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# Hypersurfaces in $P^{n}$ WITH 1-PARAMETER SYMMETRY GROUPS II 

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## Introduction

We are interested in hypersurfaces $V \subset P^{n}(\mathbb{C})$ defined by homogeneous equations $f\left(x_{0}, \ldots, x_{n}\right)=0$ of degree $d$. We say that $V$ is quasi-smooth if $V$ has isolated singularities and is not a cone. If $V$ admits a subgroup $G$ of $\mathrm{PGL}_{n+1}$ of symmetries with $r=\operatorname{dim} G \geq 1$, we call $V r$-symmetric.

In [4] we gave a detailed discussion of 1-symmetric quasi-smooth hypersurfaces in the case when $G$ is semi-simple. The main object of this paper is to give a corresponding analysis when $G$ is unipotent.

Our first main result Theorem 2.4 lists the possible cases. Let $G$ be a unipotent group of type given by the sequence $R=\left\{r_{1} \geq r_{2} \geq \ldots\right\}$ (i.e. the Jordan blocks have sizes $r_{i}+1$; we omit zeros in writing $R$ ). Then we have one of the following:

$$
\begin{array}{ll}
\text { Case 2: } & d \geq 3, R=\{2\}, \\
\text { Case 4: } \quad d=3, R=\{4\}, & \text { Case 3: 21: } d=3, R=\{3\}, \\
\text { Case }=\{2,1\} .
\end{array}
$$

We find that Case 21 splits into two: one a subcase of Case 2 ; the other we rename Case 5 .

Our second main conclusion is the calculation of the total Milnor number $\mu(V)$ (the sum of the Milnor numbers $\mu_{P_{i}}(V)$ at all singular points $P_{i}$ of $V$ ). The result is, where the $V_{i}$ are auxiliary varieties defined ad hoc in each case:

| Case | $\mu(V)$ |
| :---: | :---: |
| 2 | $\frac{1}{2}(d-2)(2 d-1)(d-1)^{m}+\mu\left(V_{2}\right)+\mu\left(V_{3}\right)$ |
| 3 | $22.3^{m}$ |
| 4 | $11.2^{m}+\mu\left(V_{3}\right)$ |
| 5 | $25.2^{m}+\mu\left(V_{3}\right)$ |

There is a 'main' singular point $P$. Provided in Case 2 that $V_{2}$ is non-singular and in Case 5 that $V_{3}$ is, the Milnor number $\mu_{P}(V)$ is the first term of the sum and the singularity of $V$ at $P$ is semi-quasi-homogeneous.

The first two sections are devoted to preparation and the proof of Theorem 2.4. We then pause for a brief review of some important background results, holding for all quasi-smooth 1 -symmetric hypersurfaces; in particular, we recall that $\tau(V) \leq(d-1)^{n-2}\left(d^{2}-3 d+3\right)$, and attains this value if and only if $f$ is annihilated
by vector fields $\xi$ of degree 1 and $\eta$ of degree $d-2$, not a multiple of $\xi$. We will call $V$ oversymmetric in this case. Moreover, $f$ is 2 -symmetric if and only if it is oversymmetric with $d=3$, and is never 3 -symmetric. We briefly recall the enumeration of oversymmetric hypersurfaces in the semi-simple case. We also give a number of auxiliary methods of calculation of Milnor numbers, so as not to interrupt the main discussion.

After a brief recall of the invariant theory of the nilpotent actions we discuss Cases 2-5 in successive sections; in each case we discuss the geometry of the action, show how to reduce $f$ to a convenient normal form, analyse the conditions on $f$ for $V$ to be quasi-smooth, find the singular points, and study the total Milnor number $\mu(V)$ and the nature of the singularities presented. We proceed to discussion of the Tjurina number $\tau(V)$, and show that $V$ is always oversymmetric in Cases 3 and 21, never in Case 4, while in Case 2 by Theorem 5.7 it occurs if and only if either (a) $V_{3}$ is a cone, or (b) after change of co-ordinates if necessary, $\partial \phi / \partial B$ and $\partial \phi / \partial X$ both vanish along $X=B=0$.

In a final section we recapitulate the complete list of the five 2-symmetric cases in more detail.

## 1 Unipotent actions on vector spaces and algebras

If $N$ is a nilpotent endomorphism of a finite dimensional vector space $K$, we can choose co-ordinates to put $N$ into Jordan canonical form, and count the sizes of the blocks. If the block sizes are $\lambda_{1}, \ldots, \lambda_{t}$, arranged in non-increasing order, then $n=\sum_{i} \lambda_{i}$. If we write $\nu_{k}:=\operatorname{rank} N^{k-1}-\operatorname{rank} N^{k}$, then $\nu$ is the partition conjugate to $\lambda$, so both partitions are independent of the choice of co-ordinates. Our usual notation will be to set $r_{i}:=\lambda_{i}-1$ and let $R$ be the sequence of $r_{i}$, with zeroes omitted.

We recall the representation theory of the Lie algebra $s l_{2}$. Denote the canonical basis vectors of $s l_{2}$ by

$$
e:=\left(\begin{array}{ll}
0 & 1 \\
0 & 0
\end{array}\right), \quad f:=\left(\begin{array}{ll}
0 & 0 \\
1 & 0
\end{array}\right), \quad h:=\left(\begin{array}{cc}
1 & 0 \\
0 & -1
\end{array}\right) .
$$

These satisfy $[e, f]=h,[h, e]=2 e,[h, f]=-2 f$. Every (finite dimensional) $s l_{2}-$ module $M$ is a direct sum of irreducible modules, and any irreducible module of rank $s+1(s=0,1,2, \ldots)$ is isomorphic to the module $K_{s}$ with basis $x_{-s}, x_{2-s}, \ldots, x_{s-2}, x_{s}$ and action given by

$$
e . x_{r}=\frac{s-r}{2} x_{r+2}, \quad \text { f. } x_{r}=\frac{s+r}{2} x_{r-2}, \quad \text { h. } x_{r}=r x_{r} .
$$

Thus the eigenvalues of $h$ on $M$ are all integers, and we can define a grading on $M$ by assigning weight $r$ to the eigenspace belonging to the eigenvalue $r$. Then for any $r \geq 0, e^{r}$ gives an isomorphism of $M_{-r}$ on $M_{r}$ and $f^{r}$ gives an isomorphism of $M_{r}$ on $M_{-r}$.

Lemma 1.1. For any $\operatorname{sl}_{2}-$ module $M$,
(i) if $r \geq 0, f^{r}$ gives an isomorphism from the weight space $M_{r}$ to $M_{-r}$,
(ii) if $r>0$, then $\operatorname{Ker}\left(f \mid M_{r}\right)=0$,
(iii) if $r<0$, then $\operatorname{dim} \operatorname{Ker}\left(f \mid M_{r}\right)=\operatorname{dim} M_{r}-\operatorname{dim} M_{r-2}$,
(iv) if $x \in M$ and $f . x=h . x=0$, then $e . x=0$.
(i)-(iii) follow by inspection from the remarks above. It suffices to prove (iv) for each $K_{s}$. But if $s>0$ then $f . x=h . x=0$ for $x \in K_{s}$ implies $x=0$, while if $s=0$ then $e . x=h . x=f . x=0$ for any $x \in K_{0}$.

Lemma 1.2. The action of a nilpotent endomorphism $N$ on a (finite dimensional) vector space $K$ can be extended to an sl $l_{2}$ action, with $N$ acting as $f$.

Proof. Express $K$ as the direct sum of monogenic modules: say $K_{s}^{\prime}$ has basis $x, N . x, N^{2} . x, \ldots, N^{s} . x$ with $N^{s+1} . x=0$. If we set, for $0 \leq i \leq s, N^{i} x=: \frac{s!}{(s-i)!} x_{s-2 i}$, then we have $N x_{s-2 i}=(s-i) x_{s-2 i-2}$ or, writing $r=s-2 i, N x_{r}=\frac{s+r}{2} x_{r-2}$. We can now set, for each $r$, e. $x_{r}=\frac{s-r}{2} x_{r+2}, h . x_{r}=r x_{r}$.

The action, and the grading it defines, are not determined solely by the nilpotent action. However, if we define the weight filtration by letting $F_{v} K$ be the sum of the eigenspaces of $h$ belonging to eigenvalues $\leq v$, we have

$$
F_{v} K=\sum_{p \in \mathbb{Z}} \operatorname{Ker} N^{p} \cap \operatorname{Im} N^{p-v-1} .
$$

It suffices to check this on the modules $K_{s}^{\prime}$. Then Ker $N^{p}$ has basis $\left\{N^{i} x \mid i \geq\right.$ $s+1-p\}$ and $\operatorname{Im} N^{p-v-1}$ has basis $\left\{N^{i} x \mid i \geq p-v-1\right\}$. Thus $N^{i} x \in F_{v} K$ if and only if, for some $p, i \geq \max (s+1-p, p-v-1)$, i.e. $i+v+1 \geq p \geq s+1-i$, thus if and only if $s-2 i \leq v$.

A linear operator $L\left(x_{i}\right)=\sum a_{i, j} x_{j}$ on a vector space $K$ defines a linear differential operator $D_{L}:=\sum a_{i, j} x_{j} \partial / \partial x_{i}$, which acts on the symmetric algebra $S(K)$ of $K$, and induces the action of $L$ on $K$. We also regard $D_{L}$ as a vector field on $K$, and then denote it by $\xi_{L}$.

Over a field of characteristic zero, we can also form the 1-parameter group $\{\operatorname{Exp}(t L)\}$ of automorphisms of $K$, which inherits an action on $S(K)$. If we have a Lie algebra $\mathfrak{g}$ of linear automorphisms of $K$, the exponentials generate a group $G$ of automorphisms of $K$, and hence of $S(K)$, and the induced action of an element $L$ of the Lie algebra is that of $D_{L}$.

We have seen how to extend a nilpotent operator $N$ on $K$ to an action of $s l_{2}$ : this now extends to an action of $s l_{2}$ on the symmetric algebra $S(K)$ of $K$, which in turn we can restrict to the homogeneous part $M:=S_{d}(K)$ of degree $d$.

If further $K$ splits as a direct sum $K^{\prime} \oplus K^{\prime \prime}$ with each of $K^{\prime}, K^{\prime \prime}$ invariant under $N$, extending as above to an $s l_{2}$-action leaves each summand $s l_{2}$-invariant. The induced actions of $S L_{2}$ and $s l_{2}$ on $S(K)$ now preserve each of the subspaces $M=S_{d^{\prime}}\left(K^{\prime}\right) \otimes S_{d^{\prime \prime}}\left(K^{\prime \prime}\right)$. Corresponding remarks apply to a direct sum of three or more summands.

Applying Lemma 1.1 to $M$, we find

Theorem 1.3. Let $N$ be a nilpotent endomorphism of $K$. Then
(i) for any $r \geq 0, D_{N}^{r}$ gives an isomorphism from the weight space $M_{r}$ to $M_{-r}$,
(ii) if $w>0$, then $\operatorname{Ker}\left(D_{N} \mid M_{w}\right)=0$.
(iii) if $w<0$, then $\operatorname{dim} \operatorname{Ker}\left(D_{N} \mid M_{w}\right)=\operatorname{dim} M_{w}-\operatorname{dim} M_{w-2}$.
(iv) $\operatorname{Ker}\left(D_{N} \mid M_{0}\right)$ is the space of invariants of $S L_{2}$ acting on $M$.

## 2 Restrictions on unipotent actions

Let $K$ be a finite dimensional vector space over $\mathbb{C}$ with a nilpotent endomorphism $N$, of type given by the sequence $R=\left\{r_{1} \geq r_{2} \geq \cdots\right\}$. We consider homogeneous functions $f$ on $K$, of degree $d$, annihilated by $D_{N}$, or equivalently, invariant under the 1-parameter group $G_{N}=\{\operatorname{Exp}(t N)\}$. We seek the conditions under which the hypersurface $V$ in $P(K)$ defined by $f$ is quasi-smooth. In this section we will enumerate the possibilities for $(d ; R)$.

Let $\left\{x_{i}\right\}(1 \leq i \leq n)$ be variables with assigned weights $w\left(x_{i}\right)=w_{i}$, arranged in non-decreasing order. Define a filtration of the polynomial ring $\mathbb{C}[x]$ by letting $f \in F(v)$ if $f$ is a linear combination of monomials of weights $\leq v$.

Lemma 2.1. Let $f_{j}, j=1, \ldots, m$ be polynomials of degree $D$ in the $x_{i}$ with $f_{j} \in$ $F\left(W_{j}\right)$; we suppose $W_{1} \leq \cdots \leq W_{m}$. Suppose that the set $Z$ of common zeroes of the $f_{j}$ in affine $n$-space has dimension $\leq k$. Then $D w_{i} \leq W_{i+m+k-n}$ for $i=1, \ldots, n-k$.

Proof. If all the $f_{j}$ vanish on $\left\{x_{1}=\cdots=x_{n-k-1}=0\right\}$ then $\operatorname{dim} Z \geq k+1$; so one of the $f_{j}$, say $f_{j_{1}}$, contains a monomial in $x_{n-k}, \ldots, x_{n}$ alone, and so $D w_{n-k} \leq W_{j_{1}}$. If all but $f_{j_{1}}$ vanish on $\left\{x_{1}=\cdots=x_{n-k-2}\right\}$ then $\operatorname{dim} Z \geq k+1$; so another of the $f_{j}$, say $f_{j_{2}}$, contains a monomial in $x_{n-k-1}, \ldots, x_{n}$, and so $D w_{n-k-1} \leq W_{j_{2}}$. Continuing in this way we find distinct $j_{1}, j_{2}, \ldots, j_{n-k} \in\{1, \ldots, m\}$ s.t. $D w_{n-k+1-i} \leq W_{j_{i}}$.

Since the numbers $j_{s}$ for $1 \leq s \leq n-k-i+1$ are all distinct, at least one of them, say $j_{\ell}$, is $\leq m-n+k+i$, by the pigeonhole principle. Hence $D w_{i} \leq$ $D w_{n-k+1-\ell} \leq W_{j_{\ell}} \leq W_{m-n+k+i}$.

Corollary 2.2. Let $f$ be homogeneous of degree $d$ in the variables $x_{i}$; suppose $f \in$ $F(W)$ and that the singular set of the variety $V$ has dimension $\leq k-1$. Then $(d-1) w_{i} \leq W-w_{n+1-i-k}$ for $1 \leq i \leq n-k$. In particular, if $W=0$ and the set of weights $w_{i}$ is symmetric about 0 , we have $(d-1) w_{i} \leq w_{i+k}$.

For set $f_{i}:=\partial f / \partial x_{i}$. Then $f_{i}$ has degree $d-1$ and $f_{i} \in F\left(W-w_{i}\right)$. Rearranging these numbers in increasing order gives $W_{j}=W-w_{n+1-j}$. The singular set of $V$ has dimension $\leq k-1$ if and only if the locus of common zeros of the $f_{i}$ in affine space has dimension at most $k$. Applying the lemma shows that in this case, $(d-1) w_{i} \leq W_{i+k}$ for $i=1, \ldots, n-k$, i.e. $(d-1) w_{i} \leq W-w_{n+1-i-k}$.

Lemma 2.3. Let $f$ be homogeneous of degree $d$ in the variables $p_{i}$; suppose that each monomial occurring in $f$ has weight $\leq 0$; suppose also that the hypersurface $f=0$ is quasi-smooth. Then
(i) $f$ contains a monomial of degree $d-1$ in the two variables of highest weight,
(ii) $f$ contains two monomials, each of degree $d-1$ in the three variables of highest weight, with the other factors different.

Proof. (i) Write $\mathfrak{m}_{2}$ for the ideal generated by all variables other than the two of highest weight. If $f \in \mathfrak{m}_{2}^{2}$, the hypersurface $f=0$ is singular along the line corresponding to these two co-ordinates. Otherwise, $f$ must contain a monomial of degree $d-1$ in them and containing just one other co-ordinate.
(ii) Write $\mathfrak{m}_{3}$ for the ideal generated by all variables other than the three of highest weight. Each term in $f$ not belonging to $\mathfrak{m}_{3}^{2}$ has degree $d-1$ in these and contains just one other co-ordinate. If this other co-ordinate is the same in all cases, say $X$, we can write $f=X \phi_{d-1}\left(p_{1}, p_{2}, p_{3}\right)+R$, with $R \in \mathfrak{m}_{3}^{2}$. But then the hypersurface $f=0$ is singular along the curve $\phi_{d-1}\left(p_{1}, p_{2}, p_{3}\right)=0$ in the plane defined by $\mathfrak{m}_{3}$.

We now apply these results to the problem of hypersurfaces invariant by a unipotent group.

Theorem 2.4. Let $V: f=0$ be a quasi-smooth hypersurface of degree $d>2$ in projective space, which is invariant under the action of a unipotent group of type given by the sequence $R=\left\{r_{1} \geq r_{2} \geq \ldots\right\}$. Then we have one of the the following:

$$
\begin{array}{llll}
\text { Case 2: } & d \geq 3, R=\{2\}, & \text { Case 3: } & d=4, R=\{3\}, \\
\text { Case 4: } & d=3, R=\{4\}, & \text { Case 21: } & d=3, R=\{2,1\} .
\end{array}
$$

Proof. For each $i$ we have basis elements of weights $-r_{i}, 2-r_{i}, \ldots, r_{i}$, and by Theorem 1.3, $f$ is a linear combination of monomials of weight $\leq 0$. Thus the hypotheses of the special case of Corollary 2.2 are satisfied.

If there is just one generator of positive weight, $r_{1}$ is 1 or 2 , and other $r_{i}=0$. If $r_{1}=1$, the ring of invariants is polynomial in the generators of weight $\leq 0$, so is independent of $x_{n}$, and defines a cone. If $r_{1}=2$, we have Case 2 of the theorem.

By Corollary 2.2 , we have $(d-1) w_{n-1} \leq w_{n}$. Now if $r_{2}=r_{1}$, we have $w_{n-1}=$ $w_{n}=r_{1}$, so $r_{1} \geq r_{1}(d-1)$, a contradiction. If $r_{2}=r_{1}-1$, we have $w_{n}=r_{1}, w_{n-1}=$ $r_{1}-1$, so $r_{1} \geq\left(r_{1}-1\right)(d-1)$ and $r_{1} \leq \frac{d-1}{d-2}$. If $d>3$ this implies $r_{1}=1$, a possibility we excluded above; if $d=3$ we may also have $r_{1}=2$. If now $r_{3}=1$, we have $w_{n-2}=1$, contradicting $(d-1) w_{n-2} \leq w_{n-1}$. Thus $r_{3}=0$, and we have Case 21 of the Theorem.

Otherwise we necessarily have $w_{n}=r_{1}$ and $w_{n-1}=r_{1}-2$, whence $r_{1} \geq$ $\left(r_{1}-2\right)(d-1)$ and $r_{1} \leq \frac{2(d-1)}{d-2}$. This gives $r_{1} \leq 4$ if $d=3, r_{1} \leq 3$ if $d=4$ and $r_{1} \leq 2$ if $d>4$.

The cases $r_{1}=2$ were considered above. If $d=4$, the remaining possibility is $r_{1}=3$, so that $w_{n}=3, w_{n-1}=1$. Since $3 w_{n-2} \leq w_{n-1}$, we have $w_{n-2}=0$, and Case 3 of the Theorem. It remains to consider the cases $d=3$ and $r_{1}$ equal to 3 or 4 .

If $r_{1}=3$, then $w_{1}=3, w_{2}=1$, so again $w_{3}=0$ and $r_{2}=0$. There is just one non-trivial Jordan block, which has size 4 and weights $-3,-1,1,3$ : denote the corresponding variables by $x_{0}, x_{1}, x_{2}, x_{3}$, and write $M$ for the space of homogeneous cubics in them. By Lemma $2.3, f$ must contain a term $x_{0} x_{2}^{2}$, which has weight -1 . Now apply Theorem 1.3 to $M$. Since $\operatorname{dim} M_{-3}=\operatorname{dim} M_{-1}=3, \operatorname{Ker}\left(D \mid M_{-1}\right)=0$. We thus have a contradiction. In fact the ring of invariants $\operatorname{Ker}(D)$ is given explicitly in Lemma 4.1, and the homogeneous invariants of degree 3 are linear in both the variables of positive weight.

If $d=3$ and $r_{1}=4$, then $w_{n}=4, w_{n-1}=2$, so $2 w_{n-2} \leq w_{n-1}=2$, hence $r_{2} \leq 1$. If $r_{2}=0$, we have Case 4 of the theorem.

If $\left(d, r_{1}, r_{2}\right)=(3,4,1)$, we have co-ordinates of positive weights $4,2,1$, so by Corollary 2.2 no others, so $r_{3}=0$. Write $K=K_{1} \oplus K_{2} \oplus K_{3}$, where $K_{1}$ is the Jordan block of $N$ of size 5 , with co-ordinates of respective weights $-4,-2,0,2,4$, which we denote $x_{0}, x_{1}, x_{2}, x_{3}, x_{4}, K_{2}$ the Jordan block of size 2 , with co-ordinates $y_{0}, y_{1}$ of weights $-1,1$; all the rest have weight 0 .

By Lemma 2.3, $f$ must contain the monomial $x_{0} x_{3}^{2}$ and the monomial $x_{1} y_{1}^{2}$. Write $N_{w}$ for the vector space spanned by monomials of weight $w$ of degree 1 on $K_{1}$ and degree 2 on $K_{2}$. Since each monomial of degree 2 in $K_{2}$ has weight 2,0 or -2 , there is a unique co-ordinate on $K_{1}$ with which we can multiply to attain weight 0 ; likewise to attain weight -2 . Thus $\operatorname{dim} \operatorname{Ker}\left(D \mid N_{0}\right)=\operatorname{dim} N_{0}-\operatorname{dim} N_{-2}=0$, so no appropriate invariant function exists.

It also follows from Lemma 2.3 that
Corollary 2.5. If $r_{1}=3$ (Case 3), $f$ must contain the monomial $x_{2}^{3} x_{0}$; if $r_{1}=4$ (Case 4), $f$ must contain the monomial $x_{3}^{2} x_{0}$; and if $r_{1}=2, r_{2}=1$ (Case 21), $f$ must contain the monomial $y_{1}^{2} x_{0}$.

## 3 Toolkit

Before we start detailed investigation of the cases listed above, we first recall some general results, and then collect some methods of calculation of Milnor numbers, so as not to break the thread of exposition in the following sections .

Let $V$ be quasi-smooth, with equation $f=0$ of degree $d>2$ in $P^{n}$. Recall that we call $f$, and the hypersurface $V$, oversymmetric if $f$ is annihilated by vector fields $\xi$ of degree 1 and $\eta$ of degree $d-2$ which is not a multiple of $\xi$. When $d=3$ this is equivalent to requiring $V$ to be 2 -symmetric. We recall the important result

Theorem 3.1. Suppose $V$ quasi-smooth of degree $d$ and $\xi$ a vector field of degree $r \leq d-2$ with $\xi(f)=0$. Then $\tau(V) \leq(d-1)^{n}-r(d-1-r)(d-1)^{n-2}$, and equality holds if and only if there is a vector field $\eta$ of degree $d-1-r$ with $\eta(f)=0$ and independent of $\xi$. Moreover when this holds, any vector field annihilating $f$ is a linear combination of $\xi, \eta$ and Hamiltonian vector fields.

This is the content of [6, Theorems 4.7, 4.9] when expressed in geometrical terms. Taking $r=1$, we obtain

Theorem 3.2. Suppose $V$ quasi-smooth and 1-symmetric of degree d with $\xi(f)=0$. Then $\tau(V) \leq(d-1)^{n-2}\left(d^{2}-3 d+3\right)$, and equality holds if and only if $V$ is oversymmetric, with a second vector field $\eta$. When this holds, any vector field annihilating $f$ is a linear combination of $\xi, \eta$ and Hamiltonian vector fields.

This gives the maximal value of $\tau$ for 1-symmetric, and conjecturally for all quasi-smooth hypersurfaces.

Corollary 3.3. The hypersurface $V$ cannot be 3-symmetric; it is 2-symmetric if and only if it is oversymmetric and $d=3$.

For by [6, Lemma 5.2], if $f$ is annihilated by vector fields $\xi, \xi^{\prime}$ with $\xi \wedge \xi^{\prime} \neq 0$, of degrees $r, r^{\prime}$ we must have $r+r^{\prime} \geq d-1$. If $V$ is 2-symmetric, we have $r=r^{\prime}=1$, hence $d=3$ and $V$ is oversymmetric; the converse is immediate. It follows from the theorem that now any vector field annihilating $f$ is a linear combination of $\xi, \xi^{\prime}$ and Hamiltonian vector fields; hence if linear, is a linear combination of $\xi$ and $\xi^{\prime}$.

The vector field $\xi$ is the infinitesimal generator of a linear group $G$. The cases when $G$ is semi-simple were discussed in our earlier paper [4], and the complete list of the oversymmetric cases was given in $[6, \S 5.3]$, and more fully in our survey article [5]. The symmetry group may be taken to act diagonally, so is determined by its weights. Either the only two non-zero weights are $\pm 1$, and the intersection of $V$ with the zero weight space is a cone; or there are just three non-zero weights, and the set of weights is obtained by adding zeros to a set of four weights; these must admit the monomials $x_{1}^{d}$, either (B) $x_{0} x_{2}^{d-1}$ or (C) $x_{0} x_{2}^{d-2} x_{3}$, and either $\left(\lambda_{r}\right) x_{0}^{r} x_{3}^{d-r}$, $\left(\mu_{r}\right) x_{0}^{r} x_{2} x_{3}^{d-r-1}$, or $\left(\nu_{r}\right) x_{0}^{r} x_{1} x_{3}^{d-r-1}$ for some $r$.

In this article we complete the list by determining all the cases when $G$ is unipotent.

We turn to calculations of Milnor numbers. We begin with Thom's splitting theorem (alias the Morse lemma with parameters). As we will need a precise version, we outline a proof.

Lemma 3.4. (a) Let $f\left(x_{1}, \ldots, x_{r}, y_{1}, \ldots, y_{k}\right)$ have 2-jet a non-degenerate quadratic form in $x_{1}, \ldots, x_{r}$. Suppose that (locally) the solution of the equations $\partial f / \partial x_{i}=0$ $(1 \leq i \leq r)$ is given by $x_{i}=\alpha_{i}\left(y_{1}, \ldots, y_{k}\right)(1 \leq i \leq r)$. Then $f$ is right-equivalent to $g(y)+\sum_{1}^{r} \pm z_{i}^{\prime 2}$, where $g(y)=f\left(\alpha_{1}(y), \ldots, \alpha_{r}(y), y_{1}, \ldots, y_{k}\right)$.
(b) Suppose $f\left(t, x, y_{1}, \ldots, y_{k}\right)$ is singular at the origin, with non-zero coefficient of $t x$ and that $\partial f / \partial t$ vanishes along $x=0$. Then $f$ is right-equivalent to $t x+$ $f\left(0,0, y_{1}, \ldots, y_{k}\right)$.
Proof. (a) It follows from our hypothesis that the hypersurfaces $\partial f / \partial x_{i}=0$ intersect transversely at $O$, so there is a solution of the form given. Substitute $x_{i}=z_{i}+\alpha_{i}(y)$ giving $f(x, y)=F(z, y)$, say. Then $\partial F / \partial z_{i}=\partial f / \partial x_{i}$ vanishes along $z_{1}=\cdots=$ $z_{r}=0$, and $F(0, y)=g(y)$. Hence $F(z, y)-g(y) \in\left\langle z_{1}, \ldots, z_{r}\right\rangle^{2}$. It follows in turn that we can write it as $\sum_{1}^{r} z_{i} h_{i}(y, z)$ with $h_{i} \in\left\langle z_{1}, \ldots, z_{r}\right\rangle$, and as $\sum_{i, j=1}^{r} z_{i} z_{j} k_{i, j}(y, z)$, where it follows from our hypothesis that the matrix $k_{i, j}(0)$ is non-singular. Now by 'completing the square' $r$ times we can write this in the form $\sum_{1}^{r} \pm z_{i}^{\prime 2}$, where $\left(z_{1}^{\prime}, \ldots, z_{r}^{\prime}, y_{1}, \ldots, y_{k}\right)$ can be taken as local co-ordinates at $O$.
(b) Since $\partial f / \partial t$ vanishes along $x=0$, we can write $\partial f / \partial t=x a(t, x, y)$ for some $C^{\infty}$-function $a$. Hence $f(t, x, y)-f(0, x, y)=\int_{0}^{t} x a(t, x, y) d t$, hence has the form $x b(t, x, y)$ for some $C^{\infty}$-function $b$. As also $f(0, x, y)-f(0,0, y)$ is divisible by $x$, we can write $f(t, x, y)=x c(t, x, y)+f(0,0, y)$. Now $c$ vanishes at the origin and has non-zero coefficient of $t$; thus the co-ordinate change $t^{\prime}=c(t, x, y)$ gives the desired equivalence.

Write $\chi_{n}(d)$ for the Euler characteristic of a smooth hypersurface of degree $d$ in $P^{n+1}$ : then (see e.g. [1, p. 152])

$$
\begin{equation*}
\chi_{n}(d)=n+2+\frac{(-1)^{n}}{d}\left((d-1)^{n+2}-(-1)^{n}\right) . \tag{1}
\end{equation*}
$$

When $n=-2,-1,0$ this formula gives $0,0, d$ respectively, so remains correct. The cone over such a hypersurface in $P^{n+2}$ admits a $\mathbb{C}^{*}$-action which is free except at the fixed points, which consist of the hypersurface itself and an isolated point. Thus this cone has Euler characteristic $\chi_{n}(d)+1$. If $V$ is a hypersurface of degree $d$ in $P^{n+1}$ with isolated singularities, then (see e.g. [1, p. 162]) $\chi(V)=\chi_{n}(d)+(-1)^{n-1} \mu(V)$.

In weighted projective space, suppose $f$, of degree $d$ with respect to weights $w_{i}$ $(0 \leq i \leq n+1)$ with sum $W$, defines a smooth hypersurface $V$, so $\operatorname{dim} V=n$. We have the following theorem of Steenbrink.

Theorem 3.5. (see [1, Theorem B34] and [8]). The mixed Hodge numbers of the primitive cohomology of $V$ are given by $h_{0}^{i, n-i}(V)=\operatorname{dim} M(f)_{d(i+1)-W}$, where $M(f)$ is the Milnor algebra

$$
M(f)=\mathbb{C}\left[x_{0}, \ldots, x_{n+1}\right] /\left\langle\partial f / \partial x_{0}, \ldots, \partial f / \partial x_{n+1}\right\rangle
$$

$M(f)$ has Euler-Poincaré polynomial $p(t)=\prod_{i=0}^{n+1}\left(1-t^{d-w_{i}}\right) /\left(1-t^{w_{i}}\right)$. Thus the primitive Betti number $h_{0}^{n}(V)$, which is the sum of the Hodge numbers $h_{0}^{i, n-i}(V)$, is equal to the sum of the coefficients of $p(t)$ in degrees congruent to $-W$ modulo $d$. This sum is given by $\frac{1}{d} \sum \epsilon^{W} p(\epsilon)$, where $\epsilon$ runs through the $d^{t h}$ roots of unity. It follows that

$$
\begin{equation*}
\chi(V)=n+1+\frac{(-1)^{n}}{d} \sum_{\epsilon^{d}=1} \epsilon^{W} p(\epsilon) \tag{2}
\end{equation*}
$$

We now evaluate this in the two cases we will need.
Corollary 3.6. A non-singular hypersurface $V$ of degree $d$ in weighted projective space with weights $w_{0}=2$ and $w_{i}=1$ for $1 \leq i \leq n+1$ has

$$
\chi(V)=n+2+\frac{(-1)^{n}}{d}\left\{\frac{d-2}{2}(d-1)^{n+1}+(-1)^{n+1}\right\}
$$

if $d$ is odd, and is $\frac{1}{2}$ less than this if $d$ is even. In particular, if $d=3$ we have $n+2+\frac{(-1)^{n}}{3}\left\{2^{n}-(-1)^{n}\right\}$.

In this case $W=n+3$, so $p(t)=\left(1-t^{d-2}\right)\left(1-t^{d-1}\right)^{n+1} /\left(1-t^{2}\right)(1-t)^{n+1}$. If $\epsilon=1$ we evaluate $p$ by l'Hôpital's rule, obtaining $p(1)=\frac{1}{2}(d-2)(d-1)^{n+1}$. If $\epsilon=-1$, which is only possible if $d$ is odd, we have $\left(1-\epsilon^{d-1}\right) /(1-\epsilon)=1$, so $p(-1)=(-1)^{n+1} \frac{1}{2}(d-2)$. Otherwise we have $\left(1-\epsilon^{d-1}\right) /(1-\epsilon)=-\epsilon^{-1}$ and $\left(1-\epsilon^{d-2}\right) /\left(1-\epsilon^{2}\right)=-\epsilon^{-2}$, so $p(\epsilon)=\epsilon^{W}\left(-\epsilon^{-2}\right)\left(-\epsilon^{-1}\right)^{n+1}$, which reduces to $(-1)^{n}$. We thus obtain $n+1+(-1)^{n} \frac{1}{d}\left\{\frac{1}{2}(d-2)(d-1)^{n+1}+(-1)^{n}(d-1)\right\}$ if $d$ is odd, and $n+1+(-1)^{n} \frac{1}{d}\left\{\frac{1}{2}(d-2)(d-1)^{n+1}+(-1)^{n+1} \frac{1}{2}(d-2)+(-1)^{n}(d-2)\right\}$ if $d$ is even, which reduces to the values stated.

If $V$ has isolated singularities then, as above, we must add $(-1)^{n-1} \mu(V)$ to this expression. If the singularities occur at smooth points of the ambient weighted projective space, this is proved as before, using the additive nature of $\chi$; at other points, we may take it as the definition of $\mu$.

Corollary 3.7. A non-singular hypersurface $V$ of degree 6 in weighted projective space with weights $w_{0}=w_{1}=1$ and $w_{i}=2$ for $1<i \leq n+1$ has

$$
\chi(V)=n+1+\frac{1}{6}(-1)^{n}\left(26 \cdot 2^{n}+4(-1)^{n}\right)
$$

In this case, $W=2 n+4$ and $M(\phi)$ has Euler-Poincaré polynomial

$$
p(t)=\left(1-t^{5}\right)^{2}\left(1-t^{4}\right)^{n} /(1-t)^{2}\left(1-t^{2}\right)^{n}=\left(1+t+t^{2}+t^{3}+t^{4}\right)^{2}\left(1+t^{2}\right)^{n} .
$$

We have $p(1)=5^{2} .2^{n}$ and $p(-1)=2^{n}$. If $\epsilon^{6}=1$ and $\epsilon^{2} \neq 1,\left(1-\epsilon^{5}\right) /(1-\epsilon)=-\epsilon^{-1}$ and $\left(1-\epsilon^{4}\right) /\left(1-\epsilon^{2}\right)=-\epsilon^{-2}$, so $p(\epsilon)=(-1)^{n} \epsilon^{-2-2 n}$. Hence $\sum \epsilon^{2 n-4} p(\epsilon)=26.2^{n}+$ $2^{n}+4(-1)^{n}$.

## 4 Invariants of unipotent actions

Now consider a nilpotent endomorphism $N$ of a vector space $K$; we adopt as our standard notation $\xi_{N}=\sum_{i=1}^{k} x_{i-1} \partial / \partial x_{i}$. We write $G_{N}$ for the Lie group obtained by exponentiating, $E_{t}=\operatorname{Exp}(t N)$, then $t . \mathbf{x}=E_{t} \mathbf{x}$; thus a polynomial $f$ on $K$ is annihilated by $\xi_{N}$ if and only if it is invariant under $G_{N}$. Write $\mathcal{I}_{k}$ for the ring of invariants of the group $G_{N}$ (a subring of $\mathbb{C}\left[x_{0}, \ldots, x_{k}\right]$ ). If we have a second Jordan block, denote the variables $y_{0}, \ldots, y_{l}$, set $\xi=\sum_{i=1}^{k} x_{i-1} \partial / \partial x_{i}+\sum_{i=1}^{l} y_{i-1} \partial / \partial y_{i}$, and write $\mathcal{I}_{k, l}$ for the ring of invariants.

It is a classical theorem of Weitzenbock [11] that the ring of invariants is finitely generated. Weitzenbock also determined the localisation at $x_{0}$ of the ring of invariants. Indeed, if $x_{0} \neq 0$, there is a unique choice $t_{0}=-\frac{x_{1}}{x_{0}}$ of the parameter $t$ such that $(t . \mathbf{x})_{1}=0$. Then all the $X_{i}=\left(t_{0} \cdot x\right)_{i}$ for $2 \leq i \leq k$ are invariants, and clearly $\mathcal{I}_{k}\left[x_{0}^{-1}\right]=\mathbb{C}\left[x_{0}, X_{2}, \ldots, X_{k}, x_{0}^{-1}\right]$. The argument also applies if $N$ has several Jordan blocks.

This remark can be used to compute the structure of the ring of invariants. Weitzenbock himself did this for $\operatorname{dim} K \leq 4$; a general algorithm was given by Tan [9], and a fuller account is in the book of Nowicki [7]. The results we need can be stated as follows.

Lemma 4.1. We have rings of invariants

$$
\begin{aligned}
\mathcal{I}_{2} & \cong \mathbb{C}[X, B] \\
\mathcal{I}_{3} & \cong \mathbb{C}\left[X, B, C, \Delta / X^{2} \Delta+C^{2}+B^{3}=0\right] \\
\mathcal{I}_{4} & \cong \mathbb{C}\left[X, B, U, C, E / X^{3} E=3 X^{2} B U-B^{3}-C^{2}\right] \\
\mathcal{I}_{2,1} & \cong \mathbb{C}\left[X, Y, T, B, S / X S=Y^{2} B+T^{2}\right]
\end{aligned}
$$

where $X:=x_{0}, Y:=y_{0}$, and

$$
\begin{aligned}
B & :=T_{x, x}^{2}=2 x_{0} x_{2}-x_{1}^{2}, \\
C & :=3 x_{0}^{2} x_{3}-3 x_{0} x_{1} x_{2}+x_{1}^{3}, \\
\Delta & :=-9 x_{0}^{2} x_{3}^{2}+18 x_{0} x_{1} x_{2} x_{3}-8 x_{0} x_{2}^{3}+3 x_{1}^{2} x_{2}^{2}-6 x_{1}^{3} x_{3}, \\
U & :=T_{x, x}^{4}=2 x_{0} x_{4}-2 x_{1} x_{3}+x_{2}^{2}, \\
E & :=12 x_{0} x_{2} x_{4}-9 x_{0} x_{3}^{2}+6 x_{1} x_{2} x_{3}-2 x_{2}^{3}-6 x_{1}^{2} x_{4}, \\
T & :=T_{x, y}^{1}=x_{0} y_{1}-x_{1} y_{0}, \\
S & :=x_{0} y_{1}^{2}-2 x_{1} y_{0} y_{1}+2 x_{2} y_{0}^{2} .
\end{aligned}
$$

Here $\mathcal{I}_{3}$ was given in [11], $\mathcal{I}_{4}$ in [9] and $\mathcal{I}_{2,1}$ in [7]. For Cases 2 and 3 we follow the notation of [4].

For the geometric problem, we have additional variables $w=\left(w_{1}, \ldots, w_{m}\right)$, all invariant. Thus the dimension $n=m+2, m+3, m+4$ or $m+4$ in our 4 cases respectively. Denote the corresponding elements of the ring of invariants by $W$ := $\left(W_{1}, \ldots, W_{m}\right)$.

We can use changes of co-ordinates that are compatible with $N$ to simplify our formulae.

Lemma 4.2. In Case 4, the co-ordinate changes compatible with $N$ are: $x_{4}^{\prime}=$ $\sum_{0}^{4} a_{i} x_{4-i}+\sum e_{j} w_{j}, x_{3}^{\prime}=\sum_{0}^{3} a_{i} x_{3-i}, x_{2}^{\prime}=\sum_{0}^{2} a_{i} x_{2-i}, x_{1}^{\prime}=\sum_{0}^{1} a_{i} x_{1-i}, x_{0}^{\prime}=a_{0} x_{0}$, $w_{i}^{\prime}=\sum p_{i, j} w_{j}+q_{i} x_{0}$, where $a_{0} \neq 0$ and $\left(p_{i, j}\right)$ is non-singular.

For we have taken an arbitrary element of $K$ for $x_{4}^{\prime}$; then $x_{3}^{\prime}=N x_{4}^{\prime}, x_{2}^{\prime}, x_{1}^{\prime}$ and $x_{0}^{\prime}$ are determined. Since our change of co-ordinates must respect the filtration, $w_{i}^{\prime}$ must be as stated. For the formulae to define a co-ordinate change we must have $a_{0} \neq 0$ and ( $p_{i, j}$ ) non-singular.

The results in Cases 2 and 3 are almost the same, and Case 21 is very similar.
In the next four sections we give detailed discussions of the four cases of Theorem 2.4 in turn.

## 5 Case 2

We define a map $\pi: K \rightarrow L$ taking as target co-ordinates $(W, X, B)$. This induces a map $\bar{\pi}: P(K) \rightarrow P(L)$, where $P(L)$ is the weighted projective space with all weights 1 except $w(B)=2$. The map $\bar{\pi}$ is defined except on the set $\mathcal{E}$ where all co-ordinates other than $x_{2}$ vanish. Thus $\mathcal{E}$ is a point, which we also denote $P$. The space $P(L)$ has just one singular point, where all co-ordinates except $B$ vanish: we denote it by $Q$.

As is usual for moduli spaces, we have a natural stratification. We define strata $\mathcal{S}_{i}$ in $K$ and $\overline{\mathcal{S}}_{i}$ in $L$, compatible with each other under $\pi$ and with passage to projective space.

$$
\begin{array}{ll}
\mathcal{S}_{0}: x_{0} \neq 0 ; & \overline{\mathcal{S}}_{0}: X \neq 0 \\
\mathcal{S}_{1}: x_{0}=0, x_{1} \neq 0 ; & \overline{\mathcal{S}}_{1}: X=0, B \neq 0 \\
\mathcal{S}_{2}: x_{0}=x_{1}=0 ; & \overline{\mathcal{S}}_{2}: X=B=0 .
\end{array}
$$

The set $\mathcal{F}$ of fixed points is defined by the vanishing of $x_{0}, x_{1}$, so coincides with $\mathcal{S}_{2}$. Each orbit of the action of $G_{N}$ on $K$ or on $P(K)$ outside $\mathcal{F}$ is isomorphic to an affine line; their degrees are 2,1 for $\mathcal{S}_{0}, \mathcal{S}_{1}$ respectively.

For any $(W, X, B) \in L$, we calculate $\pi^{-1}(W, X, B)$.
In each case, we have uniquely $w=W, x_{0}=X$.
If $X \neq 0, x_{1}$ is free (i.e. can be chosen arbitrarily), and $x_{2}=\left(B+x_{1}^{2}\right) / 2 X$ : we have one orbit.

If $X=0, x_{1}= \pm \sqrt{ }(-B)$, and $x_{2}$ is free. If $B \neq 0$, this gives two orbits, but if $B=0$, a line of fixed points.

From this we infer (with some care) pre-images under $\bar{\pi}$; in each case, we tabulate the Euler characteristic of the pre-image.
Lemma 5.1. [4, Lemma 6.4] The preimage $\bar{\pi}^{-1}(W, X, B)$ is as follows:
$\overline{\mathcal{S}}_{0}$ ) one orbit, $\chi=1$,
$\left(\overline{\mathcal{S}}_{1}\right)$ if $W \neq 0$, two orbits, $\chi=2$; if $W=0$ (the point $Q$ ), one orbit, $\chi=1$,
$\overline{\mathcal{S}}_{2}$ ) infinitely many point orbits, $\chi=1$.
Since the ring of invariants is a polynomial ring, any invariant function $f$ is of the form $f=\phi \circ \pi$, where $\phi=\phi(W, X, B)$ is a polynomial function on $L$. Set $\phi_{B}:=\frac{\partial \phi}{\partial B}$, $\phi_{X}:=\frac{\partial \phi}{\partial X}$.

Denote by $V$ the hypersurface in $P(K)$ defined by $f$, by $V_{1}$ the hypersurface in the weighted projective space $P(L)$ defined by $\phi$, and by $V_{2}$ and $V_{3}$ the intersections of $V_{1}$ with $X=0$ and with $X=B=0$ respectively. As in similar cases below, our notation is chosen so that each $V_{r}$ (also $V_{r}^{*}$, etc.) has dimension $m+1-r$.
Lemma 5.2. (compare [4, Lemma 5.5]) $V$ has isolated singular points if and only if $V_{1}$ has no singular points and $V_{3}$ has isolated singular points. The singular points of $f$ are $P$ and points $P_{i}$ corresponding to the singular points $Q_{i}$ of $V_{3}$ at which $\phi_{B} \neq 0$.
Proof. At a critical point of $f$, the following vanish:

$$
\frac{\partial f}{\partial w_{i}}=\frac{\partial \phi}{\partial W_{i}}, \quad \frac{\partial f}{\partial x_{0}}=\phi_{X}-2 x_{2} \phi_{B}, \quad \frac{\partial f}{\partial x_{1}}=2 x_{1} \phi_{B}, \quad \text { and } \quad \frac{\partial f}{\partial x_{2}}=-2 x_{0} \phi_{B} .
$$

If $\phi_{B}=0$, we have a critical point of $\phi$. If $W=X=B=0$, the only corresponding point in $P(K)$ is $P$. Otherwise we have a singular point of $V_{1}$. Conversely, if we have a singular point of $V_{1}$, all the points in its pre-image are singular on $V$, so are non-isolated singular points of $V$.

For a critical point of $f$ with $\phi_{B} \neq 0$, we have $x_{0}=x_{1}=0$, hence $X=B=0$, and a critical point of the restriction of $\phi$ to $X=B=0$. If $W=0$, we again have the point $P$. Otherwise we have a singular point of $V_{3}$. Conversely, if we have a singular point of $V_{3}$ at which $\phi_{B} \neq 0$, there is a unique corresponding value of $x_{2}$ giving a critical point of $f$, hence a unique corresponding singular point of $V$.

However, if we have a singular point of $V_{3}$ at which $\phi_{B}=0$, then $\phi_{X} \neq 0$ as otherwise we would have a singular point of $V_{1}$; and as $\phi_{X} \neq 0$, there is no corresponding critical point of $f$.

It also follows that if $V$ has isolated singular points, so also has $V_{2}$. For if the singular locus of $V_{2}$ had positive dimension, it would have to intersect the hypersurface $\phi_{X}=0$, and any point of intersection gives a singular point of $V_{1}$. Observe also that if $d$ is even, $Q \notin V_{2}$, for otherwise $Q$ would be a singular point of $V_{1}$. If $d$ is odd, then there can be no term $B^{d / 2}$, so $Q \in V_{2}$.

We can easily describe the singularities of $V$ at points other than $P$.
Proposition 5.3. The singularity at $P_{i}$ of $V$ corresponding to a singularity at $Q_{i}$ of $V_{3}$ at which $\phi_{B} \neq 0$ is right-equivalent to a suspension of that singularity.
Proof. We may suppose, after an allowable co-ordinate change, that at the singular point $P_{i}$ we have $0=x_{2}=\phi_{X}$. Apply Lemma 3.4(a) to $f$ with the variables $\left(x_{0}, x_{1}, x_{2}\right)$. We observe that $\partial f / \partial x_{0}, \partial f / \partial x_{1}, \partial f / \partial x_{2}$ all vanish when $x_{0}=x_{1}=0$ and $x_{2}=\phi_{X} / 2 \phi_{B}$. Substituting these values gives $g(w)=f(w, 0,0,0)$. The result follows.

We are now ready to calculate $\mu(V)$.
Theorem 5.4. We have $\mu(V)=\frac{1}{2}(d-2)(2 d-1)(d-1)^{m}+\mu\left(V_{2}\right)+\mu\left(V_{3}\right)$.
Proof. First suppose $d$ odd. Then as $Q \in V_{1}$, by Lemma 5.1 we have
$\chi(V)=\chi(\mathcal{E})+\chi\left(V_{1} \backslash V_{2}\right)+2 \chi\left(V_{2} \backslash V_{3}\right)-1+\chi\left(V_{3}\right)=\chi\left(V_{1}\right)+\chi\left(V_{2}\right)-\chi\left(V_{3}\right)$. For $V$ and $V_{3}$ we apply (1); for $V_{1}$ we apply Corollary 3.6 with $n=m$; for $V_{2}$ the same, but with $n=m-1$, and amended for singularities. Thus

$$
\begin{aligned}
& \chi(V)=\frac{(-1)^{m+1}}{d}\left\{(d-1)^{m+3}-(-1)^{m+1}\right\}+m+3+(-1)^{m} \mu(V), \\
& \chi\left(V_{1}\right)=\frac{(-1)^{m}}{d}\left\{\frac{d-2}{2}(d-1)^{m+1}+(-1)^{m+1}\right\}+m+2, \\
& \chi\left(V_{2}\right)=\frac{(-1)^{m-1}}{d}\left\{\frac{d-2}{2}(d-1)^{m}+(-1)^{m}\right\}+m+1+(-1)^{m} \mu\left(V_{2}\right), \\
& \chi\left(V_{3}\right)=\frac{(-1)^{m}}{d}\left\{(d-1)^{m}-(-1)^{m}\right\}+m+(-1)^{m-3} \mu\left(V_{3}\right) .
\end{aligned}
$$

Since $\chi(V)-\chi\left(V_{1}\right)-\chi\left(V_{2}\right)+\chi\left(V_{3}\right)$ vanishes, we find that $\mu(V)-\mu\left(V_{2}\right)-\mu\left(V_{3}\right)$ is equal to $d^{-1}\left\{(d-1)^{m+3}+\frac{d-2}{2}(d-1)^{m+1}-\frac{d-2}{2}(d-1)^{m}-(d-1)^{m}\right\}$, which reduces to $\frac{1}{2}(d-2)(2 d-1)(d-1)^{m}$.

In the case $d$ even, as $Q \notin V_{1}$, we obtain $\chi(V)=\chi\left(V_{1}\right)+\chi\left(V_{2}\right)-\chi\left(V_{3}\right)+1$, but the values of each of $\chi\left(V_{1}\right)$ and $\chi\left(V_{2}\right)$ are $\frac{1}{2}$ less than those above. Hence the formula in terms of $d$ is the same as before.

If $m=0$, the value of $\mu(V)$ is given by [2, Proposition 3.1]: the value $\frac{1}{2}(d-$ $2)(2 d-1)$ is correct if $d$ is even; we must add $\frac{1}{2}$ if $d$ is odd. Here if $d$ is odd $V_{2}$ is necessarily singular: indeed, $\phi$ vanishes identically on $\overline{\mathcal{S}}_{1}$.

For the case $m=1,\left[4\right.$, Prop 6.6] gives $\mu(V)=\frac{1}{2}(d-2)(2 d-1)(d-1)+k-N$, where $k=\left\lfloor\frac{d}{2}\right\rfloor$ and $N$ is the number of distinct points of $V_{2} \backslash V_{3}$ with $W \neq 0$. In this case, $V_{2}$ has dimension 0 , so $\chi\left(V_{2}\right)=\# V_{2}$; by Corollary 3.6, if $V_{2}$ is smooth, $\chi\left(V_{2}\right)=\left\lfloor\frac{1}{2}(d+1)\right\rfloor$. To reconcile these we need to interpret $\mu\left(V_{3}\right)$ as 1 if $V_{3} \neq \emptyset$, i.e. if the coefficient of $W^{d}$ in $\phi$ is non-zero.

When does $\mu(V)$ take its maximal value? We expect this to occur if $\mu\left(V_{2}\right)$ and $\mu\left(V_{3}\right)$ are both as large as possible, hence when $V_{2}$ and $V_{3}$ are both cones. Geometry imposes restrictions as follows.

As already observed, $Q \in V_{2}$ if and only if $d$ is odd.
If $V_{2}$ is a cone with vertex not in $V_{3}$, then it is a cone on $V_{3}$. Since it must have isolated singularities, $V_{3}$ must be non-singular.

If $V_{2}$ is a cone with vertex different from $Q$, then $Q \notin V_{2}$ (since the local geometry at $Q$ differs from that elsewhere).

Thus if $d$ is even, while $Q \notin V_{2}$, we cannot exclude the possibility that $V_{2}$ is a cone with vertex $Q^{*}$ in $V_{3}$. When this holds, $V_{3}$ also is a cone, so indeed we expect $\mu(V)$ to be maximal. The singularity of $V_{2}$ at $Q^{*}$ is equisingular to a sum of $(m-1) d^{t h}$ powers and a $\frac{1}{2} d^{t h}$ power, so $\mu\left(V_{2}\right)=\frac{1}{2}(d-1)^{m-1}(d-2)$; that of $V_{3}$ at $Q^{*}$ is equisingular to a sum of $d^{t h}$ powers, so $\mu\left(V_{3}\right)=(d-1)^{m-1}$ and $\mu(V)=\frac{1}{2}(d-1)^{m-1}\left(2 d^{3}-7 d^{2}+8 d-2\right)$. Since $\phi_{B}\left(Q^{*}\right)=0$, the only singular point of $V$ is $P$.

This case does indeed occur for all $m \geq 1$ and even $d \geq 4$ : we can take $\phi=$ $X^{d}+X W_{1}^{d-1}+B^{d / 2}+\sum_{2}^{m} W_{i}^{d}$. Then each of $V_{2}$ and $V_{3}$ is a cone, smooth except at the point $Q^{*}$ where all co-ordinates except $W_{1}$ vanish; and $V_{1}$ is non-singular.

We believe these to give the maximal values of $\mu(V)$ for all even $d \geq 4, m \geq 1$ : for $m=1$ this follows from [4, Proposition 6.6].

If $d$ is odd and $V_{2}$ is a cone, then the vertex of the cone is $Q$ and $V_{3}$ is nonsingular. In affine co-ordinates $B=1, \phi$ is equisingular to a sum of $m d^{t h}$ powers, which suggests $\mu_{Q}\left(V_{2}\right)=(d-1)^{m}$, but since we must factor out the antipodal map on affine space we actually have $\mu_{Q}\left(V_{2}\right)=\frac{1}{2}(d-1)^{m}$, thus $\mu(V)=\frac{1}{2}(d-1)^{m+1}(2 d-3)$. This case occurs for all $d=2 k+1$ : we can take $\phi=X B^{k}+X^{d}+\sum_{1}^{m} W_{i}^{d}$. Again this gives the maximum value for $\mu$ if $m=1$ : we cannot show that this holds in general.

We now treat the case $d=3$ in more detail. Here $f=a_{1} B+a_{3}$, where $a_{1}, a_{3}$ are homogeneous functions of $x_{0}, w_{1}, \ldots, w_{m}$. If $a_{1}$ is not a multiple of $x_{0}$, we can make a change of co-ordinates to write $a_{1}=w_{1}$; otherwise we can take $a_{1}=x_{0}$ (if $a_{1} \equiv 0$, $V$ has non-isolated singularities). Denote by $V_{3}^{*}$ the variety $w_{1}=a_{3}=0$ in $P^{m}$ and set $V_{4}:=V_{3} \cap V_{3}^{*}$.
Lemma 5.5. (i) If $a_{1}=x_{0}, V$ is quasi-smooth if and only if $V_{3}$ is non-singular. In this case, the only singular point is $P$.
(ii) If $a_{1}=w_{1}$, there is a bijection between singular points of $V_{2}$ and $V_{4}$; the singularity of the former is isomorphic to the suspension of the latter.
(iii) If $a_{1}=w_{1}, V_{1}$ is non-singular if and only if $V_{3}^{*}$ is non-singular.

Proof. (i) In this case, $V_{2}$ is a cone with vertex $Q$, so the result follows as above.
(ii) At a singular point $(W, 0, B)$ of $V_{2}$ we have $0=\phi_{B}=W_{1}$ and $0=\partial \phi / \partial W_{i}$ for $i \geq 2$, so $(W, 0,0)$ is a singular point of $V_{4}$ (note that $(0,0, B)$ is not a singular point of $V_{2}$ since $\partial \phi / \partial W_{1}$ does not vanish there).

Conversely, if $(W, 0,0)$ is a singular point of $V_{4}$, the point $(W, 0, B)$ is singular on $V_{2}$ if and only if $B=-\partial a_{3} / \partial W_{1}$.

Now apply Lemma 3.4(a), taking $B$ and $W_{1}$ as the preferred co-ordinates. The equations $\phi_{B}=\partial \phi / \partial W_{1}=0$ are solved by $W_{1}=0, B=-\partial a_{3} / \partial W_{1}$ : substituting these in $\phi=w_{1} B+a_{3}$ gives the restriction of $a_{3}$ to $W_{1}=0$.
(iii) The same argument as for (ii) applies here.

If $a_{1}=x_{0}, V_{2}$ is a cone with vertex $Q$, so $\mu(V)=3.2^{m}$ as above. If $a_{1}=w_{1}$, it follows that $\mu\left(V_{2}\right)=\mu\left(V_{4}\right)$, so by Theorem 5.4, $\mu(V)=5.2^{m-1}+\mu\left(V_{3}\right)+\mu\left(V_{4}\right)$.

In some cases, we can determine the nature of the singularities.
Proposition 5.6. Suppose $V$ in Case 2 with $V_{2}$ non-singular. Then the singularity of $V$ at $P$ is semi-quasi-homogeneous with degree $2 d$ and variables of weights 1,4 and 2 ( $m$ times).

Proof. We first give a direct argument, then an indirect method, which only determines the $\mu$-constant stratum, but will be usable in other cases below.

As in [4, Proposition 6.6], take local affine co-ordinates $x_{2}=1$ at $P$, and substitute $x_{0}=\frac{1}{2}\left(B+x_{1}^{2}\right)$, so that $f$ becomes $\phi\left(w, B, \frac{1}{2}\left(B+x_{1}^{2}\right)\right)$. Now assign weights 4 to $B, 1$ to $x_{1}$ and 2 to the $w_{i}$. The terms of least weight $2 d$ give $\phi\left(w, B, \frac{1}{2} x_{1}^{2}\right)$. We must check that this has an isolated critical point. At a singular point, $\partial \phi / \partial w_{i}, \phi_{B}$ and $x_{1} \phi_{X}$ vanish. Since $V_{1}$ is non-singular, $x_{1}=0$, so $(W, 0, B)$ is a singular point of $V_{2}$, contradicting our hypothesis.

For our second argument, we note that by Theorem 5.4, $\mu(V)-\mu\left(V_{3}\right)$ takes the same value for all these cases (with $d$ and $m$ fixed). Now $\mu(V)=\mu_{P}(V)+\sum_{i} \mu_{P_{i}}(V)$, and by Proposition 5.3, the values $\mu_{P_{i}}(V)=\mu_{Q_{i}}\left(V_{3}\right)$. Hence $\mu_{P}(V)$ is the same for all these cases, so all belong to the same $\mu$-constant stratum.

To determine this, first observe that we can adjoin a new variable $w_{m+1}$; then $f^{\prime}:=f+w_{m+1}^{d}$ again satisfies the conditions of Lemma 5.2, and the new singularity is obtained from the old one also by adjoining a new variable and adding its $d^{\text {th }}$ power. Hence the $\mu$-constant type of the singularity can be deduced from the case with $m$ decreased by 1 .

If $m=0$, as observed above, if $d$ is odd, $V_{2}$ is necessarily singular; if $d$ is even, the result holds. However the case $m=1$ was analysed in [4, Proposition 6.6], where we showed directly that the singularity has the type stated. The result thus follows in general.

The second method can also be applied to the case when $V_{2}$ is a cone with vertex $Q$. We see that the singularity is equisingular to a sum of $m d^{t h}$ powers and the curve singularity occurring in the case $m=0$, which can be taken to be $\prod_{1}^{k}\left(y^{2}-2 x+4 c_{i} x^{2}\right)$ if $d=2 k$ and $x \prod_{1}^{k}\left(y^{2}-2 x+4 c_{i} x^{2}\right)$ if $d=2 k+1$ (with the $c_{i}$ all distinct in each case). If $d \geq 5$, it is not quasi-homogeneous.

If $d=3$, the cases arising when $a_{1}=w_{1}$ can be enumerated in low dimensions by considering the varieties $V_{4} \subset V_{3}$. We can determine the $\mu$-constant strata of the critical points of $f$ using the fact that the terms of lowest weight are $\phi\left(w, B, \frac{1}{2} x_{1}^{2}\right)$, which reduces by splitting to $\phi\left(0, w_{2}, \ldots, w_{m}, 0, x_{1}^{2}\right)$, together with our calculation of $\mu$.

For $m=1$, we have $A_{5}$ at $P$, perhaps a further $A_{1}$.
For $m=2$, we have the cubic curve $a_{3}\left(w_{1}, w_{2}, x_{0}\right)=0$ meeting $x_{0}=0$ in $V_{3}$ and the point $Q^{*}$ where $w_{1}=x_{0}=0$ in $V_{4}$. Let $w_{1}^{r}$ be the highest power of $w_{1}$ dividing $a_{3}\left(w_{1}, w_{2}, 0\right)$. If $r=0, V_{4}=\emptyset, V_{2}$ is non-singular and $V$ has a $T_{2,3,6}$ at $P$ and a further $A_{1}$ (or $A_{2}$ ) if $V_{3}$ has a repeated point (a 3 -fold point). If $r \geq 1$, we apply the same substitution, but must now use the 2 -jet $B w_{1}$ and obtain the splitting by direct calculation. The first substitution gives the 4 -jet $w_{2}^{2} x_{1}^{2}+\alpha w_{2}^{4}$, where $\alpha=0$ if $r>1$. Thus if $r=1$ the singularity has $\mu=11$, hence type $T_{2,4,6}$, and the other two points on $V_{3}$ could coincide, giving a further $A_{1}$. If $r=2$ we have a singularity $T_{2, p, q}$ with $p, q \geq 5$ and $\mu=12$, hence $p+q=11$ so $(p, q)=(5,6)$. In the case $r=3$ we have $p+q=12$, and need a further calculation to obtain the 5 -jet, leading to $p=q=6$. Thus in each case, we have $T_{2,3+r, 6}$ at $P$.

For $m=3$, if $V_{4}$ is non-singular (i.e. 3 points), $U T_{0,0,0}^{1}$ (in the notation of [10, p. 475]) together with $-, A_{1}, A_{2}, 2 A_{1}, A_{3}, 3 A_{1}$ or $D_{4}$. If $V_{4}$ is singular, we have non-reduced 3 -jet ( $V$ or $V^{\prime}$ series) and the singularities do not have accepted names.

*     *         * 

We turn to calculation of $\tau$ : here our results are much more partial. According to Lemma 3.2, $\tau$ takes its maximal value when $f$ is oversymmetric. To find when this is applicable, we use the method of $[4, \S 6]$.

Theorem 5.7. A function in Case 2 is oversymmetric if and only if either (a) $V_{3}$ is a cone, or (b) after change of co-ordinates if necessary, $\phi_{B}$ and $\phi_{X}$ both vanish along $X=B=0$.

Proof. Since $V_{1}$ is non-singular, the sequence $\left\{\partial \phi / \partial W_{1}, \ldots, \partial \phi / \partial W_{m}, \phi_{X}, \phi_{B}\right\}$ is regular, and any vector field annihilating $\phi$ is a linear combination of the Hamiltonian fields $\partial(\phi, *) / \partial\left(W_{i}, W_{j}\right), \partial(\phi, *) / \partial\left(W_{i}, X\right), \partial(\phi, *) / \partial\left(W_{i}, B\right)$ and $\partial(\phi, *) / \partial(X, B)$. We seek a vector field $\eta$ which is a lift of a linear combination of these. We are only interested in $\eta$ modulo Hamiltonian fields: removing the corresponding linear combination of the $\partial(f, *) / \partial\left(w_{i}, w_{j}\right)$ and $\partial(f, *) / \partial\left(w_{i}, x_{0}\right)$, we can take $\sum C_{i} \partial(\phi, *) / \partial\left(W_{i}, B\right)$ $+D \partial(\phi, *) / \partial(X, B)$. Since we seek $\eta$ of degree $d-2$, we want the $C_{i}$ and $D$ to be constants. We now have

$$
\eta=\sum_{1}^{m} p_{i} \partial / \partial w_{i}+\sum_{0}^{2} q_{j} \partial / \partial x_{j}
$$

where $p_{i}=-C_{i} \phi_{B}, q_{0}=-D \phi_{B}$ and

$$
2\left(x_{2} q_{0}-x_{1} q_{1}+x_{0} q_{2}\right)=\sum C_{i} \partial \phi / \partial W_{i}+D \phi_{X} .
$$

Thus

$$
\begin{equation*}
2\left(x_{0} q_{2}-x_{1} q_{1}\right)=\sum C_{i} \partial \phi / \partial W_{i}+D \phi_{X}+2 x_{2} D \phi_{B} \tag{3}
\end{equation*}
$$

The right hand side of this equation must thus vanish identically along $X=B=0$.
First suppose $D=0$. Changing the $w$ co-ordinates, we may suppose the vector field is $\partial / \partial W_{1}$. If we set $a_{d}(W):=\phi(W, X, B)$, we need $\partial a_{d} / \partial W_{1} \equiv 0$, i.e. $a_{d}$ is independent of $W_{1}$. Expressing the condition geometrically, it holds if and only if $V_{3}$ is a cone.

If $D \neq 0$, a suitable substitution $W_{i}^{\prime}:=W_{i}+\lambda_{i} X, X^{\prime}:=X$ reduces the $C_{i}$ to zero, so it suffices to consider the vector field $\partial(\phi, *) / \partial(X, B)$. Here the condition reduces to requiring both $\phi_{B}$ and $\phi_{X}$ to vanish along $X=B=0$.

We could reformulate (b) as: there exist constants $c_{i}$ such that $\phi_{B}$ and $\phi_{X}+$ $\sum c_{i} \partial \phi / \partial W_{i}$ both vanish along $X=B=0$.

This proof shows more generally that any vector field $\eta=\sum_{1}^{m} p_{i} \partial / \partial w_{i}+$ $\sum_{0}^{2} q_{j} \partial / \partial x_{j}$ annihilating $f$ can be reduced modulo Hamiltonian vector fields to the lift of $\sum C_{i} \partial(\phi, *) / \partial\left(W_{i}, B\right)+D \partial(\phi, *) / \partial(X, B)$, where $p_{i}=-C_{i} \phi_{B}, q_{0}=-D \phi_{B}$ and (3) holds. Moreover, we may suppose $D$ and the $C_{i}$ independent of $x_{0}$ and $x_{1}$. This can be used as the starting point for further calculations of $\tau$. However, since the cases arising are diverse, we only consider $m=0, m=1$ and certain cases with $d=3$.

In the case of curves $(m=0)$, the condition frequently holds, and then $\tau=$ $d^{2}-3 d+3$ (see [2, Proposition 3.1]): otherwise, $\tau=d^{2}-3 d+2$.

The case of surfaces $(m=1)$ was treated in [4]. By Theorem 6.7 loc.cit., $\tau_{\text {tot }}(V)=(d-1)\left(d^{2}-3 d+3\right)$ if $\alpha=0$ or $\gamma=0$, and $(d-1)\left(d^{2}-3 d+3\right)-1$ otherwise; where $\alpha, \gamma$ are the coefficients of $W^{d}$ and $B W^{d-2}$ in $\phi$. Moreover (Lemma 6.5 loc.cit.) $P$ is the only singular point unless $\alpha=0 \neq \gamma$, when there is one further singular point, of type $A_{1}$. The case $\alpha=0$ corresponds to clause (a) of the Theorem; the case $\gamma=0$ to clause (b) (here we appear to require $\beta=\gamma=0$ : the difference arises because of the above normalisation of co-ordinates).

We can calculate $\tau(V)$ ad hoc in further low dimensional cases. When $d=3$, if $a_{1}=w_{1}$ the values can be inferred from the above list of $\mu$-constant strata: we have $\tau=\mu$ for $T_{2,3,6}$ and $\tau=\mu-1$ for $T_{2, p, q}$ with $\frac{1}{p}+\frac{1}{q}<\frac{1}{2}$. If $a_{1}=x_{0}$ we have

Lemma 5.8. If $V$ is in Case 2, with $d=3$ and $a_{1}=x_{0}$, then

$$
\tau(V)=2^{m+1}+\operatorname{dim}\left(\mathbb{C}\left[x_{0}, w_{1}, \ldots, w_{m}\right] /\left\langle x_{0}, \partial a_{3} / \partial x_{0}, \partial a_{3} / \partial w_{1}, \ldots, \partial a_{3} / \partial w_{n}\right\rangle\right)
$$

Proof. Here $\phi=X B+a_{3}(W, X)$, so (3) reduces to

$$
\sum_{i} C_{i} \partial a_{3} / \partial W_{i}+D \partial a_{3} / \partial X \in\langle X\rangle
$$

By Lemma 5.5(i), in this case $V_{3}$ is non-singular. Hence the restrictions of the $\partial a_{3} / \partial W_{i}$ to $X=0$ form a regular sequence in $\mathbb{C}[W]$, spanning an ideal $J$. The class of $\eta$ modulo Hamiltonian fields and multiples of $\xi$ is determined by the class of $\left.D\right|_{X=0}$ modulo $J$.

The algebra $\mathbb{C}[W] / J$ is Gorenstein of dimension $2^{m}$, with $\binom{m}{r}$ basis elements in degree $r$. If the ideal in it generated by the class of $\partial a_{3} /\left.\partial X\right|_{X=0}$ has dimension $e$, its annihilator has dimension $2^{m}-e$. Since the space of multiples of $\xi$ modulo Hamiltonian vector fields has dimension $2^{m+1}$, we obtain $\tau(V)=3.2^{m}-e$.

But $\mathbb{C}\left[x_{0}, w_{1}, \ldots, w_{m}\right] /\left\langle x_{0}, \partial a_{3} / \partial x_{0}, J\right\rangle \cong \mathbb{C}[W] /\left(J+\left\langle\partial a_{3} /\left.\partial X\right|_{X=0}\right\rangle\right)$, so has dimension $2^{m}-e$. The result follows.

If $m=2$ we can take $a_{3}\left(W_{1}, W_{2}, 0\right)=W_{1}^{3}+W_{2}^{3}$ and see easily that if the coefficient of $W_{1} W_{2} X$ in $a_{3}$ is non-zero, $e=0$ and $\tau=12$ : otherwise $e=1$ and $\tau=11$. If $m=3$ we take $W_{1}^{3}+W_{2}^{3}+W_{3}^{3}+3 \alpha W_{1} W_{2} W_{3}$ : here either $e=0$ or $e=2$. For $m \geq 4$, cases are more numerous.

We observe that while there are numerous cases where $\tau(V)$ takes its maximal value (for given dimension $n$ and degree $d$ ) but $\mu(V)$ does not, we do not know an example in the reverse direction. Indeed, $\mu$ maximal implies $\tau$ maximal for curves, surfaces of degree 4 or odd, and for cubic 3 -folds. If $V_{2}$ is a cone with vertex in $V_{3}$ then $\tau$ is maximal; but if it is a cone with vertex $Q$, while $\phi_{B}$ vanishes along $X=B=0$ we have no control on $\phi_{X}$.

We now give a more detailed discussion of the 2 -symmetric case $d=3$, following the notation of the above proof.

Proposition 5.9. In Case 2, $f$ is 2-symmetric only in the following 3 cases:
Case (b): we have $a_{1}=X . V_{2}$ is a cone with vertex $Q$ and $V_{3}$ is non-singular. After a suitable substitution $x_{2}^{\prime}:=x_{2}+\frac{1}{2} b\left(w, x_{0}\right), a_{3}$ is independent of $x_{0}$, and we may take $\eta=-2 x_{0} \partial / \partial x_{0}+x_{1} \partial / \partial x_{1}+4 x_{2} \partial / \partial x_{2}$. The singularity has $\mu=3.2^{m}$ and is quasi-homogeneous of degree 12 with respect to weights 3,6 and 4 ( $m$ times).

Case (a1): we have $a_{1}=W_{1}, V_{3}$ is a cone with vertex not on $W_{1}=0$ and $V_{4}$ non-singular. After a further substitution $x_{2}^{\prime}:=x_{2}+\frac{1}{2} b\left(w, x_{0}\right)$, may suppose $a_{3}$ independent of $w_{1}$. Then $f$ is invariant by $\eta=x_{1} \partial / \partial x_{1}+2 x_{2} \partial / \partial x_{2}-2 w_{1} \partial / \partial w_{1}$. There are two singularities, with Milnor numbers $5.2^{m-1}$ and $2^{m-1}$; both quasi-homogeneous of degree 6, the first with respect to weights 1 and 2 ( $m-1$ times); the second with respect to weights 3 and 2 ( $m-1$ times).

Case (a2): we have $a_{1}=W_{1} ; V_{3}, V_{4}$ are cones with vertex on $W_{1}=0$. After a further substitution $w_{2}^{\prime}:=b\left(w, x_{0}\right)$, we may suppose $a_{3}-x_{0} w_{2}^{2}$ independent of $w_{2}$. Then $f$ is invariant by $\eta=w_{2} \partial / \partial x_{2}-w_{1} \partial / \partial w_{2}$. The singularity has $\mu=13.2^{m-2}$, and is in the same $\mu$-constant stratum as $x^{6}+x^{2} y^{2}+y^{6}+\sum_{2}^{m} w_{i}^{3}$.

Proof. The enumeration is given in Theorem 5.7. For Case (b), we must have $a_{1}=$ $X$; it follows that $V_{2}$ is a cone with vertex $Q, V_{3}$ is non-singular and the singularity was determined above. Now write $a_{3}\left(w, x_{0}\right)=x_{0}^{2} b_{1}\left(w, x_{0}\right)+x_{0} b_{2}(w)+b_{3}(w)$. Since $\phi_{X}$ vanishes along $X=B=0, b_{2}$ vanishes identically, so the substitution $x_{2}^{\prime}=$ $x_{2}+\frac{1}{2} b_{1}\left(w, x_{0}\right)$ reduces $a_{3}$ to $b_{3}$, independent of $x_{0}$. That $f$ is now invariant under $\eta$ (so we can take $D=2$ and all $C_{i}=0$ ) follows by inspection. We could also infer the singularities from the semi-simple group action.

If Case (a) ( $V_{3}$ is a cone) occurs, $a_{1}$ cannot be $x_{0}$ (else $V_{3}$ would be non-singular), so can be taken as $w_{1}$. We must distinguish according as the vertex of the cone does or does not lie on $w_{1}=0$.

In Case (a1) it does not, so the intersection $V_{4}$ of the cone with $w_{1}=0$ is non-singular, hence so is $V_{2}$. The description of the singularities now follows from Propositions 5.3 and 5.6, or again from the group action.

After adjusting the $w$ co-ordinates, we may suppose $a_{3}(w, 0)$ independent of $w_{1}$. Then we can write $a_{3}=b_{3}+x_{0} w_{1} b_{1}\left(w, x_{0}\right)$, with $b_{3}$ independent of $w_{1}$. Again the substitution $x_{2}^{\prime}=x_{2}+\frac{1}{2} b_{1}\left(w, x_{0}\right)$ reduces $a_{3}$ to $b_{3}$. Now by inspection, $\eta f=0$ (so we may take $C_{1}=2, C_{i}=0$ for $i \neq 1$ ).

In case (a2), we may suppose $a_{3}(w, 0)$ independent of $w_{2}$, and hence that $a_{3}=$ $b_{3}+x_{0} w_{2} b_{1}\left(w, x_{0}\right)$, with $b_{3}$ independent of $w_{2}$. Now if the coefficient of $w_{2}$ in $b_{1}$ were zero, the point where all co-ordinates except $w_{2}$ vanish would be singular on $V_{3}^{*}$. Hence we can write $b_{1}=c^{2} w_{2}+c_{1}$, and substitute $w_{2}^{\prime}=c w_{2}+\frac{1}{2} c^{-1} c_{1}$, which reduces $a_{3}$ to the form $b_{3}^{\prime}+x_{0} w_{2}^{2}$, with $b_{3}^{\prime}$ independent of $w_{2}$. Thus $\eta f=0$, where $\eta=w_{2} \partial / \partial x_{2}-w_{1} \partial / \partial w_{2}$ (so we may take $C_{2}=2, C_{i}=0$ for $i \neq 2$ ).

To describe the singularity, as in Proposition 5.6, it suffices to consider the case $m=2$. Here since $V_{2}$ and $V_{3}$ are cones with the same vertex, we must select $T_{2,6,6}$ from the above list.

## 6 Case 3

We define the map $\pi: K \rightarrow L$ by $\pi\left(w, x_{0}, x_{1}, x_{2}, x_{3}\right)=(W, X, B, \Delta)$ in the notation of Lemma 4.1. This induces $\bar{\pi}: P(K) \rightarrow P(L)$, where $P(L)$ is the weighted projective space with all weights 1 except $w(B)=2, w(\Delta)=4$; the map $\bar{\pi}$ is defined except on the set $\mathcal{E}$ where all co-ordinates except $x_{2}$ and $x_{3}$ vanish: thus $\mathcal{E}$ is a projective line and $\chi(\mathcal{E})=2$; it contains the point $P$ where all co-ordinates except $x_{3}$ vanish. We define strata by

$$
\begin{array}{ll}
\mathcal{S}_{0}: x_{0} \neq 0 ; & \overline{\mathcal{S}}_{0}: X \neq 0 \\
\mathcal{S}_{1}: & x_{0}=0, x_{1} \neq 0 ; \\
\overline{\mathcal{S}}_{1}: X=0, B \neq 0 \\
\mathcal{S}_{2}: & x_{0}=x_{1}=0 ;
\end{array} \overline{\mathcal{S}}_{2}: X=B=0 . ~ l
$$

The set $\mathcal{F}$ of fixed points is given by the vanishing of $x_{0}, x_{1}, x_{2}$. Each orbit of the action of $G_{N}$ on $K \backslash \mathcal{F}$ (or on $P(K) \backslash \mathcal{F}$ ) is isomorphic to an affine line; their degrees are $3,2,1$ for $\mathcal{S}_{0}, \mathcal{S}_{1}, \mathcal{S}_{2} \backslash \mathcal{F}$ respectively.

We now describe the pre-image under $\pi$ of any $(W, X, B, \Delta) \in L$. In each case, $w=W$.
$\left(\mathcal{S}_{0}\right)$ if $X \neq 0, x_{0}=X, x_{1}$ is free, $x_{2}=\left(B+x_{1}^{2}\right) / 2 X, x_{3}=\left(C+3 X x_{1} x_{2}-x_{1}^{3}\right) / 3 X^{2}$ where $C= \pm \sqrt{ }\left(-X^{2} \Delta-B^{3}\right)$.
$\left(\mathcal{S}_{1}\right)$ if $X=0, B \neq 0, x_{1}= \pm \sqrt{ }(-B), x_{2}$ is free, and $x_{3}=\left(3 x_{1}^{2} x_{2}^{2}-\Delta\right) / 6 x_{1}^{3}$.
$\left(\mathcal{S}_{2}\right)$ if $X=B=0$ : if $\Delta \neq 0$, the pre-image is empty; if $\Delta=0, x_{1}=0$ and $x_{2}, x_{3}$ are arbitrary. If $x_{2} \neq 0$ we have a non-trivial orbit; if $x_{2}=0$ we have fixed points.

From this we infer (again with some care)
Lemma 6.1. The preimage $\bar{\pi}^{-1}(W, X, B, \Delta)$ is as follows:
$\left(\mathcal{S}_{0}\right)$ if $B^{3}+X^{2} \Delta=0$, one orbit, $\chi=1$, if not, two orbits, $\chi=2$,
$\left(\mathcal{S}_{1}\right)$ if $W=0$, one orbit, $\chi=1$, if not, two orbits, $\chi=2$,
$\left(\mathcal{S}_{2}\right)$ if $\Delta=0$, a plane, $\chi=1$, if not, the empty set, $\chi=0$.
A priori the map $f$ need not factor through $\pi$. However, we have
Lemma 6.2. There is an allowable change of co-ordinates which puts $f$ in the form $f=\phi \circ \pi$. More precisely, we may take $\phi=\Delta+a_{0} B^{2}+a_{2} B+a_{4}$, where $a_{i}$ is homogeneous of degree $i$ in $W, X$.

Proof. Since $d=4$, we can write $f=a_{0}^{\prime} \Delta+a_{0} B^{2}+a_{1} C+a_{2} B+a_{4}$, where $a_{i}$ is homogeneous of degree $i$ in the invariant co-ordinates $w_{1}, \ldots, w_{m}, x_{0}$.

By Corollary 2.5 , for $V$ to be quasi-smooth, $f$ must contain the monomial $x_{2}^{3} x_{0}$; so we must have $a_{0}^{\prime} \neq 0$. We may thus suppose $a_{0}^{\prime}=1$. Substituting $x_{3}^{\prime}:=x_{3}-\frac{1}{6} a_{1}(w, x)$ gives an expression of the same form but with $a_{1}=0$. This gives $f=\Delta+a_{0} B^{2}+$ $a_{2} B+a_{4}$, which is indeed of the form $\phi \circ \pi$.

Denote by $V$ the hypersurface $f=0$ in $P(K)$, by $V_{0}$ the hypersurface $\phi=0$ in $P(L)$, by $V_{1}$ its intersection with $X=0$, and by $V_{3}$ its intersection with $X=$ $B=\Delta=0$. Write also $V_{1}^{*}$ and $V_{2}^{*}$ for the respective intersections of $V_{0}$ and $V_{1}$ with $B^{3}+X^{2} \Delta=0$. Write $L^{\prime}$ for the vector space with co-ordinates $W, X, Z$ (all of degree 1), $P\left(L^{\prime}\right)$ for the corresponding projective space, $\psi_{1}(W, X, Z):=-X Z^{3}+$ $a_{0} X^{2} Z^{2}+a_{2} X Z+a_{4}=0, V_{1}^{\prime}$ for the hypersurface defined by $\psi_{1}$ in $P\left(L^{\prime}\right)$, and $V_{2}^{\prime}$ for its intersection with $X=0$.

Lemma 6.3. The singular points of $V$ are isolated iff the hypersurfaces $V_{1}^{\prime}, V_{3}$ are both non-singular; and then the only singular point of $V$ is $P$.

Proof. By Lemma 6.2, we can take $f=\Delta+a_{0} B^{2}+a_{2} B+a_{4}$. Since $\partial f / \partial x_{3}=$ $\partial \Delta / \partial x_{3}=-6 C, C$ vanishes at all critical points of $f$.

First consider critical points of $f$ in $x_{0} \neq 0$. Since each such point lies in a nontrivial orbit, and $f$ is invariant, it follows that if $f$ has isolated critical points, there can be none with $x_{0} \neq 0$. Now in this region, the critical points of $f$ are the same as those of $x_{0}^{2} f$, which is equal to $-C^{2}-B^{3}+x_{0}^{2}\left(a_{0} B^{2}+a_{2} B+a_{4}\right)$. These coincide with the critical points of $\phi_{0}:=-B^{3}+x_{0}^{2}\left(a_{0} B^{2}+a_{2} B+a_{4}\right)$ lying in $C=0$. Now regard $\phi_{0}$ as a function $\psi_{0}$ of the variables $W, X, B$. If this has a critical point with $X \neq 0$, we certainly have a critical point of $\phi_{0}$. Conversely, if we have a critical point of $\phi_{0}$, we have $0=\partial \phi_{0} / \partial x_{2}=2 x_{0} \partial \psi_{0} / \partial B$, so $0=\partial \psi_{0} / \partial B$, and in view of this, $\partial \psi_{0} / \partial X=\partial \phi_{0} / \partial x_{0}$ vanishes, and so do the $\partial \psi_{0} / \partial W_{i}$; so we have a critical point of $\psi_{0}$. Finally, in $X \neq 0$ we may make the substitution $Z:=X^{-1} B$. The critical points correspond, and since $\psi_{0}(W, X, X Z)=X^{2} \psi_{1}(W, X, Z)$, they correspond to
those of $\psi_{1}$. Thus $f$ has no critical points in $x_{0} \neq 0$ if and only if $\psi_{1}$ has none in $X \neq 0$.

Now consider critical points of $f$ with $x_{0}=0$. As $\partial f / \partial x_{3}=-6 C$, and $C$ reduces to $-x_{1}^{3}$, we must also have $x_{1}=0$, hence $B=C=\Delta=0$. It follows that $\partial f / \partial x_{2}=\partial f / \partial x_{1}=0$. There remain the conditions

$$
0=\partial f / \partial w_{i}=\partial a_{4} / \partial w_{i}, \quad 0=\partial f / \partial x_{0}=\partial a_{4} / \partial x_{0}+2 x_{2} a_{2}-8 x_{2}^{3}
$$

If we have a singular point of $V_{3}$, there is only one further equation to determine both $x_{2}$ and $x_{3}$ so we have a non-isolated singularity of $f=0$. Thus (ii) is a necessary condition for $f$ to have isolated singularities. If it holds, then for any critical points in $x_{0}=x_{1}=0$ we have $w=0$, and now the remaining equation implies $x_{2}=0$, giving the unique critical point $P$.

It remains to consider singular points of $V_{1}^{\prime}$ lying in $X=0$. Here $\partial \psi_{1} / \partial Z$ vanishes, $\partial \psi_{1} / \partial X=-Z^{3}+a_{2} Z+\partial a_{4} / \partial X$, and $\partial \psi_{1} / \partial W_{i}=\partial a_{4} / \partial W_{i}$. Since we are now assuming that $V_{3}$ has no singular points, the vanishing of the $\partial \psi_{1} / \partial W_{i}$ implies $W=0$ and that of $\partial \psi_{1} / \partial X$ then gives $Z=0$, showing that there are indeed no such singular points.

Theorem 6.4. For $V$ quasi-smooth in Case 3, $\mu(V)=22.3^{m}$.
Proof. As before, we calculate $\chi(V)$ by decomposing $\bar{\pi}(V)$ according to the stratification, calculating the Euler characteristic of each piece, inferring those of the pre-images, and adding up.

In $\overline{\mathcal{S}}_{2}, \Delta$ vanishes on the image of $\pi$ and $f$ reduces to $a_{4}(W, 0)$. The zero locus is thus the hypersurface $V_{3}$. Hence $\chi\left(V \cap \mathcal{S}_{2}\right)=\chi\left(V_{3}\right)+\chi(\mathcal{E})=\chi_{m-2}(4)+2$.

In $\overline{\mathcal{S}}_{1}$, we can assign $W$ and $B$ and solve $\Delta=-\left(a_{0} B^{2}+a_{2} B+a_{4}\right)$. The set where $W=0$ is a single point, so contributes $\chi=1$, and the set $W \neq 0$ is the product of the punctured $B$-plane and the punctured $W$ space, so has $\chi=0$. Hence $\chi\left(V \cap \mathcal{S}_{1}\right)=1$.

In $\overline{\mathcal{S}}_{0}$, we can normalise co-ordinates by $X:=1$. Projecting $V_{0} \backslash V_{1}$ onto ( $W, B$ ) space is an isomorphism, since $\Delta=-\left(a_{0} B^{2}+a_{2} B+a_{4}\right)$ on $V_{0}$. Thus $\chi\left(V_{0} \backslash V_{1}\right)=1$. Restricting to the subset where $B^{3}+\Delta=0$ we obtain an isomorphism of $V_{1}^{*} \backslash V_{2}^{*}$ onto the set of $(W, B)$ where $B^{3}=a_{0} B^{2}+a_{2} B+a_{4}$, which we can identify in turn (replacing $B$ by $Z$ ) with the subset $V_{1}^{\prime} \backslash V_{2}^{\prime}$ of $V_{1}^{\prime}$ with $X=1$. Thus $\chi\left(V_{1}^{*}\right)-\chi\left(V_{2}^{*}\right)=$ $\chi\left(V_{1}^{\prime}\right)-\chi\left(V_{2}^{\prime}\right)$. But $V_{1}^{\prime}$ is non-singular, so $\chi\left(V_{1}^{\prime}\right)=\chi_{m}(4)$, and $V_{2}^{\prime}$ is the cone on $V_{3}$, so $\chi\left(V_{2}^{\prime}\right)$ is equal to $\chi\left(V_{3}\right)+1=\chi_{m-2}(4)+1$. We thus have

$$
\chi\left(V_{1}^{*} \backslash V_{2}^{*}\right)=\chi_{m}(4)-\chi_{m-2}(4)-1
$$

and hence $\chi\left(V \cap \mathcal{S}_{0}\right)$ is twice $\chi\left(V_{0} \backslash V_{1}\right)$ minus this, i.e. $3-\chi_{m}(4)+\chi_{m-2}(4)$.
Adding these up, we find $\chi(V)=-\chi_{m}(4)+2 \chi_{m-2}(4)+6$. Now $V$ has just one singular point and $\operatorname{dim} V=m+2$, so $\chi(V)=\chi_{m+2}(4)+(-1)^{m-1} \mu(V)$. Thus finally $\mu(V)=(-1)^{m}\left\{\chi_{m+2}(4)+\chi_{m}(4)-2 \chi_{m-2}(4)-6\right\}$, which reduces to $22.3^{m}$.

We can also calculate $\tau$.
Proposition 6.5. Any $f$ in Case 3 is oversymmetric. Hence $\tau_{\text {tot }}(V)=3^{m+1} .7$.

Proof. We can use essentially the same formula as in [4, Proposition 6.3]. We have $f=\Delta+a_{0} B^{2}+a_{2} B+a_{4}$, with $a_{i}$ homogeneous of degree $i$ in $x_{0}$ and the $w_{i}$. Set $\eta^{\prime}:=x_{1}^{2} \partial / \partial x_{1}+\left(3 x_{1} x_{2}-3 x_{0} x_{3}\right) \partial / \partial x_{2}+\left(4 x_{2}^{2}-3 x_{1} x_{3}\right) \partial / \partial x_{3}$. Then $\eta^{\prime} \Delta=0$ and $\partial \Delta / \partial x_{3}=3 \eta^{\prime} B$. Hence $f$ is annihilated by $\eta^{\prime}-\frac{1}{3}\left(2 a_{0} B+a_{2}\right) \partial / \partial x_{3}$.

We contrast $\mu(V)=22.3^{m}$ with $\tau(V)=7.3^{m+1}=21.3^{m}$. The values $\mu=22, \tau=$ 21 were obtained in [4, Prop 6.3] for the case $m=0$.

Lemma 6.6. Suppose $V$ in Case 3. Then the singularity of $V$ is semi-quasi-homogeneous of degree 12 in variables of weights $1,4,6$ and 3 ( $m$ times).

Proof. Recall that $\Delta:=-9 x_{0}^{2} x_{3}^{2}+18 x_{0} x_{1} x_{2} x_{3}-8 x_{0} x_{2}^{3}+3 x_{1}^{2} x_{2}^{2}-6 x_{1}^{3} x_{3}$. When $x_{3}=1$, we can rewrite this as $\Delta=-\left(3 x_{0}-3 x_{1} x_{2}+\frac{4}{3} x_{2}^{3}\right)^{2}+6\left(\frac{2}{3} x_{2}^{2}-x_{1}\right)^{3}$. This suggests setting $q:=x_{0}-x_{1} x_{2}+\frac{4}{9} x_{2}^{3}, p:=x_{1}-\frac{2}{3} x_{2}^{2}$, so we substitute $x_{2}:=3 y$, $x_{1}:=p+6 y^{2}, x_{0}:=q+3 p y+6 y^{3}$. This gives $\Delta=-9 q^{2}-6 p^{3}, B=6 q y-p^{2}+6 p y^{2}$ and so $f=-9 q^{2}-6 p^{3}+a_{0}\left(6 q y-p^{2}+6 p y^{2}\right)^{2}+a_{2}\left(6 q y-p^{2}+6 p y^{2}\right)+a_{4}$, where $a_{i}$ is homogeneous of degree $i$ in $w_{1}, \ldots, w_{r}, q+3 p y+6 y^{3}$.

Now assign weight 1 to $y, 3$ to each $w_{i}, 4$ to $p$ and 6 to $q$. The term of least weight in $x_{0}$ is $6 y^{3}$, of weight 3 ; the term of least weight in $B$ is $6 p y^{2}$, of weight 6 . Hence each term in $f$ has weight at least 12, and the terms of degree 12 give the sum of a term $-9 q^{2}$, which we can ignore, and $g:=-6 p^{3}+36 a_{0} p^{2} y^{4}+6 a_{2} p y^{2}+a_{4}$, where $a_{i}$ is homogeneous of degree $i$ in $w_{1}, \ldots, w_{r}, 6 y^{3}$. It remains to show that $g$ has an isolated singularity.

We compare $g$ with the function $\psi_{1}(W, X, Z):=-X Z^{3}+a_{0} X^{2} Z^{2}+a_{2} X Z+a_{4}$, and observe that formally $g(w, p, y)=\psi_{1}\left(w, 6 y^{3}, p y^{-1}\right)$. Since by Lemma 6.3, the hypersurface $V_{1}^{\prime}$ defined by $\psi_{1}=0$ is non-singular, $g$ has no singular points with $y \neq 0$. But if $y=0$, the condition $\partial g / \partial p=0$ forces $p=0$; and the restriction to $p=y=0$ defines the hypersurface $V_{3}$ which, by the same result, is also non-singular. Hence indeed $g$ has an isolated singularity, and the result follows.

## $7 \quad$ Case 4

Here we define $\pi: K \rightarrow L$ by $\pi\left(w, x_{0}, x_{1}, x_{2}, x_{3}, x_{4}\right)=(W, X, B, U, E)$ in the notation of Lemma 4.1. The induced map $\bar{\pi}: P(K) \rightarrow P(L)$ is defined except on the set $\mathcal{E}$ where all co-ordinates except $x_{3}$ and $x_{4}$ vanish: $\mathcal{E}$ is a projective line containing the point $P$ where all co-ordinates except $x_{4}$ vanish and one other orbit, and $\chi(\mathcal{E})=2$. We define strata by

| $i$ | 0 | 1 | 2 | 3 |
| :---: | :---: | :---: | :---: | :---: |
| $\mathcal{S}_{i}$ | $x_{0} \neq 0$ | $x_{0}=0, x_{1} \neq 0$ | $x_{0}=x_{1}=0, x_{2} \neq 0$ | $x_{0}=x_{1}=x_{2}=0$ |
| $\overline{\mathcal{S}}_{i}$ | $X \neq 0$ | $X=0, B \neq 0$ | $X=B=0, U \neq 0$ | $X=B=U=0$. |

The set $\mathcal{F}$ of fixed points is given by the vanishing of $x_{0}, x_{1}, x_{2}, x_{3}$. Each orbit of the action of $G_{N}$ on $K \backslash \mathcal{F}$ (or on $P(K) \backslash \mathcal{F}$ ) is isomorphic to an affine line; their degrees are $4,3,2,1$ for $\mathcal{S}_{0}, \mathcal{S}_{1}, \mathcal{S}_{2}, \mathcal{S}_{3} \backslash \mathcal{F}$ respectively. The closure of each orbit in $P(K) \backslash \mathcal{F}$ is obtained by adjoining the point $P$.

We now describe the pre-image under $\pi$ of any $(W, X, B, U, E) \in L$. In each case, $w=W$ and $x_{0}=X$.
$\left(\mathcal{S}_{0}\right) x_{1}$ is free, $x_{2}=\left(B+x_{1}^{2}\right) / 2 X, x_{3}=\left(C+3 X x_{1} x_{2}-x_{1}^{3}\right) / 3 X^{2}$ where $C=$ $\pm \sqrt{ }\left(3 X^{2} B U-X^{3} E-B^{3}\right)$, and $x_{4}=\left(U+2 x_{1} x_{3}-2 x_{2}^{2}\right) / 2 X$;
$\left(\mathcal{S}_{1}\right) x_{1}= \pm \sqrt{ }(-B), x_{2}$ is free, $x_{3}=\left(x_{2}^{2}-U\right) / 2 x_{1}$, and $x_{4}=\left(6 x_{1} x_{2} x_{3}-\right.$ $\left.2 x_{2}^{3}-E\right) / 6 x_{1}^{2}$;
$\left(\mathcal{S}_{2}\right)$ if $E^{2} \neq 4 U^{3}$ the pre-image is empty; otherwise, $x_{1}=0, x_{2}=-E / 2 U, x_{3}$ and $x_{4}$ are free;
$\left(\mathcal{S}_{3}\right)$ if $E \neq 0$ the pre-image is empty; otherwise, $x_{1}=x_{2}=0, x_{3}$ and $x_{4}$ are free (if $x_{3} \neq 0$ we have a non-trivial orbit but if $x_{3}=0$ we have fixed points).

From this we infer pre-images under $\bar{\pi}$.
Lemma 7.1. For $(W, X, B, U, E) \in \overline{\mathcal{S}}_{i}$, the value $\chi\left(\bar{\pi}^{-1}(W, X, B, U, E)\right)$ is given by:

| $i$ | Condition | $\psi_{i}$ | $\chi\left(\right.$ if $\left.\psi_{i}=0\right)$ | $\chi\left(\right.$ if $\left.\psi_{i} \neq 0\right)$ |
| :---: | :---: | :---: | :---: | :---: |
| 0 | $X \neq 0$ | $B^{3}+X^{3} E-3 X^{2} B U$ | 1 | 2 |
| 1 | $X=0, B \neq 0$ | $W$ | 1 | 2 |
| 2 | $X, B=0, U \neq 0$ | $4 U^{3}-E^{2}$ | 1 | 0 |
| 3 | $X, B, U=0$ | $E$ | 1 | 0 |

Lemma 7.2. Suppose $f$, invariant under the group, defines a hypersurface $V$ with isolated singularities. Then there is an allowable change of co-ordinates which puts $f$ in the form $f=E+3 a_{1} U+a_{3}$, where $a_{i}$ is homogeneous of degree $i$ in the invariant co-ordinates.

Proof. We have $d=3$, and so can write $f=b_{0} C+a_{0} E+b_{1} B+3 a_{1} U+a_{3}$, where $a_{i}, b_{i}$ are homogeneous of degree $i$ in $w, x_{0}$. By Corollary 2.5, for $V$ to be quasi-smooth, $f$ must contain the monomial $x_{3}^{2} x_{0}$, so we must have $a_{0} \neq 0$. We may thus take $a_{0}=1$. Now substitute $x_{4}=x_{4}^{\prime}+\frac{b_{0}}{6} x_{1}-\frac{b_{1}}{6}$ and $x_{4}=x_{3}^{\prime}+\frac{b_{0}}{6} x_{0}$. This reduces $b_{0}$ and $b_{1}$ to 0 at the expense of adding terms to $a_{3}$. We thus have $f=E+3 a_{1} U+a_{3}$, of the desired form.

It will be convenient to write $a_{3}^{*}$ for $a_{1}^{3}+a_{3}$, and to give names to varieties as follows. We define $V_{2} \subset P^{m}$ by $a_{3}(W, X)=0, V_{2}^{*}$ by $a_{3}^{*}(W, X)=0$, their respective intersections with $X=0$ by $V_{3}, V_{3}^{*}$, and $V_{3} \cap V_{3}^{*}$ by $V_{4}$. In weighted projective space with coordinates $(W, X, Z, U)$ (where $U$ has weight 2) write $\psi_{1}:=$ $-Z^{3}+3 U Z+3 a_{1}(W, X) U+a_{3}(W, X)$; denote the hypersurface $\psi_{1}=0$ by $V_{0}$, and its intersection with $X=0$ by $V_{1}$.

Lemma 7.3. Suppose $f=E+3 a_{1} U+a_{3}$ as above. Then the singular points of $V$ are isolated iff
(i) $V_{0}$, or equivalently $V_{2}^{*}$ is non-singular,
(ii) $V_{1}$, or equivalently $V_{3}^{*}$ is non-singular, and
(iii) $V_{3}$ has isolated singular points.

The singular points of $f$ are then $P$ and points $P_{i}$ corresponding to the singular points $Q_{i}$ of $V_{3}$.

Proof. Since $x_{1} \partial f / \partial x_{4}+x_{0} \partial f / \partial x_{3}=x_{1} \partial E / \partial x_{4}+x_{0} \partial E / \partial x_{3}=-6 C, C$ again vanishes at all critical points of $f$.

Each singular point of $V$ in $x_{0} \neq 0$ lies in a non-trivial orbit, and $V$ is invariant, so $V$ can have no singular point, and $f$ can have no critical point with $x_{0} \neq 0$. In this region, the critical points of $f$ are the same as those of $x_{0}^{3} f$, which is equal to $3 x_{0}^{2} B U-B^{3}-C^{2}+x_{0}^{3}\left(3 a_{1} U+a_{3}\right)$. These coincide with the critical points of $\phi_{0}:=3 x_{0}^{2} B U-B^{3}+x_{0}^{3}\left(3 a_{1} U+a_{3}\right)$ lying in $C=0$. Now regard $\phi_{0}$ as a function $\psi_{0}$ of the variables $W, X, B, U$. If $\psi_{0}$ has a critical point with $X \neq 0$, we certainly have a critical point of $\phi_{0}$. Conversely, if we have a critical point of $\phi_{0}$, set

$$
\begin{gathered}
w_{i}:=W_{i}, \quad x_{0}:=X, \quad x_{1}:=t, \quad x_{2}:=\left(B+t^{2}\right) / 2 X, \\
x_{3}:=\left(3 x_{0} x_{1} x_{2}-x_{1}^{3}\right) / 3 x_{0}^{2}, \quad x_{4}:=\left(U+2 x_{1} x_{3}-x_{2}^{2}\right) / 2 x_{0},
\end{gathered}
$$

(so $C=0$ ). Then $0=\partial \phi_{0} / \partial x_{4}=2 x_{0} \partial \psi_{0} / \partial U$, so $0=\partial \psi_{0} / \partial U$, and $0=\partial \phi_{0} / \partial x_{2}=$ $2 x_{0} \partial \psi_{0} / \partial B+2 x_{2} \partial \psi_{0} / \partial U$, so $0=\partial \psi_{0} / \partial B$, and hence again, $\partial \psi_{0} / \partial X=\partial \phi_{0} / \partial x_{0}$ and the $\partial \psi_{0} / \partial W_{i}$ all vanish; so we have a critical point of $\psi_{0}$.

In $X \neq 0$ we may make the substitution $Z:=X^{-1} B$, then $\psi_{0}(W, X, X Z)=$ $X^{3} \psi_{1}(W, X, Z)$. Then the critical points of $\psi_{0}$ correspond to those of $\psi_{1}$. Thus $f$ has no critical points in $x_{0} \neq 0$ if and only if $\psi_{1}$ has none in $X \neq 0$; equivalently, $V_{0} \backslash V_{1}$ is non-singular.

For a critical point of $\psi_{1}, 0=\partial \psi_{1} / \partial U=3\left(Z+a_{1}\right)$, so $Z=-a_{1}$ and $0=$ $\partial \psi_{1} / \partial Z=3 U-3 Z^{2}$, so $U=Z^{2}=a_{1}^{2}$. With this value of $U$, the partial derivatives of $\psi_{1}$ with respect to $X$ and the $W_{i}$ coincide with those of $a_{1}^{3}+a_{3}$. Thus $f$ has no critical point in $x_{0} \neq 0$ if and only if $a_{3}^{*}$ has none in $X \neq 0$.

For singular points on $x_{0}=0$, we have $0=\partial f / \partial x_{4}=-6 x_{1}^{2}$, so $x_{1}=0$ also, and hence $U=x_{2}^{2}, E=-2 x_{2}^{3}$, so $f$ reduces to $-2 x_{2}^{3}+3 x_{2}^{2} a_{1}(w, 0)+a_{3}(w, 0)$. We now have $\partial f / \partial x_{4}=\partial f / \partial x_{3}=0$, and

$$
\begin{gathered}
\partial f / \partial x_{2}=6 x_{2}\left(a_{1}(w, 0)-x_{2}\right), \quad \partial f / \partial x_{1}=-6 x_{3}\left(a_{1}(w, 0)-x_{2}\right) \\
\partial f / \partial x_{0}=12 x_{2} x_{4}-9 x_{3}^{2}+6 x_{4} a_{1}(w, 0)+3 x_{2}^{2} \partial a_{1} / \partial x_{0}+\partial a_{3} / \partial x_{0}(w, 0) .
\end{gathered}
$$

If $x_{2}=a_{1}$, then $\partial f / \partial w_{i}=\partial\left(a_{1}^{3}+a_{3}\right) / \partial w_{i}=\partial a_{3}^{*} / \partial w_{i}$. If the restriction of $a_{3}^{*}$ to $x_{0}=0$ has a critical point, we can assign this value to the $w_{i}$, set $x_{2}=a_{1}$, and then only have one further equation in $x_{3}$ and $x_{4}$ : thus $f$ has non-isolated critical points. Thus condition (ii) is necessary. If it holds, then if $x_{2}=a_{1}$ we have $w=0$, hence in turn $x_{2}=x_{3}=0$ and we have the unique critical point $P_{4}$. Note also that (ii) implies that $a_{3}^{*}$ has no critical point on $X=0$, thus completing the proof of the necessity of (i).

For a critical point of $f$ with $x_{2} \neq a_{1}$, we must have $x_{2}=x_{3}=0$. Then $\partial f / \partial w_{i}$ reduces to $\partial a_{3} / \partial w_{i}$ and $\partial f / \partial x_{0}$ to $6 x_{4} a_{1}+\partial a_{3} / \partial x_{0}$. Thus we have a critical point of the restriction of $a_{3}$ to $x_{0}=0$; since $a_{1} \neq x_{2}=0$, each such critical point yields a unique value of $x_{4}$ and hence critical point of $f$.

Theorem 7.4. For $V$ quasi-smooth in Case 4, we have $\mu(V)=11.2^{m}+\mu\left(V_{3}\right)$.
Proof. As before, we calculate the $\chi\left(V \cap \mathcal{S}_{i}\right)$ using Lemma 7.1.
For $\mathcal{S}_{3}$ we only have to consider $a_{3}(w, 0)=0$, which defines $V_{3}$. Hence $\chi\left(V \cap \mathcal{S}_{3}\right)=$ $\chi\left(V_{3}\right)+\chi(\mathcal{E})$.

We know that $\chi\left(V \cap \mathcal{S}_{2}\right)$ is equal to the Euler characteristic of the set of $(W, U, E)$ with $\phi(W, 0,0, U, E)=0,4 U^{3}=E^{2}$ and $U \neq 0$. Since we can solve $\phi=0$ for $E$, it suffices to consider the set of $(W, U)$ where $4 U^{3}=\left(3 a_{1} U+a_{3}\right)^{2}$ and $U \neq 0$. We cannot have $W=0$ here, as this would imply $U=0$. We can thus project on the space $P^{m-1}$ with co-ordinates $W$. The fibre consists of the roots of the cubic equation in $U$, which has discriminant $16 a_{3}^{3} a_{3}^{*}$.

In the following table, the first column defines the subset of $P^{m-1}$, the second gives its Euler characteristic, the third is the number of points in the fibre with $U \neq 0$, and the fourth the contribution to

$$
\chi\left(\left\{(W, U) \mid 4 U^{3}=\left(3 a_{1} U+a_{3}\right)^{2}, U \neq 0\right\} .\right.
$$

| $P^{m-1} \backslash\left(V_{3} \cup V_{3}^{*}\right)$ | $m-\chi\left(V_{3}\right)-\chi\left(V_{3}^{*}\right)+\chi\left(V_{4}\right)$ | 3 | $3 m-3 \chi\left(V_{3}\right)-3 \chi\left(V_{3}^{*}\right)+3 \chi\left(V_{4}\right)$ |
| :--- | :--- | :--- | :--- |
| $V_{3} \backslash V_{3}^{*}$ | $\chi\left(V_{3}\right)-\chi\left(V_{4}\right)$ | 1 | $\chi\left(V_{3}\right)-\chi\left(V_{4}\right)$ |
| $V_{3}^{*} \backslash V_{3}$ | $\chi\left(V_{3}^{*}\right)-\chi\left(V_{4}\right)$ | 2 | $2 \chi\left(V_{3}^{*}\right)-2 \chi\left(V_{4}\right)$ |
| $V_{4}$ | $\chi\left(V_{4}\right)$ | 0 | 0 |

Hence $\chi\left(V \cap \mathcal{S}_{2}\right)=3 m-2 \chi\left(V_{3}\right)-\chi\left(V_{3}^{*}\right)$.
In $\mathcal{S}_{1}(X=0, B \neq 0)$ we can again assign $W, B$ and $U$ and solve for $E$. Since the other conditions are independent of $B$, which runs through $\mathbb{C}^{*}$, we have $\chi=0$ in each case, except when $W=U=0$ which leads to the unique point with $E=0$ also, and hence to $\chi\left(V \cap \mathcal{S}_{1}\right)=1$.

Finally, for $\mathcal{S}_{0}$, while we again solve uniquely for $E$, so that $\chi\left(\phi^{-1}(0) \cap \overline{\mathcal{S}}_{0}\right)=1$, we have to distinguish according as $B^{3}+X^{3} E-3 X^{2} B U=0$ or not, hence according as $0=B^{3}-3 X^{2} B U-X^{3}\left(3 a_{1} U+a_{3}\right)$. As before, since here $X$ is non-zero, we can replace $B$ by $Z=B / X$, so obtain $0=Z^{3}-3 Z U-\left(3 a_{1} U+a_{3}\right)$, giving $V_{0}$. By Lemma 7.3, both $V_{0}$ and its intersection $V_{1}$ with $X=0$ are non-singular. Applying again Lemma 7.1, we obtain, since $\chi\left(\overline{\mathcal{S}}_{0}\right)=1$,

$$
\chi\left(V \cap \mathcal{S}_{0}\right)=\left(\chi\left(V_{0}\right)-\chi\left(V_{1}\right)\right)+2\left(1-\chi\left(V_{0}\right)+\chi\left(V_{1}\right)\right)=2-\chi\left(V_{0}\right)+\chi\left(V_{1}\right)
$$

Recall that by Corollary 4.4, if $H^{n}$ is a smooth hypersurface of dimension $n$ where one of the weights is 2 , then $\chi\left(H^{n}\right)=n+2+\frac{1}{3}\left\{(-2)^{n}-1\right\}$. Since $V_{0}, V_{1}$ have respective dimensions $m+1, m$,

$$
\chi\left(V_{0}\right)-\chi\left(V_{1}\right)=1+\frac{1}{3}\left\{(-2)^{m+1}-(-2)^{m}\right\}=1-(-2)^{m}
$$

Adding up, $\chi(V)$ is equal to

$$
\chi\left(V_{3}\right)+\chi(\mathcal{E})+3 m-2 \chi\left(V_{3}\right)-\chi\left(V_{3}^{*}\right)+1+2-\left(1-(-2)^{m}\right),
$$

and to $\chi_{m+3}(3)+(-1)^{m} \mu(V)$. We can substitute $\chi(\mathcal{E})=2, \chi\left(V_{3}^{*}\right)=\chi_{m-2}(3)$ and $\chi\left(V_{3}\right)=\chi_{m-2}(3)+(-1)^{m-1} \mu\left(V_{3}\right)$, so that

$$
\mu(V)-\mu\left(V_{3}\right)=(-1)^{m-1}\left(\chi_{m+3}(3)+2 \chi_{m-2}(3)-3 m-5\right)+\left(2^{m}-(-1)^{m}\right)
$$

Substituting $\chi_{n}(3)=\frac{(-1)^{n}}{3}\left(2^{n+2}-(-1)^{n+2}\right)+n+2$, this reduces to $11.2^{m}$.
As before, we can determine the singularities.

Proposition 7.5. Suppose $V$ in Case 4. Then the singularity of $V$ corresponding to $a$ singularity of $V_{3}$ is right-equivalent to a suspension of that singularity.

Proof. Suppose $Q_{i}$ a singular point of $V_{3}$. Then $a_{1}$ cannot be a multiple of $x_{0}$, for otherwise $V_{3}=V_{3}^{*}$ would be non-singular. We may thus set $w_{1}=a_{1}$. By the arguments above, $a_{1}$ does not vanish at $Q_{i}$. We may thus work in affine co-ordinates $w_{1}=1$.

Now apply Lemma 3.4(a) to $f$ with the variables $\left(x_{0}, x_{1}, x_{2}, x_{3}, x_{4}\right)$. We observe that $\partial f / \partial x_{0}, \partial f / \partial x_{1}, \partial f / \partial x_{2}, \partial f / \partial x_{3}, \partial f / \partial x_{4}$ all vanish when $x_{0}=x_{1}=x_{2}=$ $x_{3}=0$ and $x_{4}=-\frac{1}{2} \partial a_{3} / \partial x_{0}$. Substituting these values gives $g(w)=a_{3}(w, 0,0,0)$. The result follows.

Lemma 7.6. Suppose $V$ in Case 4. Then the singularity of $V$ at $P$ is semi-quasihomogeneous of degree 12 in variables of weights $1,6,6,6$ and 4 ( $m$ times).

Proof. If we substitute $x_{4}=1, x_{3}=2 z, x_{2}=y+3 z^{2}$, and $x_{1}=x+3 y z+3 z^{3}$, we obtain $E=6 U y-8 y^{3}-6 x^{2}$; thus $f=-6 x^{2}+6 U y-8 y^{3}+3 U a_{1}+a_{3}$.

We also obtain $U=2 x_{0}-4 x z-6 y z^{2}-3 z^{4}+y^{2}$. Substitute $x_{0}=\frac{1}{2} U+2 x z-$ $\frac{1}{2} y^{2}+3 y z^{2}+\frac{3}{2} z^{4}$ in $f$, and assign weights 1 to $z, 4$ to $y$ and to the $w_{i}, 6$ to $x$ and 8 to $U$. Then in the expression for $x_{0}$, all terms have weight $>4$ except for $\frac{3}{2} z^{4}$, of weight 4. Hence all terms in $f$ have weight at least 12, and those of exactly this weight are obtained by substituting $\frac{3}{2} z^{4}$ for $x_{0}$ in $a_{1}$ and $a_{3}$.

It thus remains only to show that the result of this substitution has an isolated singularity. Here we can ignore the summand $-6 x^{2}$; the rest is obtained from $\psi_{1}(W, X, Z, U)$ by the substitution $Z=2 y, X=\frac{3}{2} z^{4}$. But by Lemma 7.3, the hypersurface $V_{0}$ given by $\psi_{1}=0$ is non-singular. The result follows.

We observe that we also have $f=-6 x^{2}+\left(U-\frac{4}{3} y^{2}+\frac{2}{3} a_{1} y-a_{1}^{2}\right)\left(2 y+a_{1}\right)+a_{1}^{3}+a_{3}$, which we can write as $-6 x^{2}+6 U^{\prime} y^{\prime}+a_{3}^{*}\left(w, x_{0}\right)$, though in view of the substitution $x_{0}=\frac{1}{2} U^{\prime}+2 z x+3 z^{2} y^{\prime}-\frac{3}{2} z^{2} a_{1}+\frac{3}{2} z^{4}+\frac{1}{6} y^{\prime 2}-\frac{1}{2} y^{\prime} a_{1}+\frac{3}{8} a_{1}^{2}$ we must make for $x_{0}$, the simplicity of this form is misleading.

In certain cases, we can also determine $\tau$.
Proposition 7.7. Suppose $f$, in the normal form for Case 4, satisfies also
(i) $a_{1}$ is a multiple of $x_{0}$, and
(ii) $\partial a_{3} / \partial x_{0}$ vanishes when $x_{0}=0$.

Then the singularity of $f$ at $P$ is quasi-homogeneous, so $\tau_{P}(V)=\mu_{P}(V)=11.2^{m}$.
Proof. Write $f$ as $f=E+a x_{0} U+u x_{0}^{2}+c x_{0}^{3}+C(w)$, where $u$ is a non-zero linear combination of the $w_{i}$. We now define a number of vector fields. In the table, the left column gives the name, the next defines the field, and the last gives its effect on $f$. Here $\partial_{i}$ denotes $\partial / \partial x_{i}$ and $R=\left(\partial_{4} E \partial_{0}-\partial_{3} E \partial_{1}+\partial_{2} E \partial_{2}-\partial_{1} E \partial_{3}+\partial_{0} E \partial_{4}\right) / 6$.

| $H$ | $-2 x_{0} \partial_{0}-x_{1} \partial_{1}+x_{3} \partial_{3}+2 x_{4} \partial_{4}$ | $-2 a x_{0} U-2 x_{0}^{2}\left(2 u+3 c x_{0}\right)$ |
| :--- | :--- | :--- |
| $M$ | $x_{1} \partial_{1}+(3 / 2) x_{2} \partial_{2}+(3 / 2) x_{3} \partial_{3}+x_{4} \partial_{4}$ | $a x_{1} U+2 x_{0} x_{1}\left(2 u+3 c x_{0}\right)$ |
| $P$ | $x_{0} \partial_{0}+x_{1} \partial_{1}+x_{2} \partial_{2}+x_{3} \partial_{3}+x_{4} \partial_{4}$ | $3 E+3 a x_{0} U+x_{0}^{2}\left(2 u+3 c x_{0}\right)$ |
| $Q$ | $x_{0} \partial_{2}+x_{1} \partial_{3}+x_{2} \partial_{4}$ | $2 x_{0}(3 U+a B)$ |
| $R$ |  | $6 U^{2}+a\left(x_{0} E+B U\right)+x_{0} B\left(2 u+3 c x_{0}\right)$ |
| $S$ | $\left(\partial_{4} E \partial_{2}-\partial_{3} E \partial_{3}+\partial_{2} E \partial_{4}\right) / 12$ | $3 x_{0} E-3 B U+a x_{0}^{2} U$ |

Linear combinations of these give

| $Y_{1}$ | $2 x_{4} H+x_{3} M$ | $-\left(4 x_{0} x_{4}-x_{1} x_{3}\right)\left(a U+x_{0}\left(2 u+3 c x_{0}\right)\right)$ |
| :--- | :--- | :--- |
| $Y_{2}$ | $-x_{1} M+U \partial_{2}+R$ | $a x_{0} E+6\left(4 x_{0} x_{4}-x_{1} x_{3}\right) U$ |
|  |  | $+2 a B U+2 x_{0}\left(x_{0} x_{2}-x_{1}^{2}\right)\left(2 u+3 c x_{0}\right)$ |
| $Y_{3}$ | $-5 S+4 x_{0} P-\frac{7}{2} U \partial_{4}$ | $-3 x_{0} E-6 B U+4 x_{0}^{3}\left(2 u+3 c x_{0}\right)$ |
| $Y_{4}$ | $3 Q+3 x_{0} \partial_{2}-2 a x_{0} \partial_{4}$ | $18 x_{0}\left(4 x_{0} x_{4}-x_{1} x_{3}\right)-6 a x_{0}\left(x_{0} x_{2}-x_{1}^{2}\right)-4 a^{2} x_{0}^{3}$ |

Thus the vector field $Z=Y_{1}+\frac{a}{6} Y_{2}+\frac{a^{2}}{18} Y_{3}+\frac{2 b u+3 c x_{0}}{18} Y_{4}$ kills $f$, and at the point $P$, $Z$ reduces to $2 \partial / \partial x_{4}$. The result now follows by Saito's criterion.

Condition (i) is invariant under allowed changes of co-ordinates. An invariant version of (ii) is that substituting $x_{0}=0$ in $\partial a_{3} / \partial x_{0}$ gives a function in the Jacobian ideal of $a_{3}(w, 0)$. We believe both these conditions to be necessary for the result.

Since $\tau(V) \leq \mu(V)=11.2^{m}<12.2^{m}$, no function in Case 4 can be oversymmetric.

## 8 Case 21

First we normalise co-ordinates.
Lemma 8.1. There is an allowable change of co-ordinates which puts $f$ in the form $f=S+a_{1} B+a_{3}$, where $a_{i}$ is homogeneous of degree $i$ in $w, x_{0}, y_{0}$. Moreover, we may suppose that either $a_{1}=x_{0}$ or $a_{1}=w_{1}$.

Proof. Here $d=3$ and $f$ has the form $a_{0} S+a_{1} B+b_{1} T+a_{3}$, where $a_{i}$ (and $b_{i}$ ) denotes a homogeneous function of degree $i$ in $w, x_{0}, y_{0}$. It follows from Corollary 2.5 that $f$ must contain the monomial $y_{1}^{2} x_{0}$. Hence we must have $a_{0} \neq 0$, and can take $a_{0}=1$. Now substitute $y_{1}^{\prime}:=y_{1}+\frac{1}{2} b_{1}$ to reduce $b_{1}$ to 0 (the extra terms introduced can be absorbed in $a_{3}$ ), and so $f$ to $S+a_{1} B+a_{3}$.

If $a_{1}$ involves any of the $w$ co-ordinates, we can make a linear substitution among the $w$ 's to reduce $a_{1}$ to the form $q x_{0}+p y_{0}+w_{1}$, and then change again to achieve $a_{1} \equiv w_{1}$. Otherwise, we can write $a_{1}=2 p y_{0}+q x_{0}$ and use the substitution $y_{1}^{\prime}=$ $y_{1}-p x_{1}, y_{0}^{\prime}=y_{0}-p x_{0}$. This transforms $S$ to $S+\left(p^{2} x_{0}-2 p y_{0}\right) B$, so $a_{1}$ is changed to $\left(q-p^{2}\right) x_{0}$. We may thus suppose that either $a_{1} \equiv x_{0}$ or $a_{1} \equiv 0$. However, if $a_{1} \equiv 0, V$ is singular along the plane $w=y_{0}=y_{1}=0$.

Proposition 8.2. In Case 21, $V$ is always 2-symmetric.
Proof. The function $f$ is annihilated by the vector fields $\xi=x_{0} \partial / \partial x_{1}+x_{1} \partial / \partial x_{2}+$ $y_{0} \partial / \partial y_{1}$ and $\xi^{\prime}:=y_{0} \partial / \partial x_{1}+y_{1} \partial / \partial x_{2}-a_{1} \partial / \partial y_{1}$.

There are significant differences between the two cases. We now rename the case with $a_{1}=x_{0}$ as $21_{0}$, and the case $a_{1}=w_{1}$ as Case 5 .

In Case $21_{0}$, the operators $D \pm D^{\prime}$ each fall into Case 2 (other linear combinations are all in Case $21_{0}$ ). To see this, substitute

$$
x_{0}=u_{0}+v_{0}, \quad y_{0}=u_{0}-v_{0}, \quad x_{1}=u_{1}+v_{1}, \quad y_{1}=u_{1}-v_{1},
$$

Then $f$ reduces to $-8 x_{2} u_{0} v_{0}+4 u_{0} v_{1}^{2}+4 v_{0} u_{1}^{2}+a_{3}$, which we can write as $4 u_{0}\left(v_{1}^{2}-2 v_{0} x_{2}\right)$ added to a cubic in $w, u_{0}, v_{0}$ and $v_{1}$. Conversely, this is essentially the normal form for Case 2 where the cubic has the form $4 v_{0} u_{1}^{2}$ added to a cubic in $w, u_{0}, v_{0}$, and thus can be identified with Case (a2) of Proposition 5.9.

For the remainder of this section we consider only Case 5; here all non-zero linear combinations of $D$ and $D^{\prime}$ are in Case 5 . We re-name $a_{1}$ as $z_{0}$. In this notation, co-ordinates are $\left(w_{1}, \ldots, w_{m}, x_{0}, y_{0}, z_{0}, x_{1}, y_{1}, x_{2}\right)$, and we re-group the terms in $f$ as $J+a_{3}\left(x_{0}, y_{0}, z_{0}, w\right)$, where

$$
J:=2 x_{2}\left(y_{0}^{2}-x_{0} z_{0}\right)+x_{1}^{2} z_{0}-2 x_{1} y_{1} y_{0}+y_{1}^{2} x_{0}
$$

It is now more natural to treat all the differential operators on the same footing, and consider $f$ as invariant under the 2-dimensional group $G$ whose Lie algebra is spanned by $D$ and $D^{\prime}$. We see that $J$ is invariant under $G$, and it seems very likely that the ring of invariants coincides with the polynomial ring $\mathbb{C}\left[w, x_{0}, y_{0}, z_{0}, J\right]$ (clearly it contains this, and we can show that the localisation at $\left\langle x_{0}\right\rangle$ is correct), but we will not use the precise assertion.

We define a new stratification,
$\mathcal{S}_{0}: y_{0}^{2}-x_{0} z_{0} \neq 0$,
$\mathcal{S}_{1}: y_{0}^{2}-x_{0} z_{0}=0$ but $\left(x_{0}, y_{0}, z_{0}\right) \neq(0,0,0)$,
$\mathcal{S}_{2}:\left(x_{0}, y_{0}, z_{0}\right)=(0,0,0)$.
Lemma 8.3. A point in $\mathcal{S}_{2}$ is fixed under $G$; otherwise the dimension of the orbit is equal to the rank of $\left(\begin{array}{ccc}x_{0} & y_{0} & x_{1} \\ y_{0} & z_{0} & y_{1}\end{array}\right)$.

For the fixed points of $A \xi_{0}+B \xi_{1}=\left(A x_{0}+B y_{0}\right) \partial / \partial x_{1}+\left(A y_{0}+B z_{0}\right) \partial / \partial y_{1}+$ $\left(A x_{1}+B y_{1}\right) \partial / \partial x_{2}$ are those where the three coefficients vanish; this holds for some $(A, B) \neq(0,0)$ if and only if the rank of the matrix drops.

We define a projection $\pi: K \rightarrow L$ by $\pi\left(w, x_{0}, y_{0}, z_{0}, x_{1}, y_{1}, x_{2}\right)=\left(w, x_{0}, y_{0}, z_{0}\right)$, and again write $\bar{\pi}: P(K) \rightarrow P(L)$ for the induced map of projective spaces. The exceptional set where $\bar{\pi}$ is undefined is the projective plane where $w=x_{0}=y_{0}=$ $z_{0}=0$, and the pre-image of any point in the target is isomorphic to affine 3 -space. Write $\bar{\pi}_{0}$ for the restriction of $\bar{\pi}$ to the hypersurface $V$ defined by $f=0$.

Lemma 8.4. The fibres of $\bar{\pi}_{0}$ are as follows:
in $\overline{\mathcal{S}}_{0}$ each fibre is a quadric isomorphic to an affine plane,
in $\overline{\mathcal{S}}_{1}$, one or two affine planes according as $a_{3}\left(w, x_{0}, y_{0}, z_{0}\right)=0$ or $\neq 0$,
in $\overline{\mathcal{S}}_{2}$, affine 3 -space or the empty set according as $a_{3}(w, 0,0,0)=0$ or $\neq 0$.
Proof. The assertions are trivial except for $\overline{\mathcal{S}}_{1}$. Here we may write $\left(x_{0}, y_{0}, z_{0}\right)=$ $\left(t^{2}, t u, u^{2}\right)$ for some $t, u$ not both zero. Then $J$ reduces to $x_{1}^{2} z_{0}-2 x_{1} y_{1} y_{0}+y_{1}^{2} x_{0}=$ $\left(u x_{1}-t y_{1}\right)^{2}$. Thus we have the plane(s) given by $u x_{1}-t y_{1}= \pm \sqrt{ }\left(-a_{3}\right)$.

We have two hypersurfaces in $P(L)$ : the cone $\overline{\mathcal{S}}_{1} \cup \overline{\mathcal{S}}_{2}=C: y_{0}^{2}-x_{0} z_{0}=$ 0 and the variety $V_{0}$ defined by $a_{3}\left(w, x_{0}, y_{0}, z_{0}\right)=0$. Write $V_{1}:=V_{0} \cap C$, and $V_{3}:=V_{0} \cap \overline{\mathcal{S}}_{2}$ for the variety $a_{3}(w, 0,0,0)=0$. Also define $V_{1}^{*}$ as the variety $\phi(w, t, u):=a_{3}\left(w, t^{2}, t u, u^{2}\right)=0$ in weighted projective space $P\left(2^{m} 1^{2}\right)$.

It follows from the lemma that

$$
\begin{equation*}
\chi(V)=3+\chi\left(\overline{\mathcal{S}}_{0}\right)+2 \chi\left(\overline{\mathcal{S}}_{1}\right)-\chi\left(\overline{\mathcal{S}}_{1} \cap V_{0}\right)+\chi\left(\overline{\mathcal{S}}_{2} \cap V_{0}\right) . \tag{4}
\end{equation*}
$$

Proposition 8.5. The variety $V$ has isolated singularities if and only if
(i) for no $\left(w, x_{0}, y_{0}, z_{0}\right)$ in $\mathcal{S}_{1}$ do we have $\partial a_{3} / \partial w_{i}=0$ for each $i$, and the matrix A of rank 1, where

$$
A:=\left(\begin{array}{ccc}
\partial a_{3} / \partial x_{0} & \partial a_{3} / \partial y_{0} & \partial a_{3} / \partial z_{0} \\
z_{0} & -2 y_{0} & x_{0}
\end{array}\right)
$$

(ii) for any singular point of $V_{3}$ we have

$$
\left(\partial a_{3} / \partial y_{0}\right)^{2} \neq 4\left(\partial a_{3} / \partial x_{0}\right)\left(\partial a_{3} / \partial z_{0}\right) .
$$

In particular, singular points of $V_{3}$ are isolated.
When this holds, $P$ is the only singular point of $V$.
Proof. Since $\partial f / \partial x_{2}=2\left(y_{0}^{2}-x_{0} z_{0}\right)$, there are no singularities in $\mathcal{S}_{0}$.
In $\mathcal{S}_{1}$, again write $\left(x_{0}, y_{0}, z_{0}\right)=\left(t^{2}, t u, u^{2}\right)$ for $t, u$ not both zero. Then

$$
\partial f / \partial x_{1}=2 u\left(u x_{1}-t y_{1}\right), \quad \partial f / \partial y_{1}=-2 t\left(u x_{1}-t y_{1}\right),
$$

$\partial f / \partial x_{0}=y_{1}^{2}+\partial a_{3} / \partial x_{0}, \partial f / \partial y_{0}=-2 x_{1} y_{1}+\partial a_{3} / \partial y_{0}, \partial f / \partial z_{0}=x_{1}^{2}+\partial a_{3} / \partial z_{0}$, and $\partial f / \partial w_{i}=\partial a_{3} / \partial w_{i}$. It follows that for a critical point of $f, u x_{1}=t y_{1}$, and hence that the matrix A has rank 1 .

Conversely, given $\left(w, x_{0}, y_{0}, z_{0}\right)$ in $\mathcal{S}_{1}$ such that $\partial a_{3} / \partial w_{i}=0$ for each $i$, and $A$ has rank 1 , we can take $\left(x_{0}, y_{0}, z_{0}\right)=\left(t^{2}, t u, u^{2}\right)$, and let the upper row of $A$ equal $-v^{2}$ times the lower. Then if $x_{1}=t v, y_{1}=u v$ and $x_{2}$ is arbitrary, we have a critical point of $f$ : none of these critical points is isolated.

In $\mathcal{S}_{2}$ we have identically $\partial f / \partial x_{2}=\partial f / \partial x_{1}=\partial f / \partial y_{1}=0$. For a critical point of $f$ we have a critical point of $a_{3}$, and 3 further equations, for $x_{1}^{2}, x_{1} y_{1}$ and $y_{1}^{2}$, which are inconsistent unless also $\left(\partial a_{3} / \partial y_{0}\right)^{2}=4\left(\partial a_{3} / \partial x_{0}\right)\left(\partial a_{3} / \partial z_{0}\right)$.

Conversely, if there is a critical point of $a_{3}$ at which this identity holds, we can solve for $x_{1}$ and $y_{1}$ and take an arbitrary value for $x_{2}$, again obtaining a non-isolated singularity of $f$.

It remains to consider the case $w=x_{0}=y_{0}=z_{0}=0$. Here the only critical point is $x_{1}=y_{1}=0$, which is indeed isolated.

For a singular point of the intersection $V_{1}:=V_{0} \cap C$, Lagrange's multiplier rule tells us that the $\partial a_{3} / \partial w_{i}$ vanish and the matrix $A$ has rank at most 1: for singular points with $\left(x_{0}, y_{0}, z_{0}\right) \neq(0,0,0)$, this condition is necessary and sufficient. Thus (i) is equivalent to the condition that $V_{1} \cap \overline{\mathcal{S}}_{1}$, or equivalently the open set of $V_{1}^{*}$ where $(t, u) \neq(0,0)$, be non-singular.

Now $V_{1}$ is always singular along $V_{1} \cap \overline{\mathcal{S}}_{2}$. For $V_{1}^{*}$ on the other hand, it is singular at a point on $t=u=0$ only if $\partial a_{3} / \partial w_{i}=0$ for each $i$, i.e. at a singular point of $V_{3}$. Thus $V_{1}^{*}$ is non-singular if and only if $V_{3}$ is.

The singularities of $V_{1}^{*}$ and $V_{3}$ are related as follows. If we expand $\phi$ as a Taylor series, the second order terms in $t$ and $u$ are $2\left(\left(\partial a_{3} / \partial x_{0}\right) t^{2}+\left(\partial a_{3} / \partial y_{0}\right) t u+\right.$ $\left.\left(\partial a_{3} / \partial z_{0}\right) u^{2}\right)$, a form which is non-singular if and only if (ii) holds. However, since the ambient weighted projective space is singular at this point, we cannot say that one singularity is the suspension of the other.

Theorem 8.6. If $V_{3}$ is non-singular, $\mu(V)=25.2^{m}$. Moreover, the singularity of $V$ is semi-quasi-homogeneous of degree 6 in variables of weights $1,1,3,3,3$ and 2 ( $m$ times).

Proof. Although the first assertion follows from the second, we give an independent proof.

For the stratification of $P(L)$ we have $\chi\left(\overline{\mathcal{S}}_{2}\right)=m, \chi\left(\overline{\mathcal{S}}_{1}\right)=2$ and $\chi\left(\overline{\mathcal{S}}_{0}\right)=1$, since $\overline{\mathcal{S}}_{2}$ is a projective space, and forgetting the $w_{i}$ defines a projection of the others to the projective plane on $x_{0}, y_{0}, z_{0}$ with contractible fibres. Substituting in (4), and using the notations $V_{i}$ thus gives $\chi(V)=8-\chi\left(V_{1}\right)+2 \chi\left(V_{3}\right)$.

Now $\chi(V)=\chi_{m+4}(3)+(-1)^{m-1} \mu(V)$ and since $V_{3}$ is non-singular, $\chi\left(V_{3}\right)=$ $\chi_{m-2}(3)$. Since the natural projection $V_{1}^{*} \rightarrow V_{1}$ is bijective, $\chi\left(V_{1}\right)=\chi\left(V_{1}^{*}\right)$, while by Corollary 3.7, $\chi\left(V_{1}^{*}\right)=m+1+\frac{1}{6}(-1)^{m}\left(26.2^{m}+4(-1)^{m}\right)$.

Putting the above results together, we obtain

$$
\begin{aligned}
\mu(V)= & \frac{1}{3}\left(2^{m+6}-(-1)^{m}\right)+(-1)^{m}(m+6) \\
& +\frac{1}{6}\left(26.2^{m}+4(-1)^{m}\right)+(-1)^{m}(m+1) \\
& -\frac{2}{3}\left((2)^{m}-(-1)^{m}\right)+(-1)^{m-1} 2 m+8(-1)^{m-1} \\
= & 25.2^{m} .
\end{aligned}
$$

Set $x_{2}=1 / 2$ and rewrite $J$ as $\left(y_{0}-x_{1} y_{1}\right)^{2}-\left(x_{0}-x_{1}^{2}\right)\left(z_{0}-y_{1}^{2}\right)$. Substitute $x^{\prime}:=x_{0}-x_{1}^{2}$, $y^{\prime}:=y_{0}-x_{1} y_{1}$ and $z^{\prime}:=z_{0}-y_{1}^{2}$ : then $f=\left(y^{\prime 2}-x^{\prime} z^{\prime}\right)+a_{3}\left(x^{\prime}+x_{1}^{2}, y^{\prime}+x_{1} y_{1}, z^{\prime}+y_{1}^{2}, w\right)$. Now assign weights 1 to $x_{1}, y_{1}, 2$ to the $w_{i}$ and 3 to $x^{\prime}, y^{\prime}, z^{\prime}$. Then the terms of weight 6 give $g=\left(y^{\prime 2}-x^{\prime} z^{\prime}\right)+a_{3}\left(x_{1}^{2}, x_{1} y_{1}, y_{1}^{2}, w\right)$. Since $V_{3}$ is non-singular, so is $V_{1}^{*}$, so $g$ has an isolated singularity, and the result follows.

If $V_{1}^{*}$ is singular, we do not have a formula for its Euler characteristic, so must proceed differently: in fact, we resolve the singularity of $C$. Define $\widehat{P}$ as the subvariety of $P^{1} \times P(L)$, where $P^{1}$ has co-ordinates $\left(t_{0}: t_{1}\right)$, given by $t_{1} x_{0}=t_{0} y_{0}, t_{1} y_{0}=$ $t_{0} z_{0}$. If $L_{0}:=\mathcal{S}_{2}$ denotes the subspace $x_{0}=y_{0}=z_{0}=0$ of $L$, then the projection $\widehat{P} \rightarrow P(L)$ has image $C$; it is bijective over $C \backslash P\left(L_{0}\right)$, but over $P\left(L_{0}\right)$ is the projection $P^{1} \times P\left(L_{0}\right) \rightarrow P\left(L_{0}\right)$.

Define $\widehat{V} \subset \widehat{P}$ to be the subvariety given by $a_{3}\left(w, x_{0}, y_{0}, z_{0}\right)=0$ : thus it is a complete intersection of multi-degree $(1,1),(1,1),(0,3)$. The natural projection $\pi: \widehat{V} \rightarrow V_{1}$ is an isomorphism outside $V_{3}$, but a product $P^{1} \times V_{3}$ over it. In particular, $\chi(\widehat{V})=\chi\left(V_{1}\right)+\chi\left(V_{3}\right)$. By Proposition $8.5(\mathrm{i}), V_{1} \backslash V_{3}$, hence its preimage, is non-singular. Also, any singularity of $\widehat{V}$ projects to a singular point $P_{i}$ of $V_{3}$.

Lemma 8.7. Above each singular point $P_{i}$ of $V_{3}$ there are just two singular points of $\widehat{V}$, and the singularity at each is isomorphic to a suspension of the singularity of $V_{3}$ at $P_{i}$.

Proof. To study a neighbourhood of $P_{i}$, it is convenient to make a linear change of the co-ordinates $\left(x_{0}, y_{0}, z_{0}\right)$, preserving the quadratic $y_{0}^{2}-x_{0} z_{0}$, so that $\partial a_{3} / \partial x_{0}=$ $\partial a_{3} / \partial z_{0}=0, \partial a_{3} / \partial y_{0} \neq 0$ at $P_{i}$ (here we use Proposition 8.5(ii)). Then $\widehat{V}$ is nonsingular at all points of $\pi^{-1}\left(P_{i}\right)$ except those where $\left(t_{0}: t_{1}\right)$ is $(0: 1)$ or $(1: 0)$ : it suffices to consider the first.

Take affine co-ordinates in $\widehat{V}$ with $t_{1}=1$ and $w_{1}=1$. Then $y_{0}=t_{0} z_{0}, x_{0}=$ $t_{0} y_{0}=t_{0}^{2} z_{0}$. Thus $a_{3}$ lifts to $\alpha_{3}:=a_{3}\left(1, w_{2}, \ldots, w_{r}, t_{0}^{2} z_{0}, t_{0} z_{0}, z_{0}\right)$. At the point $t_{0}=z_{0}=0$, the 2 -jet has a non-zero coefficient of $t_{0} z_{0}$. We now apply Lemma 3.4(b): this we can do since $\partial \alpha_{3} / \partial t_{0}$ is divisible by $z_{0}$. It thus follows that we have a suspension of the restriction to $t_{0}=z_{0}=0$, which is just the intersection with $V_{3}$.

We can now show
Theorem 8.8. If $V$ is quasi-smooth in Case 5, then $\mu(V)=25.2^{m}+\mu\left(V_{3}\right)$.
Proof. Since $\widehat{V}$ is a complete intersection with isolated singularities, its Euler characteristic is obtained from that of a smooth complete intersection of the same multidegrees by adding $(-1)^{m-1}$ times the sum of the Milnor numbers.

While we could calculate the default value of $\chi(\widehat{V})$ directly, we can also obtain it from the above calculations in the case when $V_{3}$ is non-singular. In the proof of Theorem 8.6 we calculated values of $\chi\left(V_{1}\right)$ and $\chi\left(V_{3}\right)$ : denote them for now by $c_{1}$ and $c_{3}$. Thus $c_{1}=m+1+\frac{1}{6}(-1)^{m}\left(26.2^{m}+4(-1)^{m}\right)$ and $c_{3}==\chi_{m-2}(3)$. In the case $V_{3}$ non-singular, we have $\chi(\widehat{V})=c_{1}+c_{3}$ and $\chi(V)=8-c_{1}+2 c_{3}$, whereas $\chi(V)=\chi_{m+4}(3)+(-1)^{m-1} \mu(V)$, leading to $\mu(V)=25.2^{m}$.

In the general case, by Lemma 8.7 there are two singular points of $\widehat{V}$ in $\pi^{-1}\left(P_{i}\right)$, each a suspension of the singularity of $V_{3}$ at $P_{i}$; hence $\mu(\widehat{V})=2 \mu\left(V_{3}\right)$. Now $\chi\left(V_{3}\right)=$ $c_{3}+(-1)^{m-1} \mu\left(V_{3}\right)$. Also, by the remark just made,
$\chi(\widehat{V})=c_{1}+c_{3}+(-1)^{m-1} 2 \mu\left(V_{3}\right)$. As before, we have $\chi(V)=8-\chi\left(V_{1}\right)+$ $2 \chi\left(V_{3}\right)=8-\chi(\widehat{V})+3 \chi\left(V_{3}\right)$, which now equals $8-\left[c_{1}+c_{3}+(-1)^{m-1} 2 \mu\left(V_{3}\right)\right]+$ $3\left[c_{3}+(-1)^{m-1} \mu\left(V_{3}\right)\right]$, i.e. $8-c_{1}+2 c_{3}+(-1)^{m-1} \mu\left(V_{3}\right)$. Substituting this value in $\chi(V)=\chi_{m+4}(3)+(-1)^{m-1} \mu(V)$ gives $\mu(V)=25.2^{m}+\mu\left(V_{3}\right)$ as desired.

Observe, however, that unlike the other cases, here there is just one singular point $P$, and the values $\tau=25.2^{m}$ and $\mu=25.2^{m}+\mu\left(V_{3}\right)$ both hold for the singularity at this point.

## $9 \quad 2$-symmetric cases

Finally we list 2 -symmetric hypersurfaces. This is just the subcase $d=3$ of the list of oversymmetric cases: in the semisimple case, the weights are obtained from one of $[-1,0,1],[-2,1,2]$ and $[-2,1,4]$ by adding zeros; in the unipotent case, the subcases of Case 2 when an additional action exists were analysed in Proposition 5.9: we had three cases (a1), (a2), (b); and Case 21 splits into two subcases: Case $21_{0}$ and Case 5. Since some cases arise more than once by using different 1-parameter subgroups, it is better to give the list separately.

Theorem 9.1. If $f$, of degree $\geq 3$, such that $f=0$ is quasi-smooth, is 2-symmetric, then $f$ belongs to one of the following 5 cases (A)-(E).
(A) $f=x_{0} x_{1} x_{2}+a_{3}\left(x_{3}, \ldots, x_{n}\right)(n \geq 2)$, where $a_{3}=0$ is non-singular, with the 2-parameter action $(\lambda, \mu) \cdot\left(x_{0}, x_{1}, x_{2}, \ldots, x_{n}\right)=\left(\lambda^{-1} x_{0}, \mu^{-1} x_{1}, \lambda \mu x_{2}, \ldots, x_{n}\right)$, and 3 singular points, mutually isomorphic. We have 1 -parameter subgroups with weights
$[-a,-b, a+b]$ for any $a$ and $b$. There are 3 singularities, all isomorphic, each with $\mu=2^{n-2}$ and homogeneous of degree 3 with respect to weights $1(n-1$ times) and 2 .
(B) $f=x_{0} x_{1}^{2}+x_{0} x_{2} x_{3}+a_{3}\left(x_{3}, \ldots, x_{n}\right)(n \geq 3)$, where $a_{3}=0$ is non-singular, and has non-singular intersection with $x_{3}=0$. This is annihilated by $-2 x_{0} \partial_{0}+$ $x_{1} \partial_{1}+2 x_{2} \partial_{2}$ with non-zero weights $[-2,1,2]$, and by $x_{3} \partial / \partial x_{1}-2 x_{1} \partial / \partial x_{2}$, which is in case (a1) of Proposition 5.9. There are two singularities, with Milnor numbers $5.2^{n-3}$ and $2^{n-3}$; both homogeneous of degree 6 , the first with respect to weights 1 and 2 ( $n-3$ times); the second with respect to weights 3 and $2(n-3$ times).
(C) $f=x_{0}\left(2 x_{0} x_{2}-x_{1}^{2}\right)+a_{3}\left(x_{3}, \ldots, x_{n}\right)(n \geq 2)$, where $a_{3}=0$ is non-singular. This is annihilated by $-2 x_{0} \partial_{0}+x_{1} \partial_{1}+4 x_{2} \partial_{2}$ with non-zero weights $[-2,1,4]$, and by $x_{0} \partial / \partial x_{1}+x_{1} \partial / \partial x_{2}$, which is in case (b) of Proposition 5.9. The singularity has $\mu=3.2^{n-2}$ and is homogeneous of degree 12 with respect to weights 3,6 and 4 ( $n-2$ times).
(D) $f=x_{3}\left(2 x_{0} x_{2}-x_{1}^{2}\right)+x_{0} x_{4}^{2}+a_{3}\left(x_{0}, x_{3}, x_{5}, \ldots, x_{n}\right)(n \geq 4)$, with $a_{3}$ nonsingular. We have vector fields $x_{0} \partial / \partial x_{1}+x_{1} \partial / \partial x_{2}$ in case (a2) of Proposition 5.9 and $x_{4} \partial / \partial x_{2}-x_{3} \partial / \partial x_{4}$ in Case $21_{0}$. The singularity has $\mu=13.2^{n-4}$, and is in the same $\mu$-constant stratum as $x^{6}+x^{2} y^{2}+y^{6}+\sum_{2}^{n-4} w_{i}^{3}$.
(E) $f=2 x_{2}\left(y_{0}^{2}-x_{0} z_{0}\right)+x_{1}^{2} z_{0}-2 x_{1} y_{1} y_{0}+y_{1}^{2} x_{0}+a_{3}\left(x_{0}, y_{0}, z_{0}, w_{1}, \ldots, w_{m}\right)(n \geq 5)$, satisfying the conditions of Proposition 8.5. This is invariant by $x_{0} \partial / \partial x_{1}+x_{1} \partial / \partial x_{2}+$ $y_{0} \partial / \partial y_{1}$ and $y_{0} \partial / \partial x_{1}+y_{1} \partial / \partial x_{2}+a_{1} \partial / \partial y_{1}$; any non-zero linear combination of these is in Case 5. If $V_{3}$ is non-singular, $\mu(V)=25.2^{n-5}$ and the singularity of $V$ is semi-quasi-homogeneous of degree 6 in variables of weights $1,1,3,3,3$ and $2(n-5$ times). In general, we have $\mu(V)=25.2^{n-5}+\mu\left(V_{3}\right)$.

## References

[1] Dimca, A., Singularities and topology of hypersurfaces, Springer-Verlag, 1992.
[2] du Plessis, A. A. and C. T. C. Wall, Curves in $P^{2}(\mathbb{C})$ with 1-dimensional symmetry, Revista Mat Complutense 12 (1999) 117-132.
[3] du Plessis, A. A. and C. T. C. Wall, Applications of discriminant matrices, in Aspects des Singularités, Proc. of Lille singularities semester, online at http: //www-gat.univ-lille1.fr/~tibar/Aspects/index.htm
[4] du Plessis, A. A. and C. T. C. Wall, Hypersurfaces in $P^{n}$ with 1-parameter symmetry groups, Proc. Roy Soc. London A 456 (2000) 2515-2541.
[5] du Plessis, A. A. and C. T. C. Wall, Hypersurfaces with isolated singularities with symmetry, to appear in Proceedings of 2006 Sao Carlos conference.
[6] du Plessis, A. A. and C. T. C. Wall, Discriminants, vector fields and singular hypersurfaces, pp 351-377 in New developments in singularity theory (eds. D. Siersma, C. T. C. Wall and V. Zakalyukin), Kluwer Acad. Publ. 2001.
[7] Nowicki, A., Polynomial derivations and their rings of constants, Uniwersytet Nikolaja Kopernika, Turun 1994.
[8] Steenbrink, J. H. M., Intersection form for quasi-homogeneous singularities, Compositio Math. 34 (1977) 211-223.
[9] Tan, L., An algorithm for explicit generators of the invariants of the basic $G_{a}$-actions, Comm. in Algebra 17 (1989), 565-572.
[10] Wall, C. T. C., Notes on the classification of singularities, Proc. London Math. Soc. 48 (1984), 461-513.
[11] Weitzenbock, R., Ueber die invarianten von lineare Gruppen, Acta Math. 58 (1932), 231-293.

