# Prehomogeneous Vector Spaces and Ergodic Theory III 

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Communicated by K. Rubin
Received May 28, 1997

INTRODUCTION
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Definition (0.1). Let $G$ be a connected reductive group, $V$ a representation of $G$, and $\chi$ a non-trivial character of $G$, all defined over $k$. Then ( $G, V, \chi$ ) is called a prehomogeneous vector space if it satisfies the following properties.
(1) There exists a Zariski open orbit.
(2) There exists a non-zero polynomial $\Delta(x) \in k[V]$ such that $\Delta(g x)=\chi(g) \Delta(x)$.

Such $\Delta(x)$ is called a relative invariant polynomial. We define $V^{\mathrm{ss}}=\{x \in V \mid \Delta(x) \neq 0\}$ and call it the set of semi-stable points. If $(G, V, \chi)$ is an irreducible representation, the choice of $\chi$ is essentially unique and we may write $(G, V)$ as well. The theory of prehomogeneous vector spaces was initiated by Sato-Shintani [21] and Shintani [24]. If ( $G, V$ ) is irreducible, the classification is known (see [20]).

In parts one and two [27, 25], we considered cases (2), (5), (6), (7) in the Sato-Kimura classification [20]. In this part, we consider the prehomogeneous vector space $G=\mathrm{GL}(5) \times \mathrm{GL}(3), V=\Lambda^{2} k^{5} \otimes k^{3}$. Except for the case (29) in [20] (for which we haven't carried out our program), the remaining applicable cases are either very easy or closely related to the spin or half spin representations of the spin groups. In this sense, the present case is rather an isolated case and that's why we consider this single case separately here. However, this case is quite interesting, because it

[^0]produces a family of (irrational) cubic forms in five variables whose values at integer points are dense in $\mathbb{R}$. This may be the family of cubic forms in the lowest number of variables we can achieve by the theory of prehomogeneous vector spaces. We consider most of the remaining irreducible split cases in part four.

Oppenheim conjectured in [13] that if $Q(x)$ is a real non-degenerate indefinite quadratic form in $n \geqslant 5$ variables such that the ratio of at least one pair of coefficients is irrational, for any $\varepsilon>0$, there exists $x \in \mathbb{Z}^{n}$ such that $0<|Q(x)|<\varepsilon$. Due to the result of Lewis [9], this is equivalent to saying the set $\left\{Q(x) \mid x \in \mathbb{Z}^{n}\right\}$ is dense in $\mathbb{R}$. There were many partial results including the one by Davenport with the collaboration with others [4, 5, $6,2,19]$ for $n \geqslant 21$. It was proved in the final form by Margulis (see [11]) for $n \geqslant 3$ using ergodic theory.

We posed the question of generalizing the Oppenheim conjecture from the viewpoint of prehomogeneous vector spaces in [27]. For more detailed comments, the reader should see the introduction of [27]. Here, we briefly state what we are going to prove.

Let $H_{1}=\operatorname{SL}(5), H_{2}=\operatorname{SL}(3), H=H_{1} \times H_{2}$. It is known that $(G, V)$ is a prehomogeneous vector space (see $[24,29,27]$ for the definition of prehomogeneous vector spaces). A non-constant polynomial $\Delta(x)$ on $V$ is called a relative invariant polynomial if there exists a character $\chi$ such that $\Delta(g x)=\chi(g) \Delta(x)$. Such $\Delta(x)$ exists for our case and is essentially unique. So we define $V^{\text {ss }}=\{x \in V \mid \Delta(x) \neq 0\}$. For $x \in V_{\mathbb{R}}^{\text {ss }}$, let $H_{x \mathbb{R}+}^{\circ}$ be the identity component in classical topology of the stabilizer $H_{x \mathbb{R}}$. We will prove that if $x \in V_{\mathbb{R}}^{\mathrm{ss}}$ is "sufficiently irrational" (see Theorem (6.2) for the precise definition), $H_{x \mathbb{R}+}^{\circ} H_{\mathbb{Z}}$ is dense in $H_{\mathbb{R}}$.

What Margulis did was to prove the above statement for the case $H=\operatorname{SL}(3), V=\operatorname{Sym}^{2}\left(\mathbb{R}^{3}\right)^{*}$. Our method is based on the following theorem due to Ratner.

Theorem (0.2) (Ratner). Let $G$ be a connected Lie group and $U$ a connected subgroup of $G$ generated by unipotent elements of $G$. Then given any lattice $\Gamma \subset G$ and $x \in G / \Gamma$, there exists a connected closed subgroup $U \subset F \subset G$ such that $\overline{U x \Gamma}=F x \Gamma$. Moreover, $F / F \cap x \Gamma x^{-1}$ has a finite invariant measure.

Note that in the above theorem, the definition of a lattice contains the condition that $G / \Gamma$ has a finite volume. The first statement was called Raghunathan's topological conjecture, and the second statement was proved by Ratner in conjunction with Raghunathan's topological conjecture. Raghunathan's topological conjecture was published by Dani [3] for one dimensional unipotent groups and was generalized to groups generated by unipotent elements by Margulis [10]. The proof for the general case
was given by Ratner in a series of papers [14-17]. For these, there is an excellent survey article by Ratner [18].

Note that in the above theorem, if $G$ is an algebraic group over $\mathbb{Q}$ and $\Gamma$ is an arithmetic lattice, the group $F$ becomes an algebraic group defined over $\mathbb{Q}$. For this the reader should see Proposition (3.2) [22, pp. 321-322]. It is also proved in Proposition (3.2) [22, pp. 321-322] that the radical of $F$ is a unipotent subgroup. In [22], only one lattice is considered, but one can deduce the above statement for any lattice commensurable with the lattice in [22] by a simple argument using Ratner's theorem.

We describe an application of the density of $H_{x \mathbb{R}+}^{\circ} H_{\mathbb{Z}}$ in $H_{\mathbb{R}}$. For any non-zero point $x$ in a vector space, we denote the point in the corresponding projective space determined by $x$ by $[x]$. Let $V_{1}$ be a five dimensional vector space defined over $\mathbb{Q}$. We fix a rational basis $\left\{m_{0}, \ldots, m_{4}\right\}$ for $V_{1}$.

Let

$$
\begin{aligned}
& Q(a)=a_{0} a_{4}-\frac{1}{4} a_{1} a_{3}+\frac{1}{12} a_{2}^{2}, \\
& F(a)=72 a_{0} a_{2} a_{4}+9 a_{1} a_{2} a_{3}-2 a_{2}^{3}-27 a_{0} a_{3}^{2}-27 a_{1}^{2} a_{4}
\end{aligned}
$$

for $a=\sum_{i=0}^{4} a_{i} m_{i}$. If we identify $V_{1}$ with the space of binary quartic forms by $a \rightarrow a_{0} v_{1}^{4}+\cdots+a_{4} v_{2}^{4}$ ( $v_{1}, v_{2}$ are variables), $Q, F$ correspond to quadratic and cubic $\operatorname{SL}(2)$-invariant polynomials.

If $g \in \mathrm{GL}\left(V_{1}\right)_{\mathbb{R}} \cong \mathrm{GL}(5)_{\mathbb{R}}$, it naturally acts on $\mathbb{P}\left(\operatorname{Sym}^{2} V_{1}^{*}\right)_{\mathbb{R}}$ and $\mathbb{P}\left(\operatorname{Sym}^{3} V_{1}^{*}\right)_{\mathbb{R}}$. Note that $(g Q)(a)=Q\left(g^{-1} a\right),(g F)(a)=F\left(g^{-1} a\right)$. Then the following theorem follows from the consideration of $\overline{H_{x \mathbb{R}+}} H_{\mathbb{Z}}^{\circ}$.

Theorem (0.3). Suppose $g[Q] \notin \mathbb{P}\left(\operatorname{Sym}^{2} V_{1}^{*}\right)_{\mathbb{Q}}$. Then the set of values of the cubic polynomial $F\left(g^{-1} a\right)$ at primitive integer points in $\mathbb{Z}^{5}$ is dense in $\mathbb{R}$.

In Section 1, we consider various identifications concerning tensor products of vector spaces. If $x \in V_{\mathbb{R}}^{\mathrm{ss}}$, by Ratner's theorem (Theorem (0.2)), there exists a closed connected subgroup $H_{x \mathbb{R}+}^{\circ} \subset F \subset H_{\mathbb{R}}$ such that $\overline{H_{x \mathbb{R}+}^{\circ} H_{\mathbb{Z}}}=F H_{\mathbb{Z}}$. In Section 2, we construct equivariant maps from $V^{\text {ss }}$ to various $H$-varieties.

We can summarize how we construct equivariant maps in Section 2 in the following manner.


The map $V^{\text {ss }} \rightarrow V_{1}^{*} \otimes \operatorname{Sym}^{2} V_{2}$ is similar to the one in [26], and $V_{1}^{*} \otimes \operatorname{Sym}^{2} V_{2} \rightarrow \operatorname{Sym}^{2} V_{2}^{*}$ is simply the Castling transform in [20]. It turns out that if $x \in V^{\text {ss }}$, the corresponding quadratic form in three variables is non-degenerate. Using this quadratic form, we can identify $V_{2}$ with its dual, and hence getting a map $V_{1}^{*} \otimes \operatorname{Sym}^{2} V_{2} \rightarrow V_{1}^{*} \otimes$ $\operatorname{Hom}\left(V_{2}, V_{2}\right)$. Regarding an element of $V_{1}^{*} \otimes \operatorname{Hom}\left(V_{2}, V_{2}\right)$ as a $3 \times 3$ matrix $M(v)$ with entries in the space of linear forms in five variables $v=\left(v_{1}, \ldots, v_{5}\right)$, we can consider $\operatorname{tr}\left(M(v)^{2}\right)$, $\operatorname{det} M(v)$, which are a quadratic form and a cubic form in five variables respectively. This rather indirect way is how cubic forms in five variables arise from points in $V^{\text {ss }}$.

In Section 3, we prove that these equivariant maps are well defined and are non-trivial. These equivariant maps correspond to families of such $F$ 's with the property that it $X_{F}$ is the corresponding $H$-variety, $F$ has a unique fixed point in $X_{F}$. Part of our consideration resembles the argument in [23]. In Section 4, we describe the orbit space to determine when $H_{x \mathbb{R}+}^{\circ}$ is generated by unipotent elements. In Section 5, we classify all $F$ s as above. In Section 6, we prove Theorem (0.3).

## 1. PRELIMINARIES

We are going to do a lot of computations in Section 3 regarding symmetric tensor products of vector spaces. We fix various normalizations for that purpose in this section.

Let $W$ be a vector space over $k$ with a basis $\left\{e_{1}, \ldots, e_{n}\right\}$. Let $W^{*}$ be the dual space with the dual basis $\left\{f_{1}, \ldots, f_{n}\right\}$. For $a_{1}, \ldots, a_{d} \in W$, we define

$$
\left[a_{1}, \ldots, a_{d}\right]_{d}=\frac{1}{d!} \sum_{\sigma \in \mathfrak{S}_{d}} a_{\sigma(1)} \otimes \cdots \otimes a_{\sigma(d)},
$$

where $\mathfrak{\Im}_{d}$ is the group of permutations of $\{1, \ldots, d\}$. We identify $\operatorname{Sym}^{d} W$ with the subspace of $W^{\otimes d}$ spanned by elements of the form $\left[a_{1}, \ldots, a_{d}\right]_{d}$. Similarly, we identify $\operatorname{Sym}^{d} W^{*}$ with a subspace of $\left(W^{*}\right)^{\otimes d}$. For $a_{1}, \ldots, a_{d_{1}}, a_{d_{1}+1}, \ldots, a_{d_{1}+d_{2}} \in W$, we define

$$
\left[a_{1}, \ldots, a_{d_{1}}\right]_{d_{1}}\left[a_{d_{1}+1}, \ldots, a_{d_{1}+d_{2}}\right]_{d_{2}}=\left[a_{1}, \ldots, a_{d_{1}+d_{2}}\right]_{d_{1}+d_{2}} .
$$

By this product, $\oplus \operatorname{Sym}^{*} W$ becomes an associative algebra. Since $\left[a_{1}, \ldots, a_{d}\right]_{d}=a_{1} \cdots a_{d}$, we use this usual notation of product from now on.

Since $\left(W^{*}\right)^{\otimes d}$ and $W^{\otimes d}$ are dual spaces of each other, there is a natural pairing between $\operatorname{Sym}^{d} W$ and $\operatorname{Sym}^{d} W^{*}$. If $a=a_{1} \cdots a_{d} \in \operatorname{Sym}^{d} W$, $b=b_{1} \cdots b_{d} \in \operatorname{Sym}^{d} W^{*}$, we normalize this pairing by

$$
(a, b)_{d}=(b, a)_{d}=\frac{1}{d!} \sum_{\sigma \in \mathbb{E}_{d}} b_{1}\left(a_{\sigma(1)}\right) \cdots b_{d}\left(a_{\sigma(d)}\right)
$$

Then if $i_{1}+\cdots+i_{n}=d$,

$$
\left(e_{1}^{i_{1}} \cdots e_{n}^{i_{n}}, f_{1}^{i_{1}} \cdots f_{n}^{i_{n}}\right)_{d}=\frac{i_{1}!\cdots i_{n}!}{d!}
$$

Therefore, $\operatorname{Sym}^{d} W^{*}$ can be identified with the dual space of $\operatorname{Sym}^{d} W$.
The map

$$
W \ni a \rightarrow i_{d}(a)=a \otimes a \otimes \cdots \otimes a \in \operatorname{Sym}^{d} W
$$

is a polynomial map. So if $f \in \operatorname{Sym}^{d} W^{*}, f(a)=f\left(i_{d}(a)\right)$ is a polynomial map from $W$ to $k$ and is homogeneous of degree $d$. We can identify $\operatorname{Sym}^{d} W^{*}$ with the space of degree $k$ forms on $W$ by this correspondence. If $a=\sum_{i=1}^{n} a_{i} e_{i}$,

$$
i_{d}(a)=\sum \frac{d!}{i_{1}!\cdots i_{n}!} a_{1}^{i_{1}} \cdots a_{n}^{i_{n}} e_{1}^{i_{1}} \cdots e_{n}^{i_{n}}
$$

where the sum is over all $\left(i_{1}, \ldots, i_{n}\right)$ such that $i_{1}+\cdots+i_{n}=d$. So if $f=f_{1}^{i_{1}} \cdots f_{n}^{i_{n}}, f(a)=a_{1}^{i_{1}} \cdots a_{n}^{i_{n}}$. Therefore, $f$ corresponds to the monomial $a_{1}^{i_{1}} \cdots a_{n}^{i_{n}}$.

If $G \subset \mathrm{GL}(W)$ is a subgroup, $G$ acts on $W^{*}$ by $(g f)(v)=f\left(g^{-1} v\right)$ for $g \in G, f \in W^{*}$. Whenever we consider the contragredient representation, we consider this action.

## 2. DEFINITIONS OF EQUIVARIANT MAPS

Let $G, V, H$ be as in the introduction. We construct $H$-equivariant maps from $V$ or $V^{\text {ss }}$ to various $H$-varieties in this section.

Let $W=k^{2}$ be the space of two dimensional column vectors. Let $\left\{e_{1}, e_{2}\right\}$ be the standard basis of $W$. Consider the usual action of GL(2) on $W$. This induces an action of GL(2) on $\operatorname{Sym}^{d} W$ and $\operatorname{Sym}^{d} W^{*}$ for any $d$. We define new actions of GL(2) on $\operatorname{Sym}^{2} W, \operatorname{Sym}^{4} W$ by $g \cdot x=(\operatorname{det} g)^{-1} g x$, $(\operatorname{det} g)^{-2} g x$ where $g \cdot x$ is the new action and $g x$ is the usual action. Note that scalar matrices act trivially and therefore, this defines an action of PGL(2) on $\operatorname{Sym}^{2} W, \operatorname{Sym}^{4} W$.

Let $V_{1}=\operatorname{Sym}^{4} W, \quad V_{2}=\operatorname{Sym}^{2} W, \quad V=\wedge^{2} V_{1} \otimes V_{2}$ and $G_{1}=\mathrm{GL}\left(V_{1}\right) \cong$ $\mathrm{GL}(5), G_{2}=\mathrm{GL}\left(V_{2}\right) \cong \mathrm{GL}(3), G=G_{1} \times G_{2}$. Then $G$ acts on $V$ in the usual manner and the above action of PGL(2) defines a homomorphism PGL(2) $\rightarrow G$. By Schur's lemma, this is an imbedding. In fact, $\operatorname{Ker}\left(\operatorname{PGL}(2) \rightarrow G_{1}\right)$ $=\operatorname{Ker}\left(\operatorname{PGL}(2) \rightarrow G_{2}\right)=\{1\}$. So we regard $\operatorname{PGL}(2)$ as a subgroup of $G$. Let $\tilde{T}=\operatorname{Ker}(G \rightarrow \mathrm{GL}(V))$. By Schur's lemma again,

$$
\widetilde{T}=\left\{\left(t I_{5}, t^{-2} I_{3}\right) \mid t \in \mathrm{GL}(1)\right\} \cong \mathrm{GL}(1) .
$$

If $\left(t I_{5}, t^{-2} I_{3}\right) \in \operatorname{PGL}(2)$, it acts trivially on $V_{1}, V_{2}$. So $t=1$. Therefore, $\operatorname{PGL}(2) \cap \widetilde{T}=\{1\}$.

Let $l_{0}=e_{1}^{2}, l_{1}=e_{1} e_{2}, l_{2}=e_{2}^{2}$ and $m_{0}=e_{1}^{4}, \ldots, m_{4}=e_{2}^{4}$. Then $\left\{l_{0}, l_{1}, l_{2}\right\}$, $\left\{m_{0}, \ldots, m_{4}\right\}$ are bases of $V_{2}, V_{1}$ respectively.

We define a linear map $\phi_{1}: V \rightarrow \bigwedge^{4} V_{1} \otimes \operatorname{Sym}^{2} V_{2} \cong V_{1}^{*} \otimes \operatorname{Sym}^{2} V_{2}$ by

$$
\begin{equation*}
V \ni \sum_{i=1}^{N} p_{i} \otimes q_{i} \rightarrow \frac{1}{2} \sum_{i, j=1}^{N} p_{i} \wedge p_{j} \otimes q_{i} q_{j} \tag{2.1}
\end{equation*}
$$

for $p_{1}, \ldots, p_{N} \in \bigwedge^{2} V_{1}, q_{1}, \ldots, q_{N} \in V_{2}$. Regarding $\operatorname{Sym}^{2} V_{2}$ as a subspace of $V_{2} \otimes V_{2}$, we denote the element of $V_{1}^{*} \otimes V_{2} \otimes V_{2}$ which corresponds to $\phi_{1}(x)$ by $\bar{\phi}_{1}(x)$. Regarding $V_{1}^{*}$ as the contragredient representation of $V_{1}, V_{1}^{*} \otimes \operatorname{Sym}^{2} V_{1}$ is a representation of $G$.

The following lemma can be proved as in [20, p. 80], and the proof is left to the reader.

Lemma (2.2). For $g=\left(g_{1}, g_{2}\right) \in G, \phi_{1}(g x)=\operatorname{det} g_{1} g \phi_{1}(x), \bar{\phi}_{1}(g x)=$ $\operatorname{det} g_{1} g \bar{\phi}_{1}(x)$.

If $x=\sum_{0 \leqslant i<j \leqslant 4} m_{i} \wedge m_{j} \otimes x_{i j}$ with $x_{i j} \in V_{2}$,

$$
\phi_{1}(x)=\frac{1}{2} \sum_{0 \leqslant i<j \leqslant 4} \sum_{0 \leqslant k<l \leqslant 4} m_{i} \wedge m_{j} \wedge m_{k} \wedge m_{l} \otimes x_{i j} x_{k l} .
$$

Let $m_{0}^{*}, \ldots, m_{4}^{*} \in \wedge^{4} V_{1}$ be elements such that $m_{i} \wedge m_{j}^{*}=\delta_{i j} m_{0} \wedge$ $\cdots \wedge m_{4}\left(\delta_{i j}\right.$ is Kronecker's delta). Explicitly,

$$
m_{0}=m_{1} \wedge m_{2} \wedge m_{3} \wedge m_{4}, \quad m_{1}=-m_{0} \wedge m_{2} \wedge m_{3} \wedge m_{4}, \text { etc. }
$$

We identify $\wedge^{4} V_{1}$ with the dual space of $V_{1}$ by the pairing

$$
V_{1} \times \bigwedge^{4} V_{1} \ni(a, b) \rightarrow a \wedge b
$$

and choosing $m_{0} \wedge \cdots \wedge m_{4}$ as the basis element of $\wedge^{5} V_{1}$. Then $\left\{m_{0}^{*}, \ldots, m_{4}^{*}\right\}$ can be regarded as the dual basis of $\left\{m_{0}, \ldots, m_{4}\right\}$.

Let

$$
\begin{align*}
& \operatorname{Pfaff}_{0}(x)=x_{12} x_{34}-x_{13} x_{24}+x_{14} x_{23}, \\
& \operatorname{Pfaff}_{1}(x)=-\left(x_{02} x_{34}-x_{03} x_{24}+x_{04} x_{23}\right), \\
& \operatorname{Pfaff}_{2}(x)=x_{01} x_{34}-x_{03} x_{14}+x_{04} x_{13},  \tag{2.3}\\
& \operatorname{Pfaff}_{3}(x)=-\left(x_{01} x_{24}-x_{02} x_{14}+x_{04} x_{12}\right) \\
& \operatorname{Pfaff}_{4}(x)=x_{01} x_{23}-x_{02} x_{13}+x_{03} x_{12} .
\end{align*}
$$

Then

$$
\begin{equation*}
\phi_{1}(x)=\sum_{i=0}^{4} m_{i}^{*} \otimes \operatorname{Pfaff}_{i}(x) . \tag{2.4}
\end{equation*}
$$

The quadratic polynomials $\operatorname{Pfaff}_{0}(x), \ldots, \operatorname{Pfaff}_{4}(x)$ are the Pfaffians of $4 \times 4$ main minors of $x$ if we regard $x$ as an alternating $5 \times 5$ matrix with entries in $V_{1}$. This idea was used in [26] for the case $G=\mathrm{GL}(5) \times \mathrm{GL}(4)$, $V=\wedge^{2} k^{5} \otimes k^{4}$ to parametrize quintic extensions of a given ground field.

Next we consider a linear map $\phi_{2}: V_{1}^{*} \otimes \operatorname{Sym}^{2} V_{2} \rightarrow \bigwedge^{5} \operatorname{Sym}^{2} V_{2} \cong$ Sym $^{2} V_{2}^{*}$ defined by

$$
\begin{equation*}
V_{1}^{*} \otimes \operatorname{Sym}^{2} V_{2} \ni \sum_{i=0}^{4} m_{i}^{*} \otimes p_{i} \rightarrow p_{0} \wedge \cdots \wedge p_{4} \in \bigwedge^{5} \operatorname{Sym}^{2} V_{2} \cong \operatorname{Sym}^{2} V_{2}^{*} \tag{2.5}
\end{equation*}
$$

Note that $\phi_{2}$ is the Castling transform discussed in [20].
Definition (2.6). $\quad \Phi_{1}=\left(1 / 3^{4}\right) \phi_{2} \circ \phi_{1}$.
$\Phi_{1}$ is a map from $V$ to $\operatorname{Sym}^{2} V_{2}^{*}$. The following lemma can also be proved as in [20, p. 80], and the proof is left to the reader.

Lemma (2.7). $\quad \Phi_{1}(g x)=\left(\operatorname{det} g_{1}\right)^{4}\left(\operatorname{det} g_{2}\right) g_{2} \Phi_{1}(x)$.
This map $\Phi_{1}$ was also considered in [12] (using $\phi_{1}$ also) for a different purpose. We will show in Section 3 that the discriminant of $\Phi_{1}(x)$ is not identically zero and $V^{\text {ss }}$ consists of $x$ 's such that $\Phi_{1}(x)$ is non-degenerate.

For $x \in V$, let $\bar{\Phi}_{1}(x)(\alpha, \beta)$ be the symmetric bilinear form on $V_{2}$ associated with $\Phi_{1}(x)$. In other words,

$$
\bar{\Phi}_{1}(x)(\alpha, \beta)=\frac{1}{2}\left(\Phi_{1}(x)(\alpha+\beta)-\Phi_{1}(x)(\alpha)-\Phi_{1}(x)(\beta)\right)
$$

for $\alpha, \beta \in V_{2}$.
We define a linear map $j_{x}: V_{1}^{*} \otimes V_{2} \otimes V_{2} \rightarrow \operatorname{Hom}\left(V_{2}, V_{1}^{*} \otimes V_{2}\right)$ by

$$
\begin{equation*}
j_{x}(a \otimes \alpha \otimes \beta)(\gamma)=\Phi_{1}(x)(\alpha, \gamma) a \otimes \beta \tag{2.8}
\end{equation*}
$$

for $a \in V_{1}^{*}, \alpha, \beta, \gamma \in V_{2}$. If $f \in \operatorname{Hom}\left(V_{2}, V_{1}^{*} \otimes V_{2}\right)$, we define $g f \in \operatorname{Hom}\left(V_{2}\right.$, $\left.V_{1}^{*} \otimes V_{2}\right)$ by

$$
(g f)(\alpha)=g f\left(g_{2}^{-1} \alpha\right)
$$

where we are considering the action of $g$ on the element $f\left(g_{2}^{-1} \alpha\right)$.
Lemma (2.9). $\quad j_{g x}(g(a \otimes \alpha \otimes \beta))=\left(\operatorname{det} g_{1}\right)^{4}\left(\operatorname{det} g_{2}\right) g j_{x}(a \otimes \alpha \otimes \beta)$ for all $a \in V_{1}^{*}, \alpha, \beta \in V_{2}$.

Proof. Let $\gamma \in V_{2}$. Then by Lemma (2.7),

$$
\begin{aligned}
j_{g x}(g(a \otimes \alpha \otimes \beta))(\gamma) & =\bar{\Phi}_{1}(g x)\left(g_{2} \alpha, \gamma\right) g_{1} a \otimes g_{2} \beta \\
& =\left(\operatorname{det} g_{1}\right)^{4}\left(\operatorname{det} g_{2}\right)\left(g_{2} \bar{\Phi}_{1}(x)\right)\left(g_{2} \alpha, \gamma\right) g_{1} a \otimes g_{2} \beta \\
& =\left(\operatorname{det} g_{1}\right)^{4}\left(\operatorname{det} g_{2}\right) \bar{\Phi}_{1}(x)\left(\alpha, g_{2}^{-1} \gamma\right) g_{1} a \otimes g_{2} \beta \\
& =\left(\operatorname{det} g_{1}\right)^{4}\left(\operatorname{det} g_{2}\right) g\left(\bar{\Phi}_{1}(x)\left(\alpha, g_{2}^{-1} \gamma\right) a \otimes \beta\right) \\
& =\left(\operatorname{det} g_{1}\right)^{4}\left(\operatorname{det} g_{2}\right) g\left(j_{x}(a \otimes \alpha \otimes \beta)\left(g_{2}^{-1} \gamma\right)\right) \\
& =\left(\operatorname{det} g_{1}\right)^{4}\left(\operatorname{det} g_{2}\right)\left(g\left(j_{x}(a \otimes \alpha \otimes \beta)\right)\right)(\gamma) .
\end{aligned}
$$

This proves the lemma.
Definition (2.10). $\quad \phi_{3}(x)=j_{x}\left(\bar{\phi}_{1}(x)\right)$.
Apparently, $\phi_{3}$ is a map from $V$ to $\operatorname{Hom}\left(V_{2}, V_{1}^{*} \otimes V_{2}\right)$.
Lemma (2.11). $\quad \phi_{3}(g x)=\left(\operatorname{det} g_{1}\right)^{5}\left(\operatorname{det} g_{2}\right) g \phi_{3}(x)$.
Proof.

$$
\begin{aligned}
\phi_{3}(g x) & =j_{g x}\left(\bar{\phi}_{1}(g x)\right) \\
& =\operatorname{det} g_{1} j_{g x}\left(g \bar{\phi}_{1}(x)\right) \\
& =\left(\operatorname{det} g_{1}\right)^{5}\left(\operatorname{det} g_{2}\right) g j_{x}\left(\bar{\phi}_{1}(x)\right) \\
& =\left(\operatorname{det} g_{1}\right)^{5}\left(\operatorname{det} g_{2}\right) g \phi_{3}(x) .
\end{aligned}
$$

By the basis $\left\{l_{0}, l_{1}, l_{2}\right\}$ for $V_{2}$, we can regard $\phi_{3}(x)$ as a $3 \times 3$ matrix with entries in $V_{1}^{*}$. Then the action of $g=\left(g_{1}, g_{2}\right) \in \operatorname{GL}\left(V_{1}\right) \times \mathrm{GL}(3)$ is obtained by considering $g_{2} \phi_{3}(x) g_{2}^{-1}$ and then applying $g_{1}$ entry-wise. Therefore,

$$
\begin{equation*}
\Phi_{2}(x)=\operatorname{tr}\left(\phi_{3}(x)^{2}\right) \in \operatorname{Sym}^{2} V_{1}^{*}, F_{x}=\operatorname{det} \phi_{3}(x) \in \operatorname{Sym}^{3} V_{1}^{*} \tag{2.12}
\end{equation*}
$$

define maps $x \rightarrow \Phi_{2}(x), F_{x}$ from $V$ to $\operatorname{Sym}^{2} V_{1}^{*}, \operatorname{Sym}^{3} V_{1}^{*}$.

The following lemma is an easy corollary of Lemma (2.11).
Lemma (2.13). (1) $\quad \Phi_{2}(g x)=\left(\operatorname{det} g_{1}\right)^{10}\left(\operatorname{det} g_{2}\right)^{2} g_{1} \Phi_{2}(x)$.
(2) $\quad F_{g x}(a)=\left(\operatorname{det} g_{1}\right)^{15}\left(\operatorname{det} g_{2}\right)^{3} F_{x}\left(g_{1}^{-1} a\right)$ for all $a \in V_{1}$.

For later purposes, we describe how to compute $\Phi_{2}(x), F_{x}$. We have already described how to compute $\phi_{1}(x), \Phi_{1}(x)$ in (2.4), (2.5). Let $\left\{p_{0}, p_{1}, p_{2}\right\}$ be the dual basis of $\left\{l_{0}, l_{1}, l_{2}\right\}$. Suppose

$$
\bar{\phi}_{1}(x)=\sum_{i, j=0}^{2} a_{i j} \otimes l_{i} \otimes l_{j}, \overline{\Phi_{1}}(x)=\sum_{t, s=0}^{2} b_{t s} p_{t} \otimes p_{s}
$$

with $a_{i j} \in V_{1}^{*}, b_{t s} \in k$ for all $i, j, t, s$ and $a_{i j}=a_{j i}, b_{t s}=b_{s t}$. We denote the matrices $\left(a_{i j}\right),\left(b_{t s}\right)$ also by $\bar{\phi}_{1}(x), \bar{\Phi}_{1}(x)$. Let $\gamma=\sum_{s=0}^{2} \gamma_{s} l_{s}$. Then

$$
\begin{aligned}
\phi_{3}(x)(\gamma) & =j_{x}\left(\bar{\phi}_{1}(x)\right)(\gamma) \\
& =\sum_{i, j, s=0}^{2} j_{x}\left(a_{i j} \otimes l_{i} \otimes l_{j}\right)\left(\gamma_{s} l_{s}\right) \\
& =\sum_{i, j, s=0}^{2} \bar{\Phi}_{1}(x)\left(l_{i} \otimes l_{s}\right) a_{i j} \gamma_{s} \otimes l_{j} \\
& =\sum_{i, j, s=0}^{2} a_{i j} b_{i s} \gamma_{s} \otimes l_{j} \\
& =\sum_{i, j, s=0}^{2} a_{j i} b_{i s} \gamma_{s} \otimes l_{j} .
\end{aligned}
$$

Note that $\sum_{i=0}^{2} a_{j i} b_{i s}$ is the $(j, s)$-entry of the matrix product $\bar{\phi}_{1}(x)$ $\bar{\Phi}_{1}(x)$. So if we regard $\phi_{3}(x)$ as a $3 \times 3$ matrix with entries in $V_{1}^{*}$, we get the following relation

$$
\begin{equation*}
\phi_{3}(x)=\bar{\phi}_{1}(x) \bar{\Phi}_{1}(x) . \tag{2.14}
\end{equation*}
$$

## 3. EQUIVARIANT MAPS AT $W$

In this section, we prove that the equivariant maps we constructed in Section 2 are well defined and are non-trivial by evaluating them at a point $w$, which we will define in (3.10).

Let $\left\{l_{0}, l_{1}, l_{2}\right\},\left\{m_{0}, \ldots, m_{4}\right\}$ be the bases of $\operatorname{Sym}^{2} W, \operatorname{Sym}^{4} W$ we defined in Section 2. Let $\left\{p_{0}, p_{1}, p_{2}\right\}$ be the dual basis of $\left\{l_{0}, l_{1}, l_{2}\right\}$. Let $\mathfrak{b}$ be the Lie algebra of PGL(2). Then $\mathfrak{b}$ is the Lie algebra of SL(2) also. We consider $V$ as a representation of $\mathfrak{h}$ also. Let $\Lambda$ be the fundamental dominant
weight of $\mathfrak{h}$. We denote the irreducible representation of $\mathfrak{h}$ with highest weight $d \Lambda$ also by $d \Lambda$. Then by considering weights, $\Lambda^{2} V_{1} \cong 6 \Lambda \oplus 2 \Lambda$ and

$$
\begin{equation*}
6 \Lambda \otimes 2 \Lambda \cong 8 \Lambda \oplus 6 \Lambda \oplus 4 \Lambda, \quad 2 \Lambda \otimes 2 \Lambda \cong 4 \Lambda \oplus 2 \Lambda \oplus k, \tag{3.1}
\end{equation*}
$$

where $k$ is the trivial representation.
Therefore, $V$ contains the trivial representation precisely once.
Note that $V_{2}^{*} \cong \mathfrak{h}$ as an $\mathfrak{h}$-module. We identify $\operatorname{Sym}^{4} W^{*}$ as the space of homogeneous polynomials of degree four in two variables $v={ }^{t}\left(v_{1} v_{2}\right)$ $\left(v\right.$ corresponds to $\left.v_{1} e_{1}+v_{2} e_{2}\right)$. We identify $a=a(v)=a_{0} v_{1}^{4}+\cdots+a_{4} v_{2}^{4}$ with $\left(a_{0}, \ldots, a_{4}\right)$.

Let

$$
H_{0}=\left(\begin{array}{ll}
0 & 1  \tag{3.2}\\
0 & 0
\end{array}\right), \quad H_{1}=\left(\begin{array}{rr}
1 & 0 \\
0 & -1
\end{array}\right), \quad H_{2}=\left(\begin{array}{rr}
0 & 0 \\
-1 & 0
\end{array}\right) .
$$

An easy way to compute Lie algebra actions is to consider values in the ring of dual numbers $k[\varepsilon] /\left(\varepsilon^{2}\right)$.

Let

$$
\begin{array}{ll}
A_{0}=\left(\begin{array}{lllll}
0 & 1 & 0 & 0 & 0 \\
0 & 0 & 2 & 0 & 0 \\
0 & 0 & 0 & 3 & 0 \\
0 & 0 & 0 & 0 & 4 \\
0 & 0 & 0 & 0 & 0
\end{array}\right), & A_{0}^{\prime}=\left(\begin{array}{lll}
0 & 1 & 0 \\
0 & 0 & 2 \\
0 & 0 & 0
\end{array}\right), \\
A_{1}=\left(\begin{array}{lllll}
4 & 0 & 0 & 0 & 0 \\
0 & 2 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & -2 & 0 \\
0 & 0 & 0 & 0 & -4
\end{array}\right), \quad A_{1}^{\prime}=\left(\begin{array}{lll}
2 & 0 & 0 \\
0 & 0 & 0 \\
0 & 0 & -2
\end{array}\right),  \tag{3.3}\\
A_{2}=\left(\begin{array}{lllll}
0 & 0 & 0 & 0 & 0 \\
4 & 0 & 0 & 0 & 0 \\
0 & 3 & 0 & 0 & 0 \\
0 & 0 & 2 & 0 & 0 \\
0 & 0 & 0 & 1 & 0
\end{array}\right), \quad A_{2}^{\prime}=\left(\begin{array}{lll}
0 & 0 & 0 \\
2 & 0 & 0 \\
0 & 1 & 0
\end{array}\right) .
\end{array}
$$

Then the actions of $\mathfrak{b}$ on $V_{1}, V_{2}$ are easy to describe and with respect to the bases $\left\{e_{0}, \ldots, e_{4}\right\},\left\{l_{0}, l_{1}, l_{2}\right\}, H_{0}, H_{1}, H_{2}$ map to

$$
\left(A_{0}, A_{0}^{\prime}\right),\left(A_{1}, A_{1}^{\prime}\right),-\left(A_{2}, A_{2}^{\prime}\right)
$$

respectively.

Since the action of $g \in \operatorname{PGL}(2)$ on $a(v) \in \operatorname{Sym}^{4} W^{*}$ is $a\left(g^{-1} v\right)$, the action of $H \in \mathfrak{h}$ on $a=a(v)$ is given by

$$
a((1-\varepsilon H) v)=a(v)+\varepsilon(H a)(v) .
$$

Then easy computations show that

$$
\begin{align*}
& H_{0} a=\left(0,-4 a_{0},-3 a_{1},-2 a_{2},-a_{3}\right), \\
& H_{1} a=\left(-4 a_{0},-2 a_{1}, 0,2 a_{3}, 4 a_{4}\right),  \tag{3.4}\\
& H_{2} a=\left(a_{1}, 2 a_{2}, 3 a_{3}, 4 a_{4}, 0\right) .
\end{align*}
$$

For $a=\left(a_{0}, \ldots, a_{4}\right), b=\left(b_{0}, \ldots, b_{4}\right)$, we define

$$
\begin{equation*}
Q(a, b)=a_{0} b_{4}-\frac{1}{4} a_{1} b_{3}+\frac{1}{6} a_{2} b_{2}-\frac{1}{4} a_{3} b_{1}+a_{4} b_{0} . \tag{3.5}
\end{equation*}
$$

Then $Q$ is a non-degenerate symmetric bilinear form, invariant under the action of PGL(2). This implies $Q(H a, b)$ is an alternating form for any $H \in \mathfrak{h}$. We regard this alternating form as an element of $\bigwedge^{2}\left(\operatorname{Sym}^{4} W^{*}\right)^{*} \cong$ $\wedge^{2} \operatorname{Sym}^{4} W$.

Lemma (3.6). The map $H \rightarrow f_{H}(a, b)=Q(H a, b)$ is an $\mathfrak{h}$-homomorphism from $\mathfrak{h}$ to $\wedge^{2} \operatorname{Sym}^{4} W$.

Proof. Note that the action of $H \in \mathfrak{h}$ on an element $f(a, b)$ in $\wedge^{2}\left(\operatorname{Sym}^{4} W^{*}\right)^{*}$ is given by $(H f)(a, b)=-f(H a, b)-f(a, H b)$. So if $H, H^{\prime} \in \mathfrak{h}$,

$$
\begin{aligned}
\left(H^{\prime} f_{H}\right)(a, b) & =-f_{H}\left(H^{\prime} a, b\right)-f_{H}\left(a, H^{\prime} b\right) \\
& =-Q\left(H H^{\prime} a, b\right)-Q\left(H a, H^{\prime} b\right) \\
& =-Q\left(H H^{\prime} a, b\right)+Q\left(H^{\prime} H a, b\right) \\
& =Q\left(\left[H^{\prime}, H\right] a, b\right) \\
& =f_{\left[H^{\prime}, H\right]}(a, b) .
\end{aligned}
$$

This defines an $\mathfrak{h}$-homomorphism $\mathfrak{h} \rightarrow \bigwedge^{2} V_{1}$. Regarding this homomorphism as an element of $\bigwedge^{2} V_{1} \otimes \mathfrak{h}^{*} \cong \bigwedge^{2} V_{1} \otimes V_{2}=V$, we get a fixed point of $V$ under the action of $\operatorname{PGL}(2)$. We compute this element explicitly.

Note that the linear map defined by

$$
\begin{equation*}
H_{0} \rightarrow \frac{1}{2} v_{2}^{2}, \quad H_{1} \rightarrow v_{1} v_{2}, \quad H_{2} \rightarrow \frac{1}{2} v_{1}^{2} \tag{3.7}
\end{equation*}
$$

is an $\mathfrak{h}$-homomorphism.

By (3.4),

$$
\begin{align*}
& Q\left(H_{0} a, b\right)=a_{0} b_{3}-\frac{1}{2} a_{1} b_{2}+\frac{1}{2} a_{2} b_{1}-a_{3} b_{0}, \\
& Q\left(H_{1} a, b\right)=-4 a_{0} b_{4}+\frac{1}{2} a_{1} b_{3}-\frac{1}{2} a_{3} b_{1}+4 a_{4} b_{0},  \tag{3.8}\\
& Q\left(H_{2} a, b\right)=a_{1} b_{4}-\frac{1}{2} a_{2} b_{3}+\frac{1}{2} a_{3} b_{2}-a_{4} b_{1} .
\end{align*}
$$

Note that $\left\{e_{1}^{2}, 2 e_{1} e_{2}, e_{2}^{2}\right\}$ is the dual basis of $\left\{v_{1}^{2}, v_{1} v_{2}, v_{2}^{2}\right\}$, and $\left(m_{0}, v_{1}^{4}\right)_{4}$ $=1,\left(m_{1}, v_{1}^{3} v_{2}\right)_{4}=\frac{1}{4}$, etc. We identify $\wedge^{2} V_{1}$ with the space of alternating bilinear forms on $V_{1}^{*}$ by assuming

$$
m \wedge m^{\prime}(a, b)=(m, a)_{1}\left(m^{\prime}, b\right)_{1}-(m, b)_{1}\left(m^{\prime}, a\right)_{1}
$$

for $m, m^{\prime} \in V_{1}, a, b \in V_{1}^{*}$. So by the corresponding $H \rightarrow f_{H}$,

$$
\begin{align*}
& H_{0} \rightarrow 4 m_{0} \wedge m_{3}-12 m_{1} \wedge m_{2} \\
& H_{1} \rightarrow-4 m_{0} \wedge m_{4}+8 m_{1} \wedge m_{3}  \tag{3.9}\\
& H_{2} \rightarrow 4 m_{1} \wedge m_{4}-12 m_{2} \wedge m_{3}
\end{align*}
$$

Since $\left\{H_{0}, H_{1}, H_{2}\right\}$ corresponds to $\left\{\frac{1}{2} v_{2}^{2}, v_{1} v_{2}, \frac{1}{2} v_{1}^{2}\right\}$ and $\left\{2 l_{2}, 2 l_{1}, 2 l_{0}\right\}$ is its dual basis, this correspondence can be regarded as the element $8 w$ where

$$
\begin{align*}
w= & \left(m_{0} \wedge m_{3}-3 m_{1} \wedge m_{2}\right) \otimes l_{2}+\left(-m_{0} \wedge m_{4}+2 m_{1} \wedge m_{3}\right) \otimes l_{1} \\
& +\left(m_{1} \wedge m_{4}-3 m_{2} \wedge m_{3}\right) \otimes l_{0} . \tag{3.10}
\end{align*}
$$

These considerations show the following proposition.
Proposition (3.11). The element $w \in V$ is fixed by PGL(2).
In [20, p. 95], instead of $w$, the element

$$
w^{\prime}=\left(m_{0} \wedge m_{1}+m_{2} \wedge m_{3}, m_{1} \wedge m_{2}+m_{3} \wedge m_{4}, m_{0} \wedge m_{2}+m_{1} \wedge m_{4}\right)
$$

was considered (we shifted the indices in [20] because we are using indices $0, \ldots, 4)$. However, by replacing $m_{0}, \ldots, m_{4}$ in $w$ by $m_{4}, m_{2}, m_{0}, m_{1}, m_{3}$ respectively, and multiplying scalars to basis elements of $\wedge^{2} V_{1}$, we get the above element $w^{\prime}$. Therefore, we are considering essentially the same element as in [20].

In [20, p. 96], the Lie algebra of $G_{w^{\prime}}^{\circ} / \widetilde{T}\left(G_{w^{\prime}}^{\circ}\right.$, is the identity component of the stabilizer) is computed and is isomorphic to the Lie algebra of PGL(2). Therefore, this is the case for $w$ also. Since we are assuming ch $k=0$, this implies $G_{w}^{\circ}=\operatorname{PGL}(2) \times \widetilde{T}$ if $k$ is algebraically closed (see [7]). For arbitrary $k$, we still have the inclusion $\operatorname{PGL}(2) \times \widetilde{T} \subset G_{w}^{\circ}$. Since this is an isomorphism over $\bar{k}$, we get the following proposition

Proposition (3.12). $\quad G_{w}^{\circ}=\operatorname{PGL}(2) \times \widetilde{T}$.
By the basis $\left\{e_{0}, \ldots, e_{4}\right\}$, we regard $V$ as the space of $5 \times 5$ alternating matrices with entries in $V_{2}$. Then

$$
w=\left(\begin{array}{ccccc}
0 & 0 & 0 & l_{2} & -l_{1} \\
0 & 0 & -3 l_{2} & 2 l_{1} & l_{0} \\
0 & 3 l_{2} & 0 & -3 l_{0} & 0 \\
-l_{2} & -2 l_{1} & 3 l_{0} & 0 & 0 \\
l_{1} & -l_{0} & 0 & 0 & 0
\end{array}\right) .
$$

By the definition (2.3),

$$
\begin{align*}
& \operatorname{Pfaff}_{0}(w)=-3 l_{0}^{2}, \quad \operatorname{Pfaff}_{1}(w)=-3 l_{0} l_{1}, \\
& \operatorname{Pfaff}_{2}(w)=-2 l_{1}^{2}-l_{0} l_{2},  \tag{3.13}\\
& \operatorname{Pfaff}_{3}(w)=-3 l_{1} l_{2}, \quad \operatorname{Pfff}_{4}(w)=-3 l_{2}^{2} .
\end{align*}
$$

Note that we are regarding them as elements of $\operatorname{Sym}^{2} V_{2}$ and not $\operatorname{Sym}^{4} W$.
By the basis $\left\{l_{0}, l_{1}, l_{2}\right\}$, we regard $\bar{\phi}_{1}(w)$ as a $3 \times 3$ matrix with entries in $V_{1}^{*}$ as in Section 2. Then

$$
\bar{\phi}_{1}(w)=-\left(\begin{array}{ccc}
3 m_{0}^{*} & \frac{3}{2} m_{1}^{*} & \frac{1}{2} m_{2}^{*}  \tag{3.14}\\
\frac{3}{2} m_{1}^{*} & 2 m_{2}^{*} & \frac{3}{2} m_{3}^{*} \\
\frac{1}{2} m_{2}^{*} & \frac{3}{2} m_{3}^{*} & 3 m_{4}^{*}
\end{array}\right) .
$$

Let

$$
\begin{equation*}
n_{0}=l_{0}^{2}, \quad n_{1}=l_{1}^{2}, \quad n_{2}=l_{2}^{2}, \quad n_{3}=l_{0} l_{1}, \quad n_{4}=l_{1} l_{2}, \quad n_{5}=l_{0} l_{2} \in \operatorname{Sym}^{2} V_{2} . \tag{3.15}
\end{equation*}
$$

Then $\left\{n_{0}, \ldots, n_{5}\right\}$ is a basis of $\operatorname{Sym}^{2} V_{2}$. Let $n_{0}^{*}, \ldots, n_{5}^{*}$ be elements of $\wedge^{5} \operatorname{Sym}^{2} V_{2} \cong \operatorname{Sym}^{2} V_{2}^{*}$ such that $n_{i} \wedge n_{j}^{*}=\delta_{i j} n_{0} \wedge \cdots \wedge n_{5}$. Then

$$
\begin{aligned}
& \operatorname{Pfaff}_{0}(w) \wedge \operatorname{Pfaff}_{1}(w) \wedge \operatorname{Pfaff}_{2}(w) \wedge \operatorname{Pfaff}_{3}(w) \wedge \operatorname{Pfaff}_{4}(w) \\
& \quad=\left(-3 n_{0}\right) \wedge\left(-3 n_{3}\right) \wedge\left(-2 n_{1}-n_{5}\right) \wedge\left(-3 n_{4}\right) \wedge\left(-3 n_{2}\right) \\
& \quad=-3^{4} n_{0} \wedge n_{3} \wedge\left(2 n_{1}+n_{5}\right) \wedge n_{4} \wedge n_{2} \\
& \quad=-3^{4}\left(2 n_{0} \wedge n_{3} \wedge n_{1} \wedge n_{4} \wedge n_{2}+n_{0} \wedge n_{3} \wedge n_{5} \wedge n_{4} \wedge n_{2}\right) \\
& \quad=3^{4}\left(n_{1}^{*}-2 n_{5}^{*}\right) .
\end{aligned}
$$

Therefore,

$$
\begin{equation*}
\Phi_{1}(w)=n_{1}^{*}-2 n_{5}^{*} . \tag{3.15}
\end{equation*}
$$

We identify $\operatorname{Sym}^{2} V_{2}^{*}$ with the dual space of $\operatorname{Sym}^{2} V_{2}$. Then with respect to the basis $\left\{p_{0}, p_{1}, p_{2}\right\}, n_{1}^{*}, n_{5}^{*}$ correspond to $p_{1}^{2}, 2 p_{0} p_{2}$. Therefore, $\Phi_{1}(w)=p_{1}^{2}-4 p_{0} p_{2}$. We regard $\bar{\Phi}_{1}(w)$ as a $3 \times 3$ matrix as in Section 2. Then

$$
\bar{\Phi}_{1}(w)=\left(\begin{array}{rrr}
0 & 0 & -2  \tag{3.16}\\
0 & 1 & 0 \\
-2 & 0 & 0
\end{array}\right) .
$$

Therefore,

$$
\phi_{3}(w)=\bar{\phi}_{1}(w) \bar{\Phi}_{1}(w)=\left(\begin{array}{ccc}
m_{2}^{*} & -\frac{3}{2} m_{1}^{*} & 6 m_{0}^{*}  \tag{3.17}\\
3 m_{3}^{*} & -2 m_{2}^{*} & 3 m_{1}^{*} \\
6 m_{4}^{*} & -\frac{3}{2} m_{3}^{*} & m_{2}^{*}
\end{array}\right) .
$$

So, we get the following proposition easily.
Proposition (3.18). Let $a=a_{0} m_{0}+\cdots+a_{4} m_{4}$. Then

$$
\begin{align*}
& \Phi_{2}(w)(a)=\left(\operatorname{tr}\left(\phi_{3}(w)^{2}\right)\right)(a)=72\left(a_{0} a_{4}-\frac{1}{4} a_{1} a_{3}+\frac{1}{12} a_{2}^{2}\right)  \tag{1}\\
& \begin{aligned}
F_{w}(a) & =\left(\operatorname{det} \phi_{3}(w)\right)(a) \\
& =72 a_{0} a_{2} a_{4}+9 a_{1} a_{2} a_{3}-2 a_{2}^{3}-27 a_{0} a_{3}^{2}-27 a_{1}^{2} a_{4}
\end{aligned}
\end{align*}
$$

By these considerations, $\Phi_{1}, \Phi_{2}$ are non-trivial maps. By (3.16), the discriminant $\Delta(x)$ of $\Phi_{1}(x)$ is a non-zero polynomial. By Lemma (2.7), $\Delta(x)$ is a non-constant polynomial. Therefore, it is a relative invariant polynomial. So we reproved that $w \in V_{k}^{\text {ss }}$. Since our case is known to be a regular prehomogeneous vector space, $V_{k}^{\text {ss }}$ is a single $G_{k}$-orbit if $k$ is algebraically closed. Therefore, $V^{\text {ss }}$ consists of $x$ 's such that $\Phi_{1}(x)$ is non-degenerate.

Since $\Phi_{2}(w)$ is non-degenerate, $\Phi_{2}(x)$ is non-degenerate for all $x \in V^{\text {ss }}$.

## 4. THE ORBIT SPACE $G_{K} \backslash V_{K}^{\text {SS }}$

In this section, we prove that $G_{k} \backslash V_{k}^{\text {ss }}$ corresponds bijectively with $\mathrm{GL}(1)_{k} \times \mathrm{GL}(3)_{k}$-equivalence classes of ternary quadratic forms over $k$.

We first recall the relation between the orbit space $G_{k} \backslash V_{k}^{\text {ss }}$ and the Galois cohomology set.

For any algebraic group $G$ over $k$, let $\mathrm{H}^{1}(k, G)$ be the first Galois cohomology set. We choose the definition so that trivial classes are those of the form $\left\{g^{-1} g^{\sigma}\right\}_{\sigma \in \operatorname{Gal}(\bar{k} / k)}\left(g \in G_{\bar{k}}\right)$ and the cocycle condition is $h_{\sigma \tau}=$ $h_{\tau} h_{\sigma}^{\tau}$ for a continuous map $\left\{h_{\sigma}\right\}_{\sigma \in \operatorname{Gal}(\bar{k} / k)}$ from $\operatorname{Gal}(\bar{k} / k)$ to $G_{\bar{k}}$.

Let $(G, V)$ be an arbitrary regular prehomogeneous vector space, and $w \in V_{k}^{\text {ss }}$. Then for any $x \in V_{k}^{\text {ss }}$, there exists $g_{x} \in G_{\bar{k}}$ such that $x=g_{x} w$. Then
$c_{x}=\left\{g_{x}^{-1} g_{x}^{\sigma}\right\}_{\sigma \in \operatorname{Gal}(\bar{k} / k)}$ determines a cohomology class in $\mathrm{H}^{1}\left(k, G_{w}\right)$ and does not depend on the choice of $g_{x}$. The following theorem is due to Igusa [8].

Theorem (4.1) (Igusa). The correspondence

$$
G_{k} \backslash V_{x}^{\mathrm{ss}} \ni x \rightarrow c_{x} \in \operatorname{Ker}\left(\mathrm{H}^{1}\left(k, G_{w}\right) \rightarrow \mathrm{H}^{1}(k, G)\right)
$$

is bijective.
Note that $\operatorname{Ker}\left(\mathrm{H}^{1}\left(k, G_{w}\right) \rightarrow \mathrm{H}^{1}(k, G)\right)$ is the set of elements $c \in \mathrm{H}^{1}\left(k, G_{w}\right)$ which map to the trivial class in $\mathrm{H}^{1}(k, G)$. In our case, $\mathrm{H}^{1}(k, G)$ is trivial. Therefore, $G_{k} \backslash V_{k}^{\text {ss }} \cong \mathrm{H}^{1}\left(k, G_{w}\right)$.

We recall the correspondence between $\mathrm{GL}(1)_{k} \times \mathrm{GL}(3)_{k}$-equivalence classes of ternary quadratic forms and quarternion algebras. Let $V_{3}=$ $\operatorname{Sym}^{2} V_{2}^{*}$. Then $\mathrm{GL}\left(V_{2}\right) \cong \mathrm{GL}(3)$ acts on $V_{3}$ in the usual manner. We let $\mathrm{GL}(1)$ act on $V_{3}$ by the usual multiplication. Then $\mathrm{GL}(1) \times \mathrm{GL}(3)$ acts on $V_{3}$. Let $\left\{l_{0}, l_{1}, l_{2}\right\}$ and $\left\{p_{0}, p_{1}, p_{2}\right\}$ be as before. Let $\bar{w}=p_{1}^{2}-4 p_{0} p_{2}$, and $Q$ the corresponding quadratic form. It is well known (and is easy to verify) that the stabilizer of $\bar{w}$ is isomorphic to $\mathrm{SO}(Q) \times \mathrm{GL}(1)$ and $\mathrm{SO}(Q) \cong \operatorname{PGL}(2)$.

Therefore, $\quad\left(\mathrm{GL}(1)_{k} \times \mathrm{GL}(3)_{k}\right) \backslash V_{3 k}^{\mathrm{ss}} \quad$ corresponds bijectively with $\mathrm{H}^{1}(k$, PGL(2)). Since PGL(2) is isomorphic to the automorphism group of the associative algebra $\mathrm{M}(2,2), \mathrm{H}^{1}(k, \mathrm{PGL}(2))$ corresponds bijectively with isomorphism classes of quarternion algebras. Given a ternary quadratic form, the corresponding quarternion algebra is the Clifford algebra associated with the quadratic form.

Now we go back to our situation. Let $G, H, V$ be as before. We consider the element $w \in V_{k}^{\text {ss }}$ which we defined in (3.10). We pointed out in (3.12) that $G_{w}^{\circ} \cong \mathrm{PGL}(2) \times \mathrm{GL}(1)$.

## Proposition (4.2). The group $G_{w}$ is connected.

Proof. We may assume that $k$ is algebraically closed. Suppose $g \in G_{k} / \widetilde{T}_{k}$. The identity component of the stabilizer of $w$ in $G / \widetilde{T}$ is isomorphic to PGL(2). The conjugation by $g$ induces an automorphism of PGL(2). Since there is no outer automorphism of PGL(2), by changing $g$ if necessary, we may assume that $g$ commutes with elements of PGL(2). Since $V_{1}, V_{2}$ are irreducible representations, by Schur's lemma, $g$ is represented by an element of the form $\left(t_{1} I_{5}, t_{2} I_{3}\right)$. This element fixes $w$ if and only if $t_{1}^{2} t_{2}=1$. So $g=1$ (in $G / \widetilde{T}$ ).

Proposition (4.3). (1) The map $\Phi_{1}: V \rightarrow V_{3}=\operatorname{Sym}^{2} V_{1}^{*}$ induces a bijection $G_{k} \backslash V_{k}^{\text {ss }} \cong\left(\mathrm{GL}(1)_{k} \times \mathrm{GL}(3)_{k}\right) \backslash V_{3 k}^{\text {ss }}$.
(2) If $x \in V_{k}^{\mathrm{ss}}$, the projections of $H_{x}$ to $G_{1}, G_{2}$ induce isomorphisms to the images. In particular, $H_{x} \cong \mathrm{SO}\left(\Phi_{1}(x)\right)$.

Proof. Let $c \in \mathrm{H}^{1}\left(k, G_{w}\right)$. Then $c$ becomes trivial in $\mathrm{H}^{1}(k, H)$ also. Let $g=\left(g_{1}, g_{2}\right) \in H_{\bar{k}}$ be the element such that $c$ is represented by $\left\{g^{-1} g^{\sigma}\right\}_{\sigma \in \operatorname{Gal}(\bar{k} / k)}$. Then the orbit in $V_{k}^{\text {ss }}$ corresponding to $c$ is $g w$. By Lemma (2.7), $\Phi_{1}(g w)=g_{2} \Phi_{1}(w)$.

Since $\mathrm{H}^{1}\left(k, G_{w}\right) \cong \mathrm{H}^{1}(k, \operatorname{PGL}(2))$ and the projection of $\operatorname{PGL}(2)$ to $G_{2}$ is an isomorphism to its image,

$$
\mathrm{H}^{1}\left(k, G_{w}\right) \ni\left\{g^{-1} g^{\sigma}\right\}_{\sigma \in \operatorname{Gal}(\bar{k} / k)} \rightarrow\left\{g_{2}^{-1} g_{2}^{\sigma}\right\}_{\sigma \in \operatorname{Gal}(\bar{k} / k)} \in \mathrm{H}^{1}(k, \operatorname{PGL}(2))
$$

is a bijection. Note that we are considering $\operatorname{PGL}(2) \subset G$ for the first element and PGL $(2) \subset G_{2}$ for the second element. Since $\left(\mathrm{GL}(1)_{k} \times \mathrm{GL}(3)_{k}\right) \backslash V_{3 k}^{\mathrm{ss}} \cong$ $\mathrm{H}^{1}(k, \operatorname{PGL}(2))$, this proves (1).

Note that $(\operatorname{PGL}(2) \times \widetilde{T}) \cap H=\operatorname{PGL}(2)$. So $H_{w} \cong \operatorname{PGL}(2)$. We already pointed out that statement (2) holds for $w$ in Section 2. Let $x \in V_{k}^{\text {ss }}$. By Lemma (2.7), the projection of $H_{x}$ to $G_{2}$ is contained in $\operatorname{SO}\left(\Phi_{1}(x)\right)$. So it is enough to prove (2) when $k$ is algebraically closed. But then $x$ is in the orbit of $w$ and (2) follows easily.

Remark (4.4). The map $\Phi_{1}$ induces a map $G_{k} \backslash V_{k}^{\text {ss }} \rightarrow \mathrm{GL}(3)_{k} \backslash V_{3 k}^{\text {ss }}$ also, but this may not be surjective. This may be regarded as the section $\left(\mathrm{GL}(1)_{k} \times \mathrm{GL}(3)_{k}\right) \backslash V_{3 k}^{\mathrm{ss}} \rightarrow \mathrm{GL}(3)_{k} \backslash V_{3 k}^{\mathrm{ss}}$ defined by $x \rightarrow(\operatorname{det} x)^{-1} x$.

## 5. INTERMEDIATE GROUPS

Let $x \in V_{\mathbb{R}}^{\mathrm{ss}}$. By Proposition (4.3), $H_{x \mathbb{R}}$ is connected in classical topology. So $H_{x \mathbb{R}+}^{\circ}=H_{x \mathbb{R}}$. If $\Phi_{1}(x)$ is definite, $H_{x \mathbb{R}}$ is compact by Proposition (4.3) also. Then $H_{x \mathbb{R}} H_{\mathbb{Z}} \subset H_{\mathbb{R}} / H_{\mathbb{Z}}$ is a compact set. Therefore, an analogue of the Oppenheim conjecture is not applicable to such points. The set of real indefinite non-degenerate ternary quadratic forms is a single $\mathrm{GL}(1)_{\mathbb{R}} \times$ $\mathrm{GL}(3)_{\mathbb{R}_{\mathbb{R}}}$-orbit. Therefore, we only consider $x \in G_{\mathbb{R}} w$.

We determine all the closed connected subgroups between $H_{x \mathbb{R}+}^{\circ}$ and $H_{\mathbb{R}}$ for all $x \in G_{\mathbb{R}} w$ for the rest of this section. This reduces to the consideration of Lie algebras. We consider an arbitrary ground field $k$ of characteristic zero and specialize to $k=\mathbb{R}$ in (5.10).

We first describe possible candidates for such subgroups. By Lemmas (2.7), (2.13), $\Phi_{1}, \Phi_{2}$ are $H$-equivariant maps. As we pointed out at the end of Section $3, \Phi_{1}(x) \in \operatorname{Sym}^{2} V_{1}^{*}, \Phi_{2}(x) \in \operatorname{Sym}^{2} V_{2}^{*}$ are non-degenerate for $x \in G_{k} w$. So let $\operatorname{SO}\left(\Phi_{1}(x)\right), \operatorname{SO}\left(\Phi_{2}(x)\right)$ be the corresponding special orthogonal groups.

In the following definition, $x \in G_{k} w$.

Definition (5.1). (1) $H_{x 1} \subset \mathrm{GL}\left(V_{1}\right), H_{x 2} \subset \mathrm{GL}\left(V_{2}\right)$ are the images of the projections of $H_{x}$ to $G_{1}, G_{2}$ respectively.
(2) $H_{x 3}=\mathrm{SO}\left(\Phi_{2}(x)\right) \subset \mathrm{GL}\left(V_{1}\right)$.

Note that both $H_{x 1}, H_{x 2}$ are isomorphic to $H_{x}$, and $H_{x} \cong \operatorname{PGL}(2)$.
Let $\mathfrak{h}$ be the Lie algebra of $\operatorname{PGL}(2)$ as before. Let $\mathfrak{h}_{1}=\operatorname{sl}(5), \mathfrak{h}_{2}=\operatorname{sl}(3)$ (Lie algebras of $\operatorname{SL}(5)$, $\operatorname{SL}(3)$ ). If $\mathfrak{f}$ is a Lie algebra between $\mathfrak{h}$ and $\mathfrak{h}_{1} \times \mathfrak{h}_{2}$, it is an $\mathfrak{h}$-module. So we first decompose $\mathfrak{h}_{1}, \mathfrak{h}_{2}$ to direct sums of irreducible $\mathfrak{h}$-modules.

Let

$$
\begin{align*}
& B=B\left(b_{0}, \ldots, b_{4}\right)=\left(\begin{array}{ccccc}
2 b_{2} & -3 b_{1} & b_{0} & 0 & 0 \\
12 b_{3} & -b_{2} & -2 b_{1} & 3 b_{0} & 0 \\
6 b_{4} & 3 b_{3} & -2 b_{2} & 3 b_{1} & 6 b_{0} \\
0 & 3 b_{4} & -2 b_{3} & -b_{2} & 12 b_{1} \\
0 & 0 & b_{4} & -3 b_{3} & 2 b_{2}
\end{array}\right), \\
& C=C\left(c_{0}, \ldots, c_{6}\right)=\left(\begin{array}{ccccc}
c_{3} & 3 c_{2} & -c_{1} & c_{0} & 0 \\
12 c_{4} & -2 c_{3} & -4 c_{2} & 0 & 4 c_{0} \\
6 c_{5} & -6 c_{4} & 0 & -6 c_{2} & 6 c_{1} \\
4 c_{6} & 0 & -4 c_{4} & 2 c_{3} & 12 c_{2} \\
0 & c_{6} & -c_{5} & 3 c_{4} & -c_{3}
\end{array}\right),  \tag{5.2}\\
& D=D\left(d_{0}, \ldots, d_{8}\right)=\left(\begin{array}{ccccc}
d_{4} & -d_{3} & d_{2} & -d_{1} & d_{0} \\
4 d_{5} & -4 d_{4} & 4 d_{3} & -4 d_{2} & 4 d_{1} \\
6 d_{6} & -6 d_{5} & 6 d_{4} & -6 d_{3} & 6 d_{2} \\
4 d_{7} & -4 d_{6} & 4 d_{5} & -4 d_{4} & 4 d_{3} \\
d_{8} & -d_{7} & d_{6} & -d_{5} & d_{4}
\end{array}\right), \\
& B^{\prime}=B^{\prime}\left(b_{0}^{\prime}, \ldots, b_{4}^{\prime}\right)=\left(\begin{array}{cccc}
b_{2}^{\prime} & -b_{1}^{\prime} & b_{0}^{\prime} \\
2 b_{3}^{\prime} & -2 b_{2}^{\prime} & 2 b_{1}^{\prime} \\
b_{4}^{\prime} & -b_{3}^{\prime} & b_{2}^{\prime}
\end{array}\right),
\end{align*}
$$

where $b_{0} \cdots \in k$.
We define

$$
\begin{align*}
& U_{2}=\left\{B\left(b_{0}, \ldots, b_{4}\right) \mid b_{0}, \ldots, b_{4} \in k\right\}, \\
& U_{3}=\left\{C\left(c_{0}, \ldots, c_{6}\right) \mid c_{0}, \ldots, c_{6} \in k\right\},  \tag{5.3}\\
& U_{4}=\left\{D\left(d_{0}, \ldots, d_{8}\right) \mid d_{0}, \ldots, d_{8} \in k\right\}, \\
& V_{2}=\left\{B^{\prime}\left(b_{0}^{\prime}, \ldots, b_{4}^{\prime}\right) \mid b_{0}^{\prime}, \ldots, b_{4}^{\prime} \in k\right\} .
\end{align*}
$$

Let $U_{1}, V_{1}$ be the images of $\mathfrak{h}$ in $\mathfrak{h}_{1}, \mathfrak{h}_{2} . U_{1}, V_{1}$ are clearly, sub $\mathfrak{h}$-modules.

Lemma (5.4). The subspaces $U_{2}, U_{3}, U_{4}, V_{2}$ are irreducible sub $\mathfrak{h}$-modules with highest weights $4 \Lambda, 6 \Lambda, 8 \Lambda, 4 \Lambda$ respectively.

Proof. By straightforward computations,

$$
\begin{align*}
{\left[A_{0}, B(1,0, \ldots, 0)\right] } & =\left[A_{0}, C(1,0, \ldots, 0)\right]=\left[A_{0}, D(1,0, \ldots, 0)\right]=0, \\
{\left[A_{0}^{\prime}, B^{\prime}(1,0, \ldots, 0)\right] } & =0, \\
{\left[A_{1}, B(1,0, \ldots, 0)\right] } & =4 B(1,0, \ldots, 0),  \tag{5.5}\\
{\left[A_{1}, C(1,0, \ldots, 0)\right] } & =6 C(1,0, \ldots, 0), \\
{\left[A_{1}, D(1,0, \ldots, 0)\right] } & =8 D(1,0, \ldots, 0), \\
{\left[A_{1}^{\prime}, B^{\prime}(1,0, \ldots, 0)\right] } & =4 B^{\prime}(1,0, \ldots, 0) .
\end{align*}
$$

Also

$$
\begin{align*}
{\left[A_{2}, B\left(b_{0}, \ldots, b_{4}\right)\right] } & =B\left(0, b_{0}, 6 b_{1}, b_{2}, 4 b_{3}\right), \\
{\left[A_{2}, C\left(c_{0}, \ldots, c_{6}\right)\right] } & =C\left(0,2 c_{0}, c_{1},-12 c_{2}, c_{3}, 10 c_{4}, 3 c_{5}\right),  \tag{5.6}\\
{\left[A_{2}, D\left(d_{0}, \ldots, d_{8}\right)\right] } & =D\left(0, d_{0}, 2 d_{1}, 3 d_{2}, \ldots, 8 d_{7}\right), \\
{\left[A_{2}^{\prime}, B^{\prime}\left(b_{0}^{\prime}, \ldots, b_{4}^{\prime}\right)\right] } & =B^{\prime}\left(0, b_{0}^{\prime}, 2 b_{1}^{\prime}, 3 b_{2}^{\prime}, 4 b_{3}^{\prime}\right) .
\end{align*}
$$

The author used MAPLE [1] to find $U_{2}, U_{3}, U_{4}$ but computed (5.5), (5.6) manually. So these computations can be managed manually in principle, but we checked (5.5) (5.6) by MAPLE also.

By (5.6), $U_{2}$ is spanned by elements of the form $\operatorname{ad}\left(A_{2}\right)^{i} B(1,0,0,0,0)$ $(\mathrm{ad}(*)$ is the adjoint representation). Since

$$
\operatorname{ad}\left(A_{0}\right) B(1,0,0,0,0)=0, \quad \operatorname{ad}\left(A_{1}\right) B(1,0,0,0,0)=4 B(1,0,0,0,0),
$$

$U_{2}$ is an irreducible sub $\mathfrak{h}$-module with highest weight $4 \Lambda$.
Other cases are similar.
Proposition (5.7). (1) $\left[U_{2}, U_{2}\right]=U_{1} \oplus U_{3}$.

$$
\begin{equation*}
\left[U_{2}, U_{3}\right]=U_{2} \oplus U_{4} . \tag{2}
\end{equation*}
$$

(3) $\left[U_{2}, U_{4}\right]=U_{3}$.
(4) $\left[U_{3}, U_{3}\right]=U_{1} \oplus U_{3}$.
(5) $\left[U_{3}, U_{4}\right]=U_{2} \oplus U_{4}$.
(6) $\left[U_{4}, U_{4}\right]=U_{1} \oplus U_{3}$.
(7) $\left[V_{2}, V_{2}\right]=V_{1}$.

Proof. We first consider (1). Since $U_{2}$ is irreducible, for any non-zero element $X \in U_{2}, U_{2}$ is generated by $X$ as an $\mathfrak{h}$-module. So [ $U_{2}, U_{2}$ ] is generated by $\left[X, U_{2}\right]$ as an $\mathfrak{h}$-module also.

By straightforward computations,

$$
\begin{align*}
& {\left[B(0,0,1,0,0), B\left(b_{0}, \ldots, b_{4}\right)\right]} \\
& \quad=-\frac{21 b_{1}}{5} A_{0}-\frac{21 b_{3}}{5} A_{2}+C\left(0,-4 b_{0},-\frac{8 b_{1}}{5}, 0,-\frac{8 b_{3}}{5},-4 b_{4}, 0\right) . \tag{5.8}
\end{align*}
$$

We chose $B(0,0,1,0,0)$ because it is diagonal.
By (5.8), $\left[U_{2}, U_{2}\right] \subset U_{1} \oplus U_{3}$. By choosing $b_{1}=b_{3}=0$ in (5.8), [ $U_{2}, U_{2}$ ] contains a non-zero element of $U_{3}$. This implies [ $U_{2}, U_{2}$ ] contains $U_{3}$. By choosing $b_{1} \neq 0$ in (5.8), $\left[U_{2}, U_{2}\right]$ contains an element of the form $X+X^{\prime}$ where $X \in U_{1}$ is non-zero and $X^{\prime} \in U_{3}$. So $X \in\left[U_{2}, U_{2}\right]$. This implies [ $U_{2}, U_{2}$ ] contains $U_{1}$ also. This proves (1).

Other cases follow from the following relations and by similar arguments. We found these relations manually. However, it can be checked by a routine program in MAPLE (which we did).

$$
\begin{aligned}
& {\left[B(0,0,1,0,0), C\left(c_{0}, \ldots, c_{6}\right)\right]} \\
& \quad=B\left(-\frac{16 c_{1}}{7},-\frac{16 c_{2}}{7}, 0,-\frac{16 c_{4}}{7},-\frac{16 c_{5}}{7}\right) \\
& \quad+D\left(0,-3 c_{0},-\frac{12 c_{1}}{7},-\frac{15 c_{2}}{7}, 0,-\frac{15 c_{4}}{7},-\frac{12 c_{5}}{7},-3 c_{6}, 0\right)
\end{aligned}
$$

$$
\begin{aligned}
& {\left[B(0,0,1,0,0), D\left(d_{0}, \ldots, d_{8}\right)\right]} \\
& \quad=C\left(-3 d_{1},-4 d_{2},-d_{3}, 0,-d_{5},-4 d_{6},-3 d_{7}\right), \\
& \quad\left[C(0,0,0,1,0,0,0), C\left(c_{0}, \ldots, c_{6}\right)\right], \\
& \quad=6 c_{2} A_{0}-6 c_{4} A_{2}+C\left(-c_{0}, c_{1}, c_{2}, 0,-c_{4},-c_{5}, c_{6}\right), \\
& {\left[C(0,0,0,1,0,0,0), D\left(d_{0}, \ldots, d_{8}\right)\right],}
\end{aligned}
$$

$$
=B\left(\frac{20 d_{2}}{7}, \frac{10 d_{3}}{7}, 0,-\frac{10 d_{5}}{7},-\frac{20 d_{6}}{7}\right)
$$

$$
+D\left(2 d_{0},-d_{1},-\frac{13 d_{2}}{7},-\frac{9 d_{3}}{7}, 0, \frac{9 d_{5}}{7}, \frac{13 d_{6}}{7}, d_{7},-2 d_{8}\right)
$$

$$
\begin{aligned}
& {\left[D(0,0,0,0,1,0,0,0,0), D\left(d_{0}, \ldots, d_{8}\right)\right]} \\
& \quad=-14 d_{3} A_{0}-14 d_{5} A_{2}+C\left(-5 d_{1}, 5 d_{2}, 3 d_{3}, 0,3 d_{5}, 5 d_{6},-5 d_{7}\right)
\end{aligned}
$$

$\left[B^{\prime}(0,0,1,0,0), B^{\prime}\left(b_{0}^{\prime}, \ldots, b_{4}^{\prime}\right)\right]$

$$
=-3 b_{1} A_{0}^{\prime}-3 b_{3} A_{2}^{\prime} .
$$

Note that the Lie algebras of $H_{w 1}, H_{w 2}$ are isomorphic to $\mathrm{sl}(2)$. We denote the Lie algebra of $H_{w 3}$ by so(5) (more precisely so(3,2)). Since $\operatorname{dim} \operatorname{so}(5)=10, \operatorname{so}(5)=U_{1} \oplus U_{3}$ by counting the dimension.

Proposition (5.9). If $\mathfrak{h} \subset \mathfrak{f} \subset \mathfrak{h}_{1} \times \mathfrak{h}_{2}$ is a Lie subalgebra, $\mathfrak{f}$ is one of the following subalgebras.

$$
\begin{gathered}
\mathfrak{h}, \mathrm{sl}(2) \times \mathrm{sl}(2), \mathrm{sl}(2) \times \mathrm{sl}(3), \operatorname{so}(5) \times \mathrm{sl}(2), \\
\mathrm{so}(5) \times \mathrm{sl}(3), \mathrm{sl}(5) \times \mathrm{sl}(2), \mathrm{sl}(5) \times \operatorname{sl}(3) .
\end{gathered}
$$

Proof. Let $\mathfrak{f}$ be as above. Note that $\mathrm{sl}(3)$ does not contain any $\mathfrak{h}$-module which is isomorphic to $U_{3}$ or $U_{4}$. Suppose $\mathfrak{j} \supset U_{4}$. Then $\mathfrak{f} \supset U_{1} \oplus U_{3}$ by Lemma (5.7)(1). So $\mathfrak{f} \supset U_{2}$ by Lemma (5.7)(5). Since $\mathfrak{f} \supset U_{1}$, $\mathfrak{f} \supset U_{1} \oplus V_{1}$. Therefore, $\mathfrak{f} \supset \operatorname{sl}(5) \times \operatorname{sl}(2)$. So $\mathfrak{f}=\operatorname{sl}(5) \times \operatorname{sl}(2)$ or $\mathrm{sl}(5) \times \mathrm{sl}(3)$.

Suppose the projection of $\mathfrak{f}$ to the first factor contains $U_{2}$. Then there exists an $\mathfrak{h}$-homomorphism $\alpha: U_{2} \rightarrow V_{2}$ such that $(x, \alpha(x)) \in \mathfrak{f}$ for all $x \in U_{2}$. By Lemma (5.7)(1), the projection of $\mathfrak{f}$ to the first factor contains $U_{3}$. Since $U_{3}$ is not equivalent to any other factor, $f$ contains $U_{3}$. By Lemma (5.7)(4), $\mathfrak{f} \supset U_{1}$. Since $\mathfrak{f} \supset \mathfrak{h}, \mathfrak{f} \supset U_{1} \oplus V_{1}$. If $x \in U_{2}, y \in U_{1},(y, 0)$ $\in \mathfrak{f}$. So $[(y, 0),(x, \alpha(x))]=([y, x], 0) \in \mathfrak{f}$. Since $\left[U_{1}, U_{2}\right]=U_{2}, \mathfrak{f} \supset U_{2}$. By Lemma (5.7)(2), $\mathfrak{f} \supset U_{4}$ and it reduces to the previous case.

Suppose $\mathfrak{f}$ does not contain $U_{4}$ and the projection to the first factor does not contain $U_{2}$. Suppose $\mathfrak{f} \supset U_{3}$. By Lemma (5.7)(4), $\mathfrak{f} \supset U_{1}$. Therefore, $\mathfrak{f}$ has so(5) as the first factor. This implies $\mathfrak{f}=\operatorname{so}(5) \times \mathrm{sl}(2)$ or $\mathrm{so}(5) \times \mathrm{sl}(3)$.

Suppose the projection of $\mathfrak{f}$ to the first factor is $U_{1}$. If $\mathfrak{f} \supset V_{2}, \mathfrak{f} \supset V_{1}$ also. Therefore, $\mathfrak{f}=\operatorname{sl}(2) \times \operatorname{sl}(3)$. Otherwise the projection of $\mathfrak{f}$ to both factors are $\mathrm{sl}(2)$. Since there is no sub $\mathfrak{h}$-module between $\mathfrak{h}$ and $U_{1} \times V_{1}, \mathfrak{f}$ is $\mathfrak{h}$ or $\operatorname{sl}(2) \times \operatorname{sl}(2)$.

Now we specialize to the field $k=\mathbb{R}$.

Proposition (5.10). Let $x \in G_{\mathbb{R}} w$ and $H_{x \mathbb{R}} \subset F \subset H_{\mathbb{R}}$ be a closed connected subgroup. Then $F$ is one of the following subgroups.

$$
\begin{gathered}
H_{x \mathbb{R}}, H_{x 1 \mathbb{R}} \times H_{x 2 \mathbb{R}}, H_{x 1 \mathbb{R}} \times \mathrm{SL}(3)_{\mathbb{R}}, \\
H_{x 3 \mathbb{R}} \times H_{x 2 \mathbb{R}}, H_{x 3 \mathbb{R}} \times \mathrm{SL}(3)_{\mathbb{R}}, \mathrm{SL}(5)_{\mathbb{R}} \times H_{x 2 \mathbb{R}}, \mathrm{SL}(5)_{\mathbb{R}} \times \mathrm{SL}(3)_{\mathbb{R}} .
\end{gathered}
$$

Proof. If $x=g w$ for $g \in G_{\mathbb{R}}, H_{x \mathbb{R}}=g H_{w \mathbb{R}} g^{-1}$, etc. So we may assume that $x=w$. Then this proposition follows from the previous proposition.

## 6. AN ANALOGUE OF THE OPPENHEIM CONJECTURE

In this section, we prove an analogue of the Oppenheim conjecture.
In the following lemma, $x \in V_{\mathbb{C}}^{\mathrm{ss}}$. We define $H_{x 1 \mathbb{C}}$, etc. as in Definition (5.1).
Lemma (6.1). (1) If $y \in V_{\mathbb{C}}$ is fixed by $H_{x \mathbb{C}}, y$ is a scalar multiple of $x$.
(2) If $y \in \operatorname{Sym}^{2} V_{2}^{*}$ is fixed by $H_{x 2 \mathbb{C}}, y$ is a scalar multiple of $\Phi_{1}(x)$.
(3) If $y \in \operatorname{Sym}^{2} V_{1}^{*}$ is fixed by $H_{x 1 \mathbb{C}}$ or $H_{x 3 \mathbb{C}}, y$ is a scalar multiple of $\Phi_{2}(x)$.

Proof. Consider (1). Let $x=g w$ with $g \in G_{\mathbb{C}}$. Then $H_{x \subseteq}=g H_{w \subseteq} g^{-1}$, and $g^{-1} y$ is fixed by $H_{w \subset}$. So we may assume $x=w$. By (3.1), $V$ contains the trivial representation of $\mathfrak{h}$ precisely once. Therefore, the set of fixed points of $H_{w \mathbb{C}}$ is of dimension one. This proves (1).

Consider the first part of (3). As in (1), we may assume $x=w$. Since $H_{w 1 \mathbb{C}} \cong \operatorname{PGL}(2)_{\mathbb{C}}$, it is enough to show that $\operatorname{Sym}^{2} V_{2 \mathbb{C}}^{*}$ contains the trivial representation of the Lie algebra $\mathfrak{h}_{\mathbb{C}}$ of PGL(2) $\mathbb{C}_{\mathbb{C}}$ precisely once. Let $\Lambda$ be the fundamental dominant weight of $\mathfrak{h}$ as before. Since $V_{1} \cong V_{1}^{*} \cong 4 \Lambda$, by considering weights, it is easy to see that

$$
\operatorname{Sym}^{2} V_{2 \mathbb{C}}^{*} \cong(8 \Lambda)_{\mathbb{C}} \oplus(4 \Lambda)_{\mathbb{C}} \oplus \mathbb{C}
$$

The second part of (3) and (4) are well known and were used in the proof of the Oppenheim conjecture for quadratic forms.

In the following theorem, let $x \in G_{\mathbb{R}} w$. Then $H_{x \mathbb{R}+}^{\circ}$ is generated by unipotent elements. Let $H_{x \mathbb{R}} \subset F \subset H_{\mathbb{R}}$ be the closed connected subgroup such that $\overline{H_{x \mathbb{R}} H_{\mathbb{Z}}}=F H_{\mathbb{Z}}$. By Ratner's theorem (Theorem (0.2)), such $F$ exists.

Theorem (6.2). (1) If $\Phi_{2}(x) \notin \mathbb{P}\left(\operatorname{Sym}^{2} V_{1}^{*}\right)_{\mathbb{Q}}, F=\operatorname{SL}(5)_{\mathbb{R}} \times H_{x 2 \mathbb{R}}$ or $F=\operatorname{SL}(5)_{\mathbb{R}} \times \operatorname{SL}(3)_{\mathbb{R}}$.
(2) If $\Phi_{1}(x) \notin \mathbb{P}\left(\operatorname{Sym}^{2} V_{2}^{*}\right)_{\mathbb{Q}}$ and $\Phi_{2}(x) \notin \mathbb{P}\left(\operatorname{Sym}^{2} V_{1}^{*}\right)_{\mathbb{Q}}, F=\operatorname{SL}(5)_{\mathbb{R}}$ $\times \operatorname{SL}(3)_{\mathbb{R}}$.

Proof. Suppose $F=H_{x \mathbb{R}}$. Then $F$ is defined over $\mathbb{Q}$. Therefore, for any $\sigma \in \operatorname{Aut}(\mathbb{C} / \mathbb{Q}), H_{x \mathbb{C}}^{\sigma}=H_{x \mathbb{C}}$. Since $H_{x \mathbb{C}}^{\sigma}=H_{x^{\sigma} \mathbb{C}}, x^{\sigma}$ is fixed by $H_{x \mathbb{C}}$. So $x^{\sigma}$ is a scalar multiple of $x$ by Lemma (6.1). Since this is the case for all $\sigma$, $[x] \in \mathbb{P}(V)_{\mathbb{Q}}$. Since $\Phi_{1}, \Phi_{2}$ are defined over $\mathbb{Q}, \Phi_{1}(x), \Phi_{2}(x)$ are $\mathbb{Q}$-rational points.

We show that $\Phi_{2}(x) \in \mathbb{P}\left(\operatorname{Sym}^{2} V_{1}^{*}\right)_{\mathbb{Q}}$ if $F=H_{x 1 \mathbb{R}} \times H_{x 2 \mathbb{R}}$ or $H_{x 1 \mathbb{R}} \times$ $\mathrm{SL}(3)_{\mathbb{R}}$. Since the argument is similar, we only consider the first case.

For any $\sigma \in \operatorname{Aut}(\mathbb{C} / \mathbb{Q})$,

$$
H_{x 1 \mathbb{C}}^{\sigma} \times H_{x 2 \mathbb{C}}^{\sigma}=H_{x^{\sigma} 1 \mathbb{C}} \times H_{x^{\sigma} 2 \mathbb{C}}=H_{x 1 \mathbb{C}} \times H_{x 2 \mathbb{C}} .
$$

Since $H_{x \mathbb{C}} \subset H_{x 3 \mathbb{C}}, \Phi_{2}\left(x^{\sigma}\right)$ is fixed by $H_{x 1 \mathbb{C}}$. This implies $\Phi_{2}\left(x^{\sigma}\right)=\Phi_{2}(x)$. Since $\Phi_{2}$ is defined over $\mathbb{Q}, \Phi_{2}\left(x^{\sigma}\right)=\Phi_{2}(x)^{\sigma}=\Phi_{2}(x)$ by Lemma (6.1). Therefore, $\Phi_{2}(x) \in \mathbb{P}\left(\operatorname{Sym}^{2} V_{1}^{*}\right)_{\mathbb{Q}}$.

By a similar argument, if $F=H_{x 1 \mathbb{R}} \times H_{x 2 \mathbb{R}}, H_{3 \mathbb{R}} \times H_{x 2 \mathbb{R}}$, or $\operatorname{SL}(5)_{\mathbb{R}} \times$ $H_{x 2 \mathbb{R}},\left[\Phi_{1}(x)\right] \in \mathbb{P}\left(\operatorname{Sym}^{2} V_{2}^{*}\right)_{\mathbb{Q}}$. Also if $F=H_{x 3 \mathbb{R}} \times H_{x 2 \mathbb{R}}$ or $H_{x 3 \mathbb{R}} \times \operatorname{SL}(3)_{\mathbb{R}}$, $\left[\Phi_{2}(x)\right] \in \mathbb{P}\left(\operatorname{Sym}^{2} V_{1}^{*}\right)_{\mathbb{Q}}$.

By these considerations, conditions in (1), (2) force $F$ to become the given subgroups.

Lemma (6.3). Let $x \in G_{\mathbb{R}} w$. Then for any non-zero real number $r$, there exists $h \in H_{\mathbb{R}}$ and a primitive integer point $a \in V_{1 \mathbb{Z}}$ such that $F_{h^{-1} x}(a)=r$.

Proof. We may assume $x=\lambda w$ where $\lambda \in \mathbb{R} \backslash\{0\}$. Since $F_{\lambda h^{-1} w}(a)$ $=\lambda^{60} F_{w}(h a)$, the above condition is equivalent to $F_{w}(h a)=\lambda^{-60} r$. Put $t=$ $-\lambda^{-20}(r / 2)^{1 / 3}$. Then

$$
h=\left(\left(\begin{array}{ccccc}
t^{-1} & 0 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 & 0 \\
0 & 0 & t & 0 & 0 \\
0 & 0 & 0 & 1 & 0 \\
0 & 0 & 0 & 0 & 1
\end{array}\right), I_{3}\right)
$$

and $a=^{t}(00100)$ satisfy the condition.
In the following theorem, $x \in G_{\mathbb{R}} w$.
Theorem (6.4). If $\left[\Phi_{2}(x)\right] \notin \mathbb{P}\left(\operatorname{Sym}^{2} V_{1}^{*}\right)_{\mathbb{Q}}$, the set of values of the cubic form $F_{x}(a)$ at primitive integer points is dense in $\mathbb{R}$.

Proof. Let $r$ be a non-zero real number. We choose $h=\left(h^{\prime}, h^{\prime \prime}\right) \in H_{\mathbb{R}}$ and $a \in V_{1 \mathbb{Z}}$ as in Lemma (6.3). By Theorem (6.2), there exist $h_{1}=\left(h_{1}^{\prime}, h_{1}^{\prime \prime}\right)$ $\in H_{x \mathbb{R}}$ and $h_{2}=\left(h_{2}^{\prime}, h_{2}^{\prime \prime}\right) \in H_{\mathbb{Z}}$ such that $h_{1}^{\prime} h_{2}^{\prime}$ is close to $h^{\prime}$. Then

$$
F_{x}\left(h_{2}^{\prime} a\right)=F_{h_{2}^{-1} x}(a)=F_{h_{2}^{-1} h_{1}^{-1} x}(a)=F_{\left(h_{2}^{\prime-1} h_{1}^{\prime}, 1\right) x}(a)
$$

is close to

$$
F_{\left(h^{\prime-1}, 1\right) x}(a)=F_{h^{-1} x}(a)=r .
$$

Note that $F_{h^{-1} x}$ does not depend on the second component of $h$.
Since $h_{2}^{\prime} a \in V_{1 \mathbb{Z}}$ is primitive, this proves the theorem.

Note that if $x=g w$ with $g=\left(g_{1}, g_{2}\right) \in G_{\mathbb{R}}$,

$$
\left[\Phi_{2}(x)\right]=g_{1}\left[\Phi_{2}(w)\right], F_{x}(a)=\left(\operatorname{det} g_{1}\right)^{15}\left(\operatorname{det} g_{2}\right)^{3} F_{w}\left(g_{1}^{-1} a\right) .
$$

Therefore, writing down $F_{w}$, etc. explicitly, we get the statement of Theorem (0.3).

Remark (6.5). We proved Theorem (6.4) as a consequence of Theorem (6.2). But we don't need our prehomogeneous vector space if we just want to prove Theorem (6.4). For that purpose, we only have to consider the situation $\mathrm{PGL}(2) \subset \mathrm{SL}(5)$ and apply Ratner's theorem using the computations in Section 5. We discuss this issue in [28].

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[^0]:    * Partially supported by NSF Grant DMS-9401391.

