the following steps:
(1) If a row ends with a white (resp., black) bead on $R_{2 n-2}$ corresponding to a column of height $\lambda_{k}$, and $k+n_{k}$ is even (resp. odd), then one bead should be moved to the next position, which is the leftmost in the next row. This is applied also on the zeroth position of the abacus, where we assume that there are infinitely many beads. This corresponds to the appearance of a new column in the Young wall represented by the abacus.
(2) Every bead on the rulers $R_{1}, \ldots, R_{2 n-4}$ is moved to right as much as possible according to the abacus rules.
(3) If there is at least one bead on the rulers $R_{1}, \ldots, R_{2 n-3}$ after performing Step 1 on the previous row, and the number of beads on $R_{2 n-2}$ is even, then one more bead is moved onto $R_{2 n-2}$. If there were beads on $R_{2 n-2}$ already, then this step does not change the parity of $k+n_{k}$ for the rightmost bead. If there were no beads on $R_{2 n-2}$ before, then this beads should take the appropriate color and it is possible to see that it can not be moved further with Step 1.
(4) Let $s$ be the total number of beads on the rulers $R_{1}, \ldots, R_{n-1}$ after performing the Steps 1-3. If $s>n-2$, then move all these beads on $R_{n-1}$. In this case some of these beads were here previously, so the color of the whole group of beads is already determined. If $s \leq n-2$, then move them onto the rulers $R_{1}, \ldots, R_{n-2}$, each as right as possible.
Proof. Step 1 enforces condition (D2). It also enforces condition (D3) when applied to the minus first row. Step 2 enforces conditions (D4) and (D5), Step 3 enforces condition (D1), and finally Step 4 enforces condition (D6).

The fibers of the map $p$ can be described as follows. Given a Young wall $Y \in \mathcal{Z}_{\Delta}^{0}$, we are allowed to move beads to the left and, occasionally, to the right, using the following rules.
(1) From the last position of the $i-t h$ row only one (resp. zero) bead can be moved to the left if $l_{i}$ is odd (resp., even).
(2) Every other bead is allowed to moved to the left in its row if the result is a valid abacus configuration.
(3) The leftmost bead in a row can be moved to the last position of the previous row. There it will take the color white if $\sum_{j<i} t_{j}+\sum_{j>i, t_{j} \neq 0 \bmod n-1} t_{j} \equiv 1 \bmod 2$, and grey otherwise.
(4) If $t_{i}-l_{i} \leq 2 n-2$, then the beads to the left from the $n-1$-st position are allowed to be moved to the right at most onto the ruler $R_{n-1}$.
(5) If $t_{i}-l_{i} \leq 2 n-2$, then any configuration, in which there is at least one bead at the ( $n-1$ )-st position, has to be counted with multiplicity two.
Let us call the beads that can be moved according to these rules movable. For a row with data $(t, l)$, let us also introduce the number $c(t, l)$, which is signed sum of the distance of the beads from the $R_{n-1}$-st ruler, where the sum runs over the movable beads, and the beads to the left of the $R_{n-1}$-st ruler are counted with negative sign, and the beads to the right of it are counted with positive sign. These numbers are listed in the table below:

|  | $l \equiv 0 \bmod 2$ | $l \equiv 1 \bmod 2$ |
| ---: | :---: | :---: |
| $0 \leq t-l \leq n-2$ | $\binom{n-1}{2}-\binom{n-1-t+l}{2}$ | $\binom{n}{2}-\binom{n-1-t+l}{2}$ |
| $n-1 \leq t-l \leq 2 n-3$ | $\binom{n-1}{2}-\binom{t-l-n+1}{2}$ | $\binom{n}{2}-\binom{t-l-n+1}{2}$ |
| $2 n-2 \leq t-l$ | $\binom{n-1}{2}$ | $\binom{n}{2}$ |

Lemma 7.9. The contribution of a row with data $(t, l)$ to the total weight of the fiber is $\xi^{-c(t, l)}$.
Proof. If $l$ is even but nonzero, then according to Corollary $7.6 l=t$ and there isn't any movable bead.

If $l$ is odd, then there is one movable bead on the $R_{2 n-2}$-nd ruler. Assume first that $t \leq n-1$. Then the expression

$$
\sum_{n_{1}=0}^{2 n-t+l-2} \sum_{n_{2}=0}^{n_{1}} \cdots \sum_{n_{t-l+1}=0}^{n_{t-l}}\left(\xi^{-1}\right)^{n_{1}+\cdots+n_{t-l+1}}=\binom{2 n-1}{t-l+1}_{\xi^{-1}}
$$

counts once every preimage, in which there is at most one bead at the $n-1$-st position. Similarly

$$
\left(\xi^{-1}\right)^{n-t+l-1} \sum_{n_{2}=0}^{2 n-t} \sum_{n_{3}=0}^{n_{2}} \ldots \sum_{n_{t-l+1}=0}^{n_{t-l}}\left(\xi^{-1}\right)^{n_{2}+\cdots+n_{t-l+1}}=\left(\xi^{-1}\right)^{n-t+l-1}\binom{2 n-1}{t-l}_{\xi^{-1}}
$$

counts once every preimage, in which there is at least one and at most two beads at the $n-1$-st position, since we moved one bead from the leftmost occupied position to the $n-1$-st position, and we fixed it there. The next term is given by

$$
\begin{gathered}
\left(\xi^{-1}\right)^{(n-t+l-1)+(n-t+l)} \sum_{n_{3}=0}^{2 n-t} \sum_{n_{4}=0}^{n_{3}} \cdots \sum_{n_{t-l+1}=0}^{n_{t-l}}\left(\xi^{-1}\right)^{n_{3}+\cdots+n_{t-l+1}} \\
=\left(\xi^{-1}\right)^{(n-t+l-1)+(n-t+l)}\binom{2 n-1}{t-l-1}_{\xi^{-1}}
\end{gathered}
$$

which counts once every preimage, in which there is at least two and at most three beads at the $n-1$-st ruler. Continuing in this fashion and summing up in the end we get

$$
\binom{2 n+1}{t-l+1}_{\xi^{-1}}+\sum_{i=1}^{t-l+1} \xi^{-\sum_{j=n-t+l-1}^{n-t+l-2+i} j}\binom{2 n-1}{t-l+1-i}_{\xi^{-1}}
$$

It can be checked that $\binom{2 n-1}{k}_{\xi^{-1}}=0$ unless $k=0$, in which case it is equal to 1 . Therefore, the only surviving part is the one with $i=t-l+1$ :

$$
\xi^{-\sum_{j=n-t+l-1}^{n-1} j}=\xi^{\binom{n}{2}-\binom{n-1-t+l}{2}}=\xi^{-c(t, l)}
$$

The proofs of the remaining two cases, when $l$ is odd, are very similar. The only difference in the case $2 n-2 \leq t-l$ is that first the preimages with zero or one extra movable beads at $R_{n-1}$ have to be counted, then the preimages with two or three extra movable beads, etc. As in the case $t \leq n-1$, the only nonzero term is the one where in the beginning all the movable beads have been shifted to the $R_{n-1}$-st ruler, and in this case the powers of $\xi^{-1}$ sum up to $\binom{n}{2}=c(t, l)$.
Corollary 7.10. Let $Y \in \mathcal{Z}_{\Delta}^{0}$ be a distinguished 0-generated Young wall described by the data $\left\{\left(t_{i}, l_{i}\right)_{i}\right\}$. Then

$$
\begin{aligned}
\left.\sum_{Y^{\prime} \in p_{*}^{-1}(Y)} \underline{q} \underline{\underline{\mathrm{wt}}\left(Y^{\prime}\right)}\right|_{q_{1}=\cdots=q_{n}=\xi, q_{0}=\xi^{2} q} & =\left.\underline{q}^{\underline{\mathrm{wt}}(Y)}\right|_{q_{1}=\cdots=q_{n}=\xi, q_{0}=\xi^{-(2 n-3)} q} \cdot \xi^{-\sum_{i} c\left(t_{i}, l_{i}\right)} \\
& =q^{\mathrm{wt}_{0}(Y)} \xi^{\sum_{j \neq 0}\left(\mathrm{wt}_{j}(Y)-\operatorname{dim} \rho_{j} \cdot \mathrm{wt}_{0}(Y)\right)-\sum_{i} c\left(t_{i}, l_{i}\right)}
\end{aligned}
$$

Lemma 7.11. The core of a 0-generated Young wall is a 0-generated Young wall.
Proof. With each reduction step (B1)-(B4) we always remove a bar. In the original Young wall, the salient blocks were only label 0 half blocks and salient block-pairs of label $n-1 / n$.

A similar reasoning as in the type $A$ case shows that no salient block of label $2, \ldots, n-2$ can appear after we perform the step (B1) until possible. The same is true with (B2), since if we can perform it on a pair of beads in a row, then we can always perform it on the beads between them. More precisely, it can be seen that even label 1 salient blocks cannot appear during these two steps because the parity conditions in (D1) and (D2) is always maintained by the reduction steps.

The parity conditions in (D1) and (D2) are maintained by the step (B3) as well. If we perform (B4) on a pair of beads, then we can perform it on this pair as long as they disappear from the abacus. Hence, when performing (B4) until possible we also get back the good parities. After the reduction is completed there cannot be any bead on the rulers $R_{1}, \ldots, R_{n-2}$, and all the beads on the rulers $R_{n}, \ldots, R_{2 n-3}$ are right-adjusted. Therefore the conditions (D4) and (D5) are also satisfied.

If the total number of beads was initially even, then since no label 1 salient block can appear, the total number of beads in the end must be even as well. So the final Young wall will satisfy (D3). If in the total number of beads in the initial abacus configuration is odd, then in the first row $t_{1}-l_{1}=2 n-4$. Hence, the beads on ruler $R_{2 n-2}$ in the first row are necessarily white and one of them can be taken away from the abacus with the step (B3), together with the beads on the other rulers using (B2). This is an odd number of beads removed from the abacus. After this the total number of beads is even, so we continue as in the even case.

Proof of Theorem 7.7. In light of Corollary 7.10 it remains to show that

$$
\xi^{\sum_{j \neq 0}\left(\mathrm{wt}_{j}(Y)-\operatorname{dim} \rho_{j} \cdot \mathrm{wt}_{0}(Y)\right)-\sum_{i} c\left(t_{i}, l_{i}\right)}=1 .
$$

We follow in the line of the proof the $A_{n}$ case.
Step 1: Reduction to 0 -generated cores. According to Lemma 7.11 the core of a 0 -generated Young wall is a 0 -generated core. It is immediate from the definition of $c(t, l)$ that the steps (B1), (B2) and (B4) leave the sum $\sum_{i} c\left(t_{i}, l_{i}\right)$ unchanged, while (B3) is a bit more complicated. Indeed,

- (B1) removes one movable bead from a row, and adds one to another on the same ruler;
- (B2) removes two movable beads, but these two contribute with opposite signs into $c(t, l)$;
- (B4) either moves non-movable beads from $R_{2 n-2}$ onto $R_{n-1}$, or beads from $R_{n-1}$ into non-movable beads on $R_{2 n-2}$.
Moving beads on $R_{n-1}$ according to (B3) does not affect the numbers $c(t, l)$. If (B3) moves a bead on the ruler $R_{2 n-2}$ between rows having $l$ of different parity, then the sum of the movable beads at the $(2 n-2)$-nd positions of the two rows is constant, so $\sum_{i} c\left(t_{i}, l_{i}\right)$ does not change after such a step. If (B3) were able to move a bead on $R_{2 n-2}$ between rows both having odd $l$ 's, then this can happen only if they have the same color and if there is no bead on $R_{n-1}$ between them. This means that either there must be an even number of beads between them and they should have same kind of top block, or odd number of beads and different kind of top blocks. Both cases are forbidden in 0 -generated Young walls due to Lemma 7.5] For the same reasons (B3) cannot move beads on $R_{2 n-2}$ between rows both having even l's.

It follows from these considerations that for all $Y \in \mathcal{Z}_{\Delta}^{0}$

$$
\xi^{\sum_{j \neq 0}\left(\mathrm{wt}_{j}(Y)-\operatorname{dim} \rho_{j} \cdot \mathrm{wt}_{0}(Y)\right)-\sum_{i} c\left(t_{i}, l_{i}\right)}=\xi^{\sum_{j \neq 0}\left(\operatorname{wt}_{j}(\operatorname{core}(Y))-\operatorname{dim} \rho_{j} \cdot \mathrm{wt}_{0}(\operatorname{core}(Y))-\sum_{i} c\left(t_{i}, l_{i}\right)\right.} .
$$

## Step 2: Reduction to distinguished 0-generated cores.

As described above, for any 0-generated core $Y$ there is a decomposition as $\lambda(Y)=\mu(Y) \cup \nu(Y)$, where $\mu(Y) \in \mathcal{C}^{1}, \nu(Y) \in \mathcal{C}^{2}$, and we have $\nu(Y)=\nu^{(0)}(Y)+\nu^{(1)}(Y)$, where $\nu^{(0)}(Y)$ and $\nu^{(1)}(Y)$ are two-cores, and parts in $\nu^{(1)}(Y)$ have colors given by their parity. Since $Y$ is 0 -generated, the largest part of $\nu(Y)$ has to be even, otherwise there is no bead at the rightmost position in the last row of the abacus. Therefore, the abacus represenation of $Y$ can be described as follows.
(1) There is no bead on the rulers $R_{k}$ for $1 \leq k \leq n-2$;
(2) On the ruler $R_{n-1}$ all the beads are of the same color, the beads are at the first, let's say, $m$ positions, exactly one bead at each.
(3) The positions in the first $m$ rows of the rulers $R_{k}$ for $n \leq k \leq 2 n-3$ are all filled up with beads, the other beads are pushed to the right as much as possible.
(4) There is at least $m$ beads on the ruler $R_{2 n-2}$, one at each of the first positions and there is no space between them. The first $m$ of these are all white. The total number of the rest is even, and half of them is white, half of them is black.
The abacus of a typical 0-generated core looks like this:

| $R_{1}$ | $\cdots$ | $R_{n-2}$ | $R_{n-1}$ | $R_{n}$ | $\cdots$ | $R_{2 n-3}$ | $R_{2 n-2}$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  | $\bigcirc$ | $\bigcirc$ | $\cdots$ | $\bigcirc$ | $\bigcirc$ |  |  |
|  | $\vdots$ | $\vdots$ |  | $\vdots$ | $\vdots$ |  |  |
|  | $\bigcirc$ | $\bigcirc$ | $\cdots$ | $\bigcirc$ | $\bigcirc$ |  |  |
|  |  | $\bigcirc$ | $\bigcirc$ | $\bigcirc$ | $\bigcirc$ |  |  |
|  |  |  | $\bigcirc$ | $\bigcirc$ | $\bigcirc$ |  |  |
|  |  |  |  | $\bigcirc$ | $\bigcirc$ |  |  |
|  |  |  |  | $\bigcirc$ | $\bigcirc$ |  |  |
|  |  |  |  | $\bigcirc$ | $\bigcirc$ |  |  |

It is not true that the core of a distinguished 0 -generated Young wall is always distinguished. But we can reduce each non-distinguished core further using the reduction map red and then taking the core of the result using the steps (B1)-(B4) again. The result of this is a distinguished

0 -generated core. The first step corresponds to shifting the beads on $R_{n-1}$ one step left. This decreases $\sum_{i} c\left(t_{i}, l_{i}\right)$ by $m$, and decreases $\sum_{j \neq 0}\left(\mathrm{wt}_{j}(Y)-\operatorname{dim} \rho_{j} \cdot \mathrm{wt}_{0}(Y)\right)$ also by $m$. The second step does not change these numbers by the considerations of Step 1. So we are done once we know the statement for distinguished 0-generated cores.

We remark here that during the first step the color of every second bead on the ruler $R_{2 n-2}$ changes. In the end the color of the beads on this ruler will alternate.

Step 3: The case of distinguished 0 -generated cores.
Let $\left(z_{1}, \ldots, z_{n}\right)$ be the integer tuple assigned to $Y$ as in 7.2. Then $z_{n-1}=z_{n}$ since there is no bead on $R_{n-1}$ and the color of the beads on $R_{2 n-2}$ alternates. In particular, $z_{n-1}+z_{n}$ is an even number.

We claim that

$$
\sum_{1=1}^{n} m_{n}=\sum_{i} c\left(t_{i}, l_{i}\right)
$$

Indeed, by Theorem $7.4(3), \sum_{1=1}^{n} m_{n}=-\sum_{i=1}^{n-2}(n-1-i) z_{i}-(n-1) c\left(z_{n-1}+z_{n}\right)-(n-1) b$. The number of movable beads on $R_{2 n-2-i}$ is $-z_{i}$ if $1 \leq i \leq n-1$, and it is $-c\left(z_{n-1}+z_{n}\right)-b$ if $i=0$. This proves the claim.

By Theorem 7.4 (1), for a core

$$
\underline{q}^{\underline{\mathrm{wt}}(Y)}=q_{1}^{m_{1}} \cdots \cdot q_{n}^{m_{n}}\left(q^{1 / 2}\right)^{\bar{m}^{\top} \cdot C \cdot \bar{m}}
$$

As in the type $A$ case, on the right hand side of this expression $q_{0}$ only appears in the $q^{1 / 2}$-term. Hence

$$
\frac{1}{2} \bar{m}^{\top} \cdot C \cdot \bar{m}=\mathrm{wt}_{0}(Y)
$$

and

$$
\left.\left(q^{1 / 2}\right)^{\bar{m}^{\top} \cdot C \cdot \bar{m}}\right|_{q_{1}=\cdots=q_{n}=\xi, q_{0}=\xi^{2} q}=q^{\mathrm{wt}_{0}(Y)}
$$

On the other hand,

$$
\left.\underline{q}^{\underline{\mathrm{wt}}(Y)}\right|_{q_{1}=\cdots=q_{n}=\xi, q_{0}=\xi^{2} q}=q^{\mathrm{wt}_{0}(Y)} \xi^{\sum_{j \neq 0}\left(\mathrm{wt}_{j}(Y)-\operatorname{dim} \rho_{j} \cdot \mathrm{wt}_{0}(Y)\right)} .
$$

Therefore
and so indeed

$$
\xi^{\sum_{j \neq 0}\left(\mathrm{wt}_{j}(Y)-\operatorname{dim} \rho_{j} \cdot \mathrm{wt}_{0}(Y)\right)}=\xi^{\sum_{i=1}^{n} m_{n}}=\xi^{\sum_{i} c\left(t_{i}, l_{i}\right)}
$$

$$
\xi^{\sum_{j \neq 0}\left(\mathrm{wt}_{j}(Y)-\operatorname{dim} \rho_{j} \cdot \mathrm{wt}_{0}(Y)\right)-\sum_{i} c\left(t_{i}, l_{i}\right)}=1 .
$$

## Appendix A. Background on representation theory

This section plays no logical role in our paper, but it provides important background. For further discussion about the role of representation theory, see the announcement [21].
A.1. Affine Lie algebras and extended basic representations. Let $\Delta$ be an irreducible finitedimensional root system, corresponding to a complex finite dimensional simple Lie algebra $\mathfrak{g}$ of rank $n$. Attached to $\Delta$ is also an (untwisted) affine Lie algebra $\tilde{\mathfrak{g}}$, but a slight variant will be more interesting for us, see e.g. [11, Sect 6]. Denote by $\widetilde{\mathfrak{g} \oplus \mathbb{C}}$ the Lie algebra that is the direct sum of the affine Lie algebra $\mathfrak{g}$ and an infinite Heisenberg algebra $\mathfrak{h e i s , ~ w i t h ~ t h e i r ~ c e n t e r s ~ i d e n t i f i e d . ~}$

Let $V_{0}$ be the basic representation of $\tilde{\mathfrak{g}}$, the level-1 representation with highest weight $\omega_{0}$. Let $\mathcal{F}$ be the standard Fock space representation of $\mathfrak{h e i s}$, having central charge 1 . Then $V=V_{0} \otimes \mathcal{F}$ is a representation of $\widetilde{\mathfrak{g} \oplus \mathbb{C}}$ that we may call the extended basic representation. By the Frenkel-Kac theorem 12], in fact

$$
V \cong \mathcal{F}^{n+1} \otimes \mathbb{C}\left[Q_{\Delta}\right]
$$

where $Q_{\Delta}$ is the root lattice corresponding to the root system $\Delta$. Here, for $\beta \in \mathbb{C}\left[Q_{\Delta}\right], \mathcal{F}^{n+1} \otimes e^{\beta}$ is the sum of weight subspaces of weight $\omega_{0}-\left(m+\frac{\langle\beta, \beta\rangle}{2}\right) \delta+\beta, m \geq 0$. Thus, we can write the character of this representation as

$$
\begin{equation*}
\operatorname{char}_{V}\left(q_{0}, \ldots, q_{n}\right)=e^{\omega_{0}}\left(\prod_{m>0}\left(1-q^{m}\right)^{-1}\right)^{n+1} \cdot \sum_{\beta \in Q_{\Delta}} q_{1}^{\beta_{1}} \cdots \cdots q_{n}^{\beta_{n}}\left(q^{1 / 2}\right)^{\langle\beta, \beta\rangle} \tag{20}
\end{equation*}
$$

where $q=e^{-\delta}$, and $\beta=\left(\beta_{1}, \ldots, \beta_{n}\right)$ is the expression of an element of the finite type root lattice in terms of the simple roots.
Example A.1. For $\Delta$ of type $A_{n}$, we have $\mathfrak{g}=\mathfrak{s l}_{n+1}, \tilde{\mathfrak{g}}=\tilde{\mathfrak{s l}}_{n+1}, \widetilde{\mathfrak{g} \oplus \mathbb{C}}=\tilde{\mathfrak{g}}_{n+1}$. In this case there is in fact a natural vector space isomorphism $V \cong \mathcal{F}$ with Fock space itself, see e.g. 41, Section 3E].
A.2. Affine crystals. The basic representations $V_{0}, V$ of $\tilde{\mathfrak{g}}, \widetilde{\mathfrak{g} \oplus \mathbb{C}}$ respectively can be constructed on vector spaces spanned by explicit "crystal" bases. Crystal bases have many combinatorial models; in types $A$ or $D$, the sets denoted with $\mathcal{Z}_{\Delta}$ in the main part of our paper provide one possible combinatorial model [25] for the crystal basis for the basic representation. More precisely, given $\Delta$ of type $A$ or $D$, there is a combinatorial condition which singles out a subset $\mathcal{Y}_{\Delta} \subset \mathcal{Z}_{\Delta}$. The basic representation $V_{0}$ of $\tilde{\mathfrak{g}}$ has a basis [34, 24] in bijection with elements of $\mathcal{Y}_{\Delta}$. The extended basic representation $V$ of $\widetilde{\mathfrak{g} \oplus \mathbb{C}}$ has a basis in bijection with elements of $\mathcal{Z}_{\Delta}$. The canonical embedding $V_{0} \subset V$, defined by the vacuum vector inside Fock space $\mathcal{F}$, is induced by the inclusion $\mathcal{Y}_{\Delta} \subset \mathcal{Z}_{\Delta}$.
A.3. Affine Lie algebras and Hilbert schemes. As before, let $\Gamma<\operatorname{SL}(2, \mathbb{C})$ be a finite subgroup and let $\Delta \subset \widetilde{\Delta}$ be the corresponding finite and affine Dynkin diagrams. It is a well-known fact that the equivariant Hilbert schemes $\operatorname{Hilb}^{\rho}\left(\left[\mathbb{C}^{2} / \Gamma\right]\right)$ for all finite dimensional representations $\rho$ of $G$ are Nakajima quiver varieties 37] associated to $\Delta$, with dimension vector determined by $\rho$, and a specific stability condition (see [13, 35] for more details for type $A$ ).

Nakajima's general results on the relation between the cohomology of quiver varities and KacMoody algebras, specialized to this case, imply 37] that the direct sum of all cohomology groups $H^{*}\left(\operatorname{Hilb}^{\rho}\left(\left[\mathbb{C}^{2} / G\right]\right)\right)$ is graded isomorphic to the extended basic representation $V$ of the corresponding extended affine Lie algebra $\widetilde{\mathfrak{g} \oplus \mathbb{C}}$ defined in A. 1 above. By [36, Section 7], these quiver varieties have no odd cohomology. Thus the character formula (20) implies Theorem 1.3 in all types $A, D$, and $E$.

## Appendix B. Joins

Recall [1] that the join $J(X, Y) \subset \mathbb{P}^{n}$ of two projective varieties $X, Y \subset \mathbb{P}^{n}$ is the locus of points on all lines joining a point of $X$ to a point of $Y$ in the ambient projective space. One well-known example of this construction is the following. Let $L_{1} \cong \mathbb{P}^{k}$ and $L_{2} \cong \mathbb{P}^{n-k-1}$ be two disjoint projective linear subspaces of $\mathbb{P}^{n}$.

Lemma B.1. The join $J\left(L_{1}, L_{2}\right) \subset \mathbb{P}^{n}$ equals $\mathbb{P}^{n}$. Moreover, the locus $\mathbb{P}^{n} \backslash\left(L_{1} \cap L_{2}\right)$ is covered by lines uniquely: for every $p \in \mathbb{P}^{n} \backslash\left(L_{1} \cap L_{2}\right)$, there exists a unique line $\mathbb{P}^{1} \cong p_{1} p_{2} \subset \mathbb{P}^{n}$ with $p_{i} \in L_{i}$, containing $p$.

Let now $H \subset \mathbb{P}^{n}$ be a hyperplane not containing the $L_{i}$, which we think of as the hyperplane "at infinity". Let $V=\mathbb{P}^{n} \backslash H \cong \mathbb{A}^{n}$. Let $\bar{L}_{i}=L_{i} \cap H$, and let $L_{i}^{o}=L_{i} \backslash \bar{L}_{i}=L_{i} \cap V$ be the affine linear subspaces in $V$ corresponding to $L_{i}$. Finally let $X=J\left(\bar{L}_{1}, L_{2}\right) \cong \mathbb{P}^{n-1}$ and $X^{o}=X \backslash(X \cap H) \cong \mathbb{A}^{n-1}$.

Lemma B.2. Projection away from $L_{1}$ defines a morphism $\phi: X \rightarrow L_{2}$, which is an affine fibration with fibres isomorphic to $\mathbb{A}^{k}$. $\phi$ restricts to a morphism $\phi^{o}: X^{o} \rightarrow L_{2}^{o}$, which is a trivial affine fibration over $L_{2}^{o} \cong \mathbb{A}^{n-k-1}$.

In geometric terms, the map $\phi$ is defined on $X \backslash L_{2}$ as follows: take $p \in X \backslash L_{2}$, find the unique line $p_{1} p_{2}$ passing through it, with $p_{1} \in \bar{L}_{1}$ and $p_{2} \in L_{2}$; then $\phi(p)=p_{2}$.

Let now $U$ be a projective subspace of $H$ which avoids $\bar{L}_{2}$. Let $U_{1} \subset U$ be a codimension one linear subspace, and $W=U \backslash U_{1}$ its affine complement. In the main text, we need the following statement.

Lemma B.3. $\chi\left(\left(J\left(L_{2}^{o}, W\right) \backslash L_{2}^{o}\right) \cap V\right)=0$.
Proof. With the same argument as in Lemma B.2 $J\left(L_{2}^{o}, W\right) \cap V$ is a fibration over $L_{2}^{o}$ with fiber Cone $(W)$, and $\left(J\left(L_{2}^{o}, W\right) \backslash L_{2}^{o}\right) \cap V$ is a fibration over $L_{2}^{o}$ with fiber Cone $(W) \backslash\{$ vertex $\}$. Since Cone $(W) \backslash\{$ vertex $\}=\mathbb{C}^{*} \times W$, the projection from $\left(J\left(L_{2}^{o}, W\right) \backslash L_{2}^{o}\right) \cap V \cong L_{2}^{o} \times W \times \mathbb{C}^{*}$ to $L_{2}^{o} \times W$ has fibers $\mathbb{C}^{*}$. The lemma follows.

We finally recall the base-change property of joins.
Lemma B.4. [1, B1.2] Let $S$ be an arbitrary scheme. Then for schemes $X, Y \subset \mathbb{P}_{S}^{n}$ and an $S$-scheme $T$, we have the following equality in $\mathbb{P}_{T}^{n}$ :

$$
J\left(X \times_{S} T, Y \times{ }_{S} T\right)=J(X, Y) \times{ }_{S} T
$$

## References

1. A. B. Altman and S. L. Kleiman, Joins of schemes, linear projections, Compos. Math. 31 (1975), no. 3, 309-343.
2. G. E. Andrews, Generalized Forbenius Partitions, Mem. Amer. Math. Soc., vol. 301, Amer. Math. Soc., Providence, 1984.
3. A. Białynicki-Birula, Some theorems on actions of algebraic groups, Ann. of Math. (1973), 480-497.
4. M. Brion, Invariant Hilbert schemes, Handbook of moduli. Vol. I, Adv. Lect. Math., vol. 24, Int. Press, Somerville, MA, 2013, pp. 64-117.
5. H. Cartan, Quotient d'un espace analytique par un groupe d'automorphismes, Algebraic geometry and topology: a symposium in honor of S. Lefschetz, Princeton University Press, Princeton, NJ, 1957, pp. 90-102.
6. J. Cheah, On the cohomology of Hilbert schemes of points, J. Algebraic Geom. 5 (1996), no. 3, 479-512.
7. A. N. de Celis, Dihedral groups and G-Hilbert schemes, Ph.D. thesis, University of Warwick, 2008.
8. R. Dijkgraaf and P. Sułkowski, Instantons on ALE spaces and orbifold partitions, J. High Energy Phys. 3 (2008), 013-013.
9. David Eisenbud and Joe Harris, 3264 and all that: A second course in algebraic geometry, Cambridge University Press, 2016.
10. G. Ellingsrud and S. A. Strømme, On the homology of the Hilbert scheme of points in the plane, Invent. Math. 87 (1987), no. 2, 343-352.
11. P. Etingof, Symplectic reflection algebras and affine Lie algebras, Mosc. Math. J 12 (2012), no. 3, 543-565.
12. I. B. Frenkel and V. G. Kac, Basic representations of affine Lie algebras and dual resonance models, Invent. Math. 62 (1980), no. 1, 23-66.
13. Sh. Fujii and S. Minabe, A Combinatorial Study on Quiver Varieties, SIGMA. Symmetry, Integrability and Geometry: Methods and Applications 13 (2017), 052.
14. F. Garvan, D. Kim, and D. Stanton, Cranks and t-cores, Invent. Math. 101 (1990), no. 1, 1-17.
15. I. G. Gordon, Quiver varieties, category $\mathcal{O}$ for rational Cherednik algebras, and Hecke algebras, Int. Math. Res. Pap. (2008), no. 3, rpn006.
16. L. Göttsche, The Betti numbers of the Hilbert scheme of points on a smooth projective surface, Math. Ann. 286 (1990), no. 1, 193-207.
17. I. Grojnowski, Instantons and affine algebras I: The Hilbert scheme and vertex operators, Math. Res. Lett. 3 (1996), no. 2, 275-292.
18. S. M. Gusein-Zade, I. Luengo, and A. Melle-Hernández, On generating series of classes of equivariant Hilbert schemes of fat points, Mosc. Math. J. 10 (2010), no. 3, 593-602.
19. Á. Gyenge, Hilbert scheme of points on cyclic quotient singularities of type ( $p, 1$ ), Period. Math. Hungar. 73 (2016), no. 1, 93-99.
20. , Hilbert schemes of points on some classes of surface singularities, Ph.D. thesis, Eötvös Lóránd University, Budapest, 2016.
21. Á. Gyenge, A. Némethi, and B. Szendrői, Euler characteristics of Hilbert schemes of points on surfaces with simple singularities, Int. Math. Res. Not. 2017 (2016), no. 13, 4152-4159.
22. R. Hartshorne, Algebraic geometry, Graduate Text in Mathematics, vol. 52, Springer, 1977.
23. Y. Ito and I. Nakamura, Hilbert schemes and simple singularities, London Mathematical Society Lecture Note Series (1999), 151-234.
24. S.-J. Kang, Crystal bases for quantum affine algebras and combinatorics of Young walls, Proc. London Math. Soc. 86 (2003), no. 1, 29-69.
25. S.-J. Kang and J.-H. Kwon, Crystal bases of the Fock space representations and string functions, J. Algebra 280 (2004), no. 1, 313-349.
26. M. Kapranov and E. Vasserot, Kleinian singularities, derived categories and Hall algebras, Math. Ann. 316 (2000), no. 3, 565-576.
27. J.-H. Kwon, Affine crystal graphs and two-colored partitions, Lett. Math. Phys. 75 (2006), no. 2, 171-186.
28. S. Lang, Algebra, Revised third edition, Graduate Texts in Mathematics, vol. 211, Springer, 2002.
29. B. Leclerc and H. Miyachi, Some closed formulas for canonical bases of Fock spaces, Represent. Theory 6 (2002), no. 2002, 290-312.
30. O. Lorscheid and Th. Weist, Quiver Grassmannians of extended Dynkin type D - Part 1: Schubert systems and decompositions into affine spaces, http://arxiv.org/abs/1507.00392, 2015.
31. $\qquad$ , Quiver Grassmannians of extended Dynkin type D - Part 2: Schubert decompositions and F-polynomials, http://arxiv.org/abs/1507.00395, 2015.
32. I. G. Macdonald, The Poincaré polynomial of a symmetric product, Math. Proc. Cambridge Philos. Soc. 58 (1962), no. 4, 563-568.
33. D. Maulik, Stable pairs and the HOMFLY polynomial, Invent. Math. 204 (2016), no. 3, 787831.
34. K. C. Misra and T. Miwa, Crystal base for the basic representation of $U_{q}(\widehat{\mathfrak{s l}}(n))$, Comm. Math. Phys. 134 (1990), no. 1, 79-88.
35. K. Nagao, Quiver varieties and Frenkel-Kac construction, J. Algebra 321 (2009), no. 12, 37643789.
36. H. Nakajima, Quiver varieties and finite dimensional representations of quantum affine algebras, J. Amer. Math. Soc. 14 (2001), no. 1, 145-238.
37. _ Geometric construction of representations of affine algebras, Proceedings of the International Congress of Mathematicians (Beijing, 2002), vol. 1, IMU, Higher Ed, 2002, pp. 423438.
38. H. Nakajima and K. Yoshioka, Instanton counting on blowup. I. 4-dimensional pure gauge theory, Invent. Math. 162 (2005), no. 2, 313-355.
39. A. Oblomkov and V. Shende, The Hilbert scheme of a plane curve singularity and the HOMFLY polynomial of its link, Duke Math. J. 161 (2012), no. 7, 1277-1303.
40. A. Savage and P. Tingley, Quiver grassmannians, quiver varieties and the preprojective algebra, Pacific J. Math. 251 (2011), no. 2, 393-429.
41. P. Tingley, Notes on Fock Space, http://webpages.math.luc.edu/~ptingley/, 2011.
42. Y. Toda, $S$-Duality for surfaces with $A_{n}$-type singularities, Math. Ann. 363 (2015), 679-699.
43. B. Young and J. Bryan, Generating functions for colored 3D Young diagrams and the Donaldson-Thomas invariants of orbifolds, Duke Math. J. 152 (2010), no. 1, 115-153.
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[^0]:    1991 Mathematics Subject Classification. Primary 14N10; Secondary 05E10.
    Key words and phrases. Hilbert scheme, singularities, Euler characteristic, generating series, Young wall.

[^1]:    ${ }^{1}$ The combinatorics introduced in this section should really be called type $\tilde{D}_{n}^{(1)}$, but we do not wish to overburden the notation. Also we have reflected the pattern in a vertical axis compared to the pictures of [25, 27].
    ${ }^{2}$ In [25, 27], these arrangements are called proper Young walls. Since we will not meet any other Young wall, we will drop the adjective proper for brevity.

[^2]:    ${ }^{3}$ This is the properness condition of [25].

[^3]:    ${ }^{4}$ Again, for a missing half block only the half blocks of the same orientation have to be missing.

[^4]:    ${ }^{5}$ Once again, we should call it type $\tilde{D}_{n}^{(1)}$, but we simplify for ease of notation.

