# The Solution to a Generalized Toda Lattice and Representation Theory 

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## Contents

## 0. Introduction

0.1. Statement of the main results
0.2. A brief description of Sections 1-7
0.3. General comments

1. Poisson commutativity of translated invariants
1.1. A formula for the Poisson bracket
1.2. The case of a subalgebra of a semi-simple Lie algebra
1.3. Relating cent $g^{x}$ to $S(g)^{G}$
1.4. The theorem $\left[I^{\prime}, J\right]=0$ for Lie summands
1.5. The case of a Borel subalgebra where $f=\sum e_{-\alpha_{i}}$
1.6. The commutativity of certain vector fields on $a_{f}^{*}$
2. The variety $Z$ of normalized Jacobi elements
2.1. The variety $b_{f}=f+b$ and invariants $I \in S(\mathcal{g})^{G}$
2.2. The vector fields $\xi_{I}, I \in S(g)^{G}$, on $Z$
2.3. The subvarieties $Z(\gamma) \subseteq Z$ and $N$-conjugacy
2.4. The Bruhat-Gelfand decomposition, the isomorphism $\beta_{v}: G_{*}^{y} \rightarrow Z(\gamma)$, and the rationality of $Z(\gamma)$ in the complex case
2.5. The foliation $Z=\cup Z(\gamma)$ in the complex case
2.6. The isomorphism $\beta_{(w):} G_{(*)}^{*} \rightarrow Z(\gamma)$ in the complex case
3. The parametrization $Z \cong H \times \boldsymbol{n}_{+}$in the real case
3.1. The polar decomposition and centralizers $g^{x}$
3.2. The relation $G^{v} \cap \bar{N} H N=G_{e^{y}}^{y}$ for $y \in Z$ and the isomorphism $Z(\gamma) \cong \mathbb{R}^{t}$
3.3. The open Weyl chamber $\hbar_{+}$and $Z$
3.4. The normalizer $\tilde{G}$ of $\mathscr{g}$ in $\mathrm{Ad} \mathscr{g} \mathbb{C}$ and the subset $G_{(*)}=s(\kappa) G_{*} \subseteq \tilde{G}$
3.5. The element $\bar{n}_{f}(w)$ and the key relation $\tilde{G}^{w} \cap G_{(*)}=G_{0}{ }^{w}$
3.6. The isomorphism $\beta_{(w)}: G_{0}{ }^{w} \rightarrow Z(\gamma)$
3.7. The isomorphism $H \times h_{+} \rightarrow Z,\left(g_{0}, w_{0}\right) \mapsto n(g) w$
[^0]
## 4. The isospectral leaf $Z(\gamma)$ as a complete flat affinely connected manifold

4.1. The union $Z=U Z(\gamma)$ as a foliation
4.2. The flat affine connection in $Z(\gamma)$
4.3. The maps $\beta_{(w)}$ and $\beta_{y}$ as isomorphisms of complete affinely connected manifolds
5. Representations and the functions $\Phi_{\lambda}\left(g_{o}, w_{o} ; t\right)$
5.1. The module $V^{\lambda}, \lambda \in D$, as a real Hilbert space
5.2. The normalization of the lowest weight vector $\varepsilon^{\kappa \lambda}=s_{0}(\kappa) v^{\lambda}$
5.3. The equation $\bar{n}(a) h(a) n(a)=s_{0}(\kappa)^{-1} \bar{n}_{-f}(w) a_{0} \bar{n}_{f}(w)^{-1}$
5.4. The polar decomposition of $h(g)^{1 / 2} n(g) \bar{n}_{f}(w)$
5.5. Formulas for $d(w)=g_{0}^{-1} p\left(g_{o}, w_{o}\right)$ in terms of $\bar{n}_{f}(w)$, as a limit, and its independence of $g_{0}$. Also formula for $g_{0}{ }^{\lambda}$ in terms of $n(g)$ and sum formula for $h(g \exp t w)^{\lambda}$
5.6. Relation to Moser's treatment of the Toda lattice in [19]
5.7. Sums over $\mathscr{P}$ and the completion of $U(n)$ and $U(\bar{n})$
5.8. The polynomials $p\left(s, w_{o}\right)$ and formulas for $\bar{n}_{f}(w), \bar{n}_{f}(w)^{-1}$, and $\bar{n}_{-f}(w)$
5.9. The numbers $c_{8, \lambda}$ and formulas for $d(w)^{\lambda}$ and $d(w)^{-\kappa \lambda}$
5.10. The definition of $\Phi_{\lambda}\left(g_{a}, w_{n} ; t\right)$ and the result $h(g \exp (-t) w)^{\lambda}=\Phi_{\lambda}\left(g_{n}, w_{a} ; t\right)$
5.11. The limiting behavior of $\log h(g \exp (-t) w)^{\mu}$ as $t \rightarrow \pm \infty$
6. The symplectic structure of $\left(Z, \omega_{Z}\right)$ and the integration of $\xi_{I}, I \in S(g)^{G}$
6.1. The symplectic structure computed in terms of Poisson brackets
6.2. The symplectic structure on coadjoint orbits
6.3. The coadjoint theory transferred to $a^{*}$
6.4. The symplectic manifold $\left(Z, \omega_{Z}\right)$ as a translated coadjoint orbit
6.5. Formula for $I_{1} \mid Z$ in terms of canonical coordinates
6.6. Complete real polarizations $F$ and integration of $\xi_{\varphi}, \dot{\varphi} \in C_{F}^{\infty}$
6.7. The formula $\beta_{(w)}(g \exp (-t) z)=\exp t \xi_{I} \cdot y$ and the formila $\varphi_{i}\left(\exp t \xi_{I} \cdot y\right)=$ $\log h(g \exp (-t) z)^{-\alpha_{i}}$
6.8. The special case where $I=I_{1}=Q / 2$
7. Denouement; the formula for $q_{i}(x(t))$
7.1. Hamilton's equations, the Hamiltonian $H=\sum p_{i}{ }^{2} / 2 m_{i}+r_{1} e^{\psi_{1}}+\cdots+r_{l} e^{\psi_{l}}$, and the bilinear form $B_{H}$
7.2. Reducing the integration of $\xi_{H}$ to the integration of $\xi_{I^{\prime}}$
7.3. The fixing of $Q$
7.4. Examples of where the $\psi_{i}$ define a Dynkin diagram and the tie-up with $\left(Z, \omega_{Z}\right)$
7.5. The main result; the formula for $q_{i}(x(t))$
7.6. The scattering of the generalized Toda lattice in terms of $g_{o}$ and $d(w)$
7.7. Application to the standard Toda lattice and the recovery of Moser's result on its scattering, $\lim _{t \rightarrow \pm \infty} q_{i}(x(t))$
7.8. The three-body problem $\sum_{i=1}^{a} p_{i}^{2} / 2 m_{i}+e^{q_{1}-q_{2}}+e^{q_{2}-q_{3}}+e^{2 q_{3}}$ worked out

## 0. Introduction

0.1. One of the main results of this paper concerns the existence of certain classical mechanical systems, generalizing the (finite, nonperiodic) Toda lattice (see, e.g., $[8,10,19,24,25]$ ) which one can explicitly integrate for all values of time. These systems are related to Dynkin diagrams. One
knows these diagrams classify semi-simple Lie groups. The integration of the system associated to a Dynkin diagram is expressible in terms of the finitedimensional representations of the corresponding group. In fact, more than that is true. The integration of the system and the theory of the finite-dimensional representation theory of semi-simple Lie groups are in a sense equivalent. Indeed, the integration of the system completely determines and is determined by the weight structure of the fundamental representations of the corresponding group.
In more detail, recall the Hamilton-Jacobi theory as it would apply to a system of particles moving on a line with respect to some potential. We envision then $n$ particles, say with masses $m_{i}, i=1, \ldots, n$, and $\mathbb{R}^{2 n}$ with linear canonical coordinate functions $p_{i}, q_{j}, i, j=1, \ldots, n$ (i.e., phase space), so that any $x \in \mathbb{R}^{2 n}$ is a (classical) state of the system and $p_{i}(x), q_{i}(x)$ are, respectively, the momentum and position of the $i$ th particle in the state $x$. Assume the potential $U$ is smooth (i.e., $U \in C^{\infty}\left(\mathbb{R}^{2 n}\right)$ ), and of course, depends only on the $q$ 's (i.e., $\left(\partial / \partial p_{i}\right) U=0$, $i=1,2, \ldots, n)$. The Hamiltonian $H \in C^{\alpha}\left(\mathbb{R}^{2 n}\right)$ of the system is then given by

$$
\begin{equation*}
H=\sum_{i=1}^{n} \frac{p_{i}{ }^{2}}{2 m_{i}}+U . \tag{0.1.1}
\end{equation*}
$$

Now if $x \in \mathbb{R}^{2 n}$ and $x(t) \in \mathbb{R}^{2 n}$ is the state of the system at time $t$ such that $x(0)=x$ then $x(t)$ is determined by Hamilton's equations. That is,

$$
\begin{equation*}
-\frac{\partial H}{\partial q_{j}}(x(t))=\frac{d p_{j}(x(t))}{d t}, \quad p_{j}(x(t))=m_{j} \frac{d q_{j}(x(t))}{d t} . \tag{0.1.2}
\end{equation*}
$$

The fundamental problem of the Hamilton-Jacobi theory is to integrate (0.1.2), that is, to determine the function $t \rightarrow x(t)$. It is clear from the second equation in (0.1.2) that it is enough to determine the position functions $t \rightarrow q_{j}(x(t))$, $j=1, \ldots, n$. For a general potential function $U$ of course one cannot do this. It is one of the main results of this paper to show that for certain very special $U$ having to do with Dynkin diagrams one can explicitly integrate (0.1.2) using representation theory.

Let $\mathscr{Q} \subseteq C^{\infty}\left(\mathbb{R}^{2 n}\right)$ be the $n$-dimensional vector space spanned by $q_{j}, j=$ $1, \ldots, n$. Let $l \leqslant n$ and let $\psi_{i} \in \mathscr{Q}, i=1, \ldots, l$, be $l$ linear independent functions. Thus we may write

$$
\begin{equation*}
\psi_{i}=\sum_{j=1}^{n} a_{i} q_{j}, \quad i=1, \ldots, l, \tag{0.1.3}
\end{equation*}
$$

where ( $a_{i j}$ ) is a constant $l \times n$ matrix. The potentials we shall consider are those of the form

$$
\begin{equation*}
U=r_{1} e^{{L_{1}}_{1}}+\cdots+r_{e^{\prime}} e^{\nu_{l}}, \tag{0.1.4}
\end{equation*}
$$

where the $r_{i}$ are some positive constants. For example, if $l=n-1, m_{j}=$ $r_{i}=1$ and $\psi_{i}=q_{i}-q_{i+1}, i=1, \ldots, l$, then the system is the Toda lattice. The interaction between the particles is that of exponential forces between nearest neighbors. (Throughout this paper "Toda lattice" will mean "finite nonperiodic Toda lattice".)

What is crucial for us in our result here is the geometry of the $\psi$ 's, that is, the lengths of the vectors $\psi_{i} \in \mathscr{Q}$ and the angles between them. We hasten to add that there is in fact a natural geometry in $\mathscr{Q}$, defined by the Hamiltonian $H$. Indeed let $[\varphi, \psi] \in C^{\infty}\left(\mathbb{R}^{2 n}\right)$ be the Poisson bracket for any pair of functions $\varphi, \psi \in C^{\infty}\left(\mathbb{R}^{2 n}\right)$. It is easy to see that if $\varphi, \psi \in \mathscr{Q}$ then $[\varphi,[\psi, H]]$ is a scalar times the identity function 1 . A positive definite symmetric bilinear form $B_{H}$ is then defined in $\mathscr{Q}$ by the relation

$$
\begin{equation*}
B_{H}(\varphi, \psi) 1=[\varphi,[\psi, H]] . \tag{0.1.5}
\end{equation*}
$$

One has in fact

$$
\begin{equation*}
B_{H}\left(q_{i}, q_{j}\right)=\delta_{i j} / m_{i} \tag{0.1.6}
\end{equation*}
$$

We will say that the $\psi$ 's define a Dynkin diagram if there exists a real split semi-simple Lie algebra $g$ of rank $l$ with an invariant bilinear form $Q$, and with a split Cartan subalgebra $h$ with simple roots $\alpha_{1}, \ldots, \alpha_{l}$ such that

$$
\begin{equation*}
B_{H}\left(\psi_{i}, \psi_{j}\right)=Q\left(\alpha_{i}, \alpha_{j}\right) \tag{0.1.7}
\end{equation*}
$$

for $i, j=1, \ldots, l$. In such a case, as one knows, the possible angles between the $\psi$ 's are $90,120,135$, or $150^{\circ}$. The Toda lattice satisfies this condition. Here $g$ is the Lie algebra of $\operatorname{Sl}(n, \mathbb{R})$. The angle between $q_{i-1}-q_{i}$ and $q_{i}-q_{i+1}$ is $120^{\circ}$. We will give some further examples. (See Section 7.4 here and [3;7, Section 30] for others. See also [3] for suggested physical interpretations.)

$$
\begin{equation*}
H=\sum_{i=1}^{n} \frac{p_{i}^{2}}{2}+e^{a_{1}-q_{2}}+\cdots+e^{q_{n-1}-q_{n}}+e^{q_{n}} \tag{1}
\end{equation*}
$$

Example (1) is like the Toda lattice except that the last particle also interacts with a fixed mass. The diagram here is that of $B_{n}$; i.e., $g$ is isomorphic to the Lie algebra of $S 0(n, n+1)$.
(2) $H=\sum_{i=1}^{4} \frac{p_{i}{ }^{2}}{2}+e^{q_{1}-q_{2}}+e^{q_{2}-q_{3}}+e^{q_{3}}+e^{\left(a_{4}-q_{1}-q_{2}-q_{3}\right) / 2}$.

This is a four-body problem where the first three particles are as in (1) but where, also, the center of mass of these three particles interacts exponentially
with a fourth particle. The Dynkin diagram here is that of the exceptional Lie algebra $F_{4}$.

Of course many different systems may correspond to the same diagram. For example, the following two four-body problems both correspond to $D_{4}$.
(3) $\sum_{i=1}^{4} \frac{p_{i}{ }^{2}}{2}+e^{q_{1}}+e^{q_{2}}+e^{q_{3}}+e^{\left(q_{4}-q_{1}-q_{2}-q_{3} / 2\right.}$.
(4) $\sum_{i=1}^{4} \frac{p_{i}{ }^{2}}{2}+e^{q_{1}-q_{2}}+e^{q_{2}-q_{3}}+e^{q_{3}-q_{4}}+e^{q_{3}+q_{4}}$.

Example (4) is like the usual Toda lattice with four particles except that the center of mass of the third and fourth particles interacts exponentially with a fixed mass. In Section 7.8 we will apply the result Theorem 7.5 to work out $q_{j}(x(t))$ for the three-body problem whose potential is $e^{q_{1}-q_{2}}+e^{q_{\mathrm{a}}-q_{3}}+e^{2 q_{3}}$. The Lie algebra in question is that of the symplectic group $\operatorname{Sp}(6, \mathbb{R})$.

Before describing the main result we would like to remark that although the condition on the $\psi$ 's is quite rigid there is in fact some flexibility. Namely, by adjusting the masses $m_{i}$ one has, given a potential of form (0.1.4), an $n$ parameter degree of freedom (see (0.1.6)), which one may use to possibly satisfy the condition.

Assume now that the $\psi$ 's define a Dynkin diagram and that $g, Q, h$, and the $\alpha$ 's are as above. We will occasionally then refer to the mechanical system whose potential is (0.1.4) as a generalized Toda lattice. Let $\pi_{i}, i=1, \ldots, l$, be the $l$-fundamental finite-dimensional irreducible representation of $g$ (or $g_{\mathrm{C}}$, its complexification). Paraphrasing Theorem 7.5 one has the following solution to Hamilton's equation for the generalized Toda lattice.

Theorem A. If $m(i)$ is the number of distinct weights of $\pi_{i}$ there exist constants $c_{k}$ and $d_{k}, k=1, \ldots, m(i)$, which depend upon the initial state $x \in \mathbb{R}^{2 n}$ of the system, whose potential is (0.1.4), and on the weights of $\pi_{i}$ (see Theorem 7.5 for the exact dependence) such that if

$$
\begin{equation*}
\Phi_{i}(t)=c_{1} e^{-t d_{1}}+c_{2} e^{-t d_{2}}+\cdots+c_{m(i)} e^{-t d_{m(i)}} \tag{0.1.8}
\end{equation*}
$$

then $\Phi_{i}(t)>0$ and except for a linear term (which vanishes if $l=n$ and all $r_{k}=1$ ) one has

$$
\begin{equation*}
q_{j}(x(t))=-\frac{2}{m_{j}} \sum_{i=1}^{l} \frac{a_{i j}}{Q\left(\alpha_{i}, \alpha_{i}\right)} \log \Phi_{i}(t) \tag{0.1.9}
\end{equation*}
$$

for $j=1, \ldots, n$. The $a_{i j}$ are given by (0.1.3).
Remark 0.1. A solution to the generalized Toda lattice has appeared in [28]. This solution (see Corollary 1 in [28]) is stated in terms of the diagonal
component $h(\exp t y)$ of the exponential of a symmetric Jacobi element $y$. The reduction of the problem to $h(\exp t y)$ is also an early side step in the proof of Theorem 7.5. It is implicit in Theorem 6.8.2 here. The reduction of the problem to $h(\exp t y)$ has been known and widely discussed by the author for a number of years. The main problem lies not with the reduction to $h(\exp t y)$ -this is relatively easy-but with the determination of $h(\exp t y)$. This is where representation theory enters. Also entering here is an analysis of the variety $Z$ (a coadjoint orbit) of normalized Jacobi elements. See Sections 2-6.1

Fixing $i$ the ordering of the weight may be made so that $c_{1}, d_{1}$ corresponds to the highest weight of $\pi_{i}$ and $c_{m(i)}, d_{m(i)}$ correspond to the lowest weight. One then has

$$
\begin{equation*}
d_{1}>d_{k}>d_{m(i)} \tag{0.1.10}
\end{equation*}
$$

for all $1>k>m(i)$. Write $d_{j i}=d_{j}$ and $c_{j i}=c_{j}$. Now if $b_{i}, i=1, \ldots, l$, are $l$ numbers and

$$
s_{j}=\frac{2}{m_{j}} \sum_{i=1}^{l} \frac{a_{i j}}{Q\left(\alpha_{i}, \alpha_{i}\right)} b_{i}
$$

then (1) put $v_{j}^{+}=s_{j}$ in case $b_{i}=d_{m(i) i}$, (2) put $v_{j}^{-}=s_{j}$ in case $b_{i}=d_{1 i}$, (3) put $u_{j}^{+}=-s_{j}$ in case $b_{i}=\log c_{m(i) i}$, and (4) put $u_{j}^{-}=-s_{j}$ in case $b_{i}=$ $\log c_{1 i}$. (One knows that $c_{m(i) i}$ and $c_{1 i}$ are positive.)

We can now express the scattering of the system as $t \rightarrow \pm \infty$ in terms of the highest and lowest weights of the fundamental representations. Ignoring a trivial linear term as in Theorem A, the following paraphrases Theorem 7.6.

Theorem B. For any $j=1, \ldots, n$ the curve $\left(t, q_{j}(x(t))\right)$ in the $t, q_{j}$ plane is asymptotic to the line $v_{j}^{-} t+u_{j}^{-}$as $t \rightarrow-\infty$ and to the line $v_{j}^{+} t+u_{j}^{+}$as $t \rightarrow+\infty$. That is, $v_{j}^{ \pm}$are the limiting velocities of the $i$ th particle as $t \rightarrow \pm \infty$ and $u_{j}{ }^{ \pm}$are the limiting phases.

Moser's result (see 4.3 in [19]) on the scattering of the Toda lattice may be readily recovered from Theorem B. See Section 7.7.

We also remark that we can reverse our considerations here. That is, if solution (0.1.9) of Hamilton's equations is given one solves for $\Phi_{i}(t)$ and, over all initial conditions, one obtains weights of the fundamental representation $\pi_{i}$. In particular one obtains the highest weight, which of course determines $\pi_{i}$.
0.2. The mechanical system in Theorem $A$ is completely integrable in a classical sense. This means the Poisson commutativity of a suitable family of functions. In Section 1 we establish the Poisson commutativity in a much more general setting than that used to establish Theorem A. In particular

[^1]one has an "isospectral flow" for all suitably translated coadjoint orbits of a parabolic subalgebra.

The technique for dealing with the usual Toda lattice involves the set of Jacobi matrices. In particular it involves a decomposition of this set into isospectral classes. See, e.g., [19, 25]. In Section 2, first of all, the notion of Jacobi matrix is generalized and one introduces, using the simple negative and positive root vectors, the notion of Jacobi element in a real split or complex semi-simple Lie algebra. One is readily reduced to the study of the $2 l$-dimensional variety $Z$ of normalized Jacobi elements. Section 2 is devoted to an analysis of $Z$, especially in the complex case. Moser used the theory of continued fractions and a method going back to Stieljes to obtain a rationality result for the case of diagonalizable Jacobi matrices. Among other structure theorems this result is generalized in Section 2. The isomorphisms $\beta_{y}$ and $\beta_{(w)}$ establish, in explicit terms, the rationality of the isospectral leaves $Z(\gamma)$. Since we can take the field to be $\mathbb{C}$ in Section 2 the elements in $Z(\gamma)$ can even be (principal) nilpotent. The results are, we believe, new even in the matrix case, and lead, for example, to a simple iterative procedure for diagonalizing a diagonalizable Jacobi matrix once one knows its eigenvalues. The results in Section 2 make notable use of the results in Section 1 of [17].

The situation in the real case is considerably more subtle. Using an interplay between the polar and Bruhat-Gelfand decompositions we obtain in Section 3 a completeness result for $Z(\gamma)$ (not true over $\mathbb{C}$ ) which later guarantees the integration of certain Hamiltonian vector fields for all values of $t$. Using an isomorphism $\beta_{(w)}: G_{0}{ }^{w} \rightarrow Z(\gamma)$ one obtains (later to be seen) the action angle coordinates in terms of the characters of a split Cartan subgroup. In fact $\beta_{(w)}$ leads to a parametrization $Z \cong H \times h_{+}$where $H \cong \mathbb{R}^{l}$ is the split Cartan subgroup and $h_{+}$is an open Weyl chamber. The isolation of $G_{0}{ }^{w}$ (which is not the identity component) among all the connected components of the centralizer $\boldsymbol{G}^{\boldsymbol{w}}$ of $w$ appears to be rather remarkable. It brings into focus a special automorphism, denoted by $m$ in Section 3.5, of $g$. This automorphism also plays a special role in the Whittaker theory (see, e.g., the element $h_{o}$ in [9, p. 106]), to which this paper will be related elsewhere.

In Section 4 it is shown that $\beta_{(w)}: G_{0}{ }^{w} \rightarrow Z(\gamma)$ is an isomorphism of complete, flat, affinely connected manifolds. This sets a correspondence between the cosets of one-parameter subgroups in $G_{0}{ }^{w}$ and the trajectories of the (Hamiltonian) vector fields $\xi_{I}$ in $Z(\gamma)$. It reduces our integration problem to the determination of $h(g \exp (-t) w)^{\lambda}$, as will be seen in Section 6.

In Section 5 we determine $h(g \exp (-t) w)^{\lambda}$ using the representation theory of $g$. A particular role is played by the explicit formulas for $\bar{n}_{f}(w)^{-1}$ and $\bar{n}_{-f}(w)$. It is also seen in Section 5 that the scattering theory is determined by the special element $d(w) \in H$. The relationship between the isospectral leaves $Z(\gamma)$ and the elements $d(w)$ is rather mysterious (to us) and (we believe) remarkable.

Section 6 is devoted to the underlying symplectic theory and coadjoint
orbits. It is shown that $Z$ itself is just one coadjoint orbit (suitably translated) of the Borel subgroup $\bar{B}$. The basic formula $\varphi_{i}(\exp t \xi \cdot y)=\log h(g \exp (-t) w)^{-\alpha_{i}}$ is established here. See (6.8.4).

Section 7 starts off with the Hamilton-Jacobi theory. It can be initially read without reference to the other sections. However, the results of the preceding sections are later brought together yielding Theorem 7.5, which gives the formula for $q_{i}(x(t))$. What is surprising here is the special role of the fundamental representations among all the others. The Cartan matrix initially enters the formula for $q_{i}(x(t))$ but since we are using the fundamental representations it is canceled out, leaving the coefficients $a_{i j}$ in (0.1.9) the same as the coefficients $a_{i j}$ in (0.1.3).
0.3. There is by now an extensive literature on both the periodic and nonperiodic Toda lattice. There is an even larger literature on completely integrable Hamiltonian systems. As examples we cite references [1, 3, 8, 10, $19,20,24,25,27,28,29$. This, of course, is only a small part of a complete list. There are also connections with semi-simple Lie groups. See, e.g., [3, 20]. We would like to remark that there is a considerable distance between establishing the Poisson commutativity of certain functions-and actually integrating and finding the solution. In fact, there is a generalization of the periodic Toda lattice in that one replaces the Dynkin diagram by the extended Dynkin diagram. Independently this has been observed in [3]. We have been able to establish complete integrability in the sense above for this case as well, but, as yet, we have not been able to solve for $q_{i}(x(t))$, which presumably, as in [25], requires Abelian integrals.

We have been principally influenced by Refs. [19, 25]. In fact the starting point for us was our recognition that the symplectic form written by van Moerbeke [25, formula (42), p. 76] on the space of symmetric Jacobi matrices (and also cited by Moser in a lecture) was in fact the symplectic structure of a coadjoint orbit of the Borel subgroup $\bar{B}$.

Translation by the element $f$ (see (1.5.4)) plays an important role here for establishing complete integrability. On the other hand our proof of what is referred to as the Kostant-Symes splitting theorem in [27] (another complete integrability statement) involved other methods. We wish to thank S. Sternberg for pointing out that the proof of Theorem 1.4 here may be readily adapted so that the splitting theorem in [27] is a special case $(f=0)$.

The coadjoint orbit $Z$ of $\bar{B}$ is not only the setting for classical mechanics as considered in this paper but it also, recalling geometric quantization, as applied in [2], defines a unitary representation $\pi_{Z}$ of $\bar{B}$. Elsewhere it will be seen that the generalized Toda lattice is completely integrable and solvable not only in the classical sense, as established here, but also in the quantum sense. That is, one can write down the simultaneous spectral resolution of a commuting family of operators which includes the Schrödinger operator
associated to the Hamiltonian (0.1.1), where $U$ is given by (0.1.4) and the $\psi$ 's define a Dynkin diagram. This uses $\pi_{Z}$ and the Whittaker theory as developed in [17].

## 1. Poisson Commutativity of Translated Invariants

1.1. Assume $A$ is a commutative (associative) algebra over a field $F$, where $F$ has characteristic zero. In fact in our applications here $F$ will be either the field of real, $\mathbb{R}$, or complex, $\mathbb{C}$, numbers. We will say that $A$ has a Poisson structure if there is a bilinear map $A \times A \rightarrow A,(a, b) \rightarrow[a, b]$ with respect to which (1) $A$ is a Lie algebra and is such that (2)

$$
\begin{equation*}
[a, b c]=[a, b] c+b[a, c] \tag{1.1.1}
\end{equation*}
$$

for any $a, b, c \in A$.
If $V$ is a finite-dimensional vector space over $F$ then $S(V)=\oplus_{k=0}^{\infty} S_{k}(V)$ will denote the symmetric algebra over $V$ with its usual grading. In particular, $V=S_{1}(V)$ and $F=S_{0}(V)$. We may regard $S(V)$ as the algebra of polynomial functions on the dual space $V^{\prime}$ to $V$. If $v \in V$ and $v^{\prime} \in V^{\prime}$ the pairing of $v$ and $v^{\prime}$ is denoted by $\left\langle v^{\prime}, v\right\rangle$ so that one has $v\left(v^{\prime}\right)=\left\langle v^{\prime}, v\right\rangle$.

Now assume that $a$ is a finite-dimensional Lie algebra over $F$. For any $g \in a^{\prime}$ let $\partial(g)$ be the derivation of degree -1 of $S(a)$ such that $\partial(g) x=\langle g, x\rangle$ for any $x \in a$. Now if $u \in S(a)$ then as a polynomial function on $a^{\prime}$ the differential $d u$ defines a polynomial map $a^{\prime} \rightarrow a, g \mapsto(d u)(g)$. Explicitly one has

$$
\begin{equation*}
(d u)(g)=\sum_{i=1}^{n}\left(\left(\partial\left(g_{i}\right) u\right)(g)\right) x_{i} \tag{1.1.2}
\end{equation*}
$$

where $x_{i} \in a, i=1, \ldots, n$, is a basis of $a$ and $g_{i} \in a^{\prime}$ is the dual basis.
Proposition 1.1. If $u, v \in S(a)$ and $g \in a^{\prime}$ then

$$
\begin{equation*}
[u, v](g)=\langle g,[(d u)(g),(d v)(g)]\rangle \tag{1.1.3}
\end{equation*}
$$

defines a Poisson structure on $S(a)$ which extends the given Lie algebra structure on a.

Proof. If $u \in a$ then clearly $(d u)(g)=u$ so that (1.1.3) agrees with the bracket structure on $a$. From the differentiation properties of the exterior derivative it is clear that (1.1.1) is satisfied. One only has to see that the Jacobi identity is satisfied for any $u, v, w \in S(a)$. If $a d u(w)=[u, w]$ this amounts to showing that

$$
\begin{equation*}
a d[u, v]=[a d u, a d v] . \tag{1.1.4}
\end{equation*}
$$

However, both sides of (1.1.4) are derivations of the associative structure of $S(a)$ by (1.1.1) and hence one has equality (1.1.4) when applied to $w=w_{1} w_{2}$ if it already holds for $w_{1}$ and $w_{2}$. The proposition then follows easily since $a$ generates $S(a)$ and (1.1.4) holds for $u, v \in a$ when applied to $w \in a$. Q.E.D.

Henceforth $S(a)$, for any finite-dimensional Lie algebra $a$ will always have the Poisson structure given by Proposition 1.1.

Remark 1.2. Since a generates $S(a)$ note that the Poisson structure on $S(a)$ is the unique one which reduces on $a$ to the given Lie algebra structure on $a$.
1.2. Now assume that $g$ is a finite-dimensional Lie algebra over $F$ and $Q(x, y) \in F$ is a nonsingular invariant symmetric bilinear form on $g$. Invariance means that $Q([x, z], y)=Q(x,[z, y])$ for $x, y, z \in g$. If $x \in g$ and $g_{x} \in g^{\prime}$ is defined so that $\left\langle g_{x}, y\right\rangle=Q(x, y)$ then $x \rightarrow g_{x}$ defines an isomorphism $g \rightarrow g^{\prime}$. We regard $S(g)$ as the algebra of polynomial functions on $g$ itself where if $u \in S(g), x \in \mathscr{g}$ then $u(x)=u\left(g_{x}\right)$. In particular if $u \in g$ then

$$
\begin{equation*}
u(x)=Q(u, x) \tag{1.2.1}
\end{equation*}
$$

For any $x \in g$ let $i(x)$ be the derivation of degree -1 on $S(g)$ so that $i(x) y=$ $Q(x, y)$ when $y \in g$. Let $u \in S(g)$. To avoid confusion with $d u$ we will let $\delta u$ be the differential of $u$ when regarded as a function on $g$. Thus $\delta u$ is the polynomial map $g \rightarrow g$ given by

$$
\begin{equation*}
(\delta u)(x)=\sum_{j}\left(\left(i\left(x_{j}\right) u\right)(x)\right) y_{j} \tag{1.2.2}
\end{equation*}
$$

for $x \in g$, where $x_{i}, y_{j}$ are two bases of $g$ such that $Q\left(x_{i}, y_{j}\right)=\delta_{i j}$.
By Proposition 1.1 one immediately has

Lemma 1.2.1. The Poisson structure in $S(g)$ is given by

$$
\begin{equation*}
[u, v](x)=Q(x,[(\delta u)(x),(\delta v)(x)]) \tag{1.2.3}
\end{equation*}
$$

for any $u, v \in S(g), x \in g$.
Now assume that $F$ is either $\mathbb{R}$ or $\mathbb{C}$ and $g$ is a semi-simple Lie algebra over $F$. Let $Q$ be a fixed invariant symmetric bilinear form on $g$ which on each simple component is a positive multiple of the Killing form. Also let $g=k+p$ be a Cartan decomposition of $g$. Thus $k$ is the Lie algebra of a maximal compact subgroup of the adjoint group of $g$ and $\beta$ is the orthogonal complement to $k$ with respect to $\operatorname{Re} Q$. Now let $\theta$ be the corresponding Cartan involution. Thus $\theta=1$ on $k$ and $\theta=-1$ on $\not 2$. For any $x \in g$ let $x^{*}=-\theta x$. Since $Q$ is negative
definite on $k$ and positive definite on $\not p$ one defines an inner product $Q_{*}$ on $g$ (i.e., $g$ is a Hilbert space over $F$ with respect to $Q_{*}$ ) by putting

$$
\begin{equation*}
Q_{*}(x, y)=Q\left(x, y^{*}\right) . \tag{1.2.4}
\end{equation*}
$$

One notes that $(\lambda x)^{*}=\lambda x^{*}$ for $\lambda \in F$ and $x \in \mathscr{g}$, where the bar denotes conjugation. The relation $Q_{*}(x, y)=\overline{Q_{*}(y, x)}$ follows immediately from the readily verified relation $Q\left(x^{*}, y^{*}\right)=\overline{Q(x, y)}$. Also, the positive definiteness of $Q_{*}$ is immediate from the fact that $\operatorname{Re} Q(x, y)=0$ for $x \in k$ and $y \subset \not \subset$.

Now if $a \subseteq_{g}$ is any subspace let $a^{*}=\left\{x^{*} \mid x \in a\right\}$. Clearly $a^{*}$ is a subspace of the same dimension as $a$. There are two other subspaces, $a^{\perp}$, and $a^{o}$, which are also associated with $a$. They both have the same dimension as $g / a$. The subspace $a^{\perp}$ is defined as the orthogonal complement to $a$ with respect to $Q_{*}$ and $a^{0}$ is the orthogonal subspace to $a$ with respect to $Q$. Thus $a \rightarrow a^{*}$, $a^{1}$, and $a^{a}$ define three involutory operations on the set of all subspaces of $g$. The following lemma clarifies the relations between them. It asserts that together with the identity operation they define an action of the Klein 4-group on the set of all subspaces.

Lemma 1.2.2. The three operations $a \rightarrow a^{*}, a^{\perp}$, and $a^{o}$ commute with one another and the composite of any two distinct operations is the third.

The proof is straightforward and is left to the reader. A particular case of Lemma 1.2.2 is the relation

$$
\begin{equation*}
\left(a^{*}\right)^{\perp}=a^{0} . \tag{1.2.5}
\end{equation*}
$$

Now if $a \subseteq g$ is a subspace we regard $S(a)$ as a subalgebra of $S(g)$. The fact that $Q_{*}$ is an inner product implies that $a$ and $a^{*}$ are non-singularly paired by $Q$.

Lemma 1.2.3. Let $u \in S(a)$. Then for any $x \in g, y \in a^{o}$ one has

$$
\begin{equation*}
u(x+y)=u(x) \tag{1.2.6}
\end{equation*}
$$

Further, the restriction map $u \rightarrow u \mid a^{*}$ defines an isomorphism of $S(a)$ onto the algebra of all polynomial functions on $a^{*}$.

Proof. Clearly $v(y)=0$ for all $v \in a$. Since $u$ is generated by elements in $a$ one has (1.2.6). On the other hand since $a$ and $a^{*}$ are non-singularly paired by $Q$ the map $v \rightarrow v \mid a^{*}$ defines an isomorphism of $S(a)$ onto the algebra of all polynomial functions on $a^{*}$.
Q.E.D.

Now let $P_{a}: g \rightarrow a$ be the orthogonal projection of $g$ on $a$ with respect
to the inner product $Q_{*}$. If $u \in S(g)$ we recall (see (1.2.2)) that $\delta u$ is a polynomial map $g \rightarrow g$. We now define $\delta_{a} u$ to be the map $g \rightarrow a$ given by

$$
\begin{equation*}
\left(\delta_{a} u\right)(x)=P_{a}((\delta u)(x)) . \tag{1.2.7}
\end{equation*}
$$

Obviously if $a=a_{1}+a_{2}$ is an orthogonal direct sum with respect to $Q_{*}$ then

$$
\begin{equation*}
\delta_{a} u=\delta_{a_{1}} u+\delta_{a_{2}} u \tag{1.2.8}
\end{equation*}
$$

In particular one has

$$
\begin{equation*}
\delta u=\delta_{a^{\prime}} u+\delta_{a^{-}} u \tag{1.2.9}
\end{equation*}
$$

Explicitly one notes
Lemma 1.2.4. Let $x_{i}, i=1, \ldots, d$, be any basis of a and let $y_{j} \in a^{*}$ be the basis of $a^{*}$ such that $Q\left(x_{i}, y_{j}\right)=\delta_{i j}$. Then $\delta_{a} u$ is the polynomial map $g \rightarrow a$ given by

$$
\begin{equation*}
\left(\delta_{a} u\right)(x)=\sum_{j=1}^{d}\left(i\left(y_{j}\right) u\right)(x) x_{j} \tag{1.2.10}
\end{equation*}
$$

for any $x \in g$.
Proof. Let $x_{i}, i=d+1, \ldots, n=\operatorname{dim} g$, be any basis of $a^{\perp}$ and let $y_{i}$, $i=d+1, \ldots, n$, be the basis of $\left(a^{\perp}\right)^{*}$ such that $Q\left(x_{i}, y_{j}\right)=\delta_{i j}$ for $i, j \geqslant d+1$. Put $w_{2}=\sum_{j=d+1}^{n}\left(i\left(y_{j}\right) u\right)(x) x_{j}$ and let $w_{1}$ be given by the right side of (1.2.10). Now since $a^{*}=\left(a^{\perp}\right)^{0}$ and $\left(a^{\perp}\right)^{*}=a^{0}$ by Lemma 1.2.2 it follows that $Q\left(x_{i}, y_{j}\right)=\delta_{i j}$ for $i, j=1, \ldots, n$. Thus $w=w_{1}+w_{2}=(\delta u)(x)$ by (1.2.9). But then $w_{1}=\left(\delta_{a} u\right)(x)$. This proves the lemma.
Q.E.D.

Lemma 1.2.5. Let $a \subseteq g$ be any subspace and let $u \in S(a)$. Then

$$
\begin{equation*}
\delta u=\delta_{a} u \tag{1.2.11}
\end{equation*}
$$

so that $\delta u$ is a polynomial map $g \rightarrow a$.
Proof. By (1.2.9) we have only to show that $\delta_{\alpha} \perp u=0$. But now applying (1.2.10) where $a^{\perp}$ replaces $a$ one notes that $i\left(y_{j}\right) u=0$ since $y_{j} \in\left(a^{\perp}\right)^{*}=a^{o}$, using Lemma 1.2.2. Thus $\delta_{a} \perp u=0$.
Q.E.D.

Now if $a$ is a Lie subalgebra of $g$ then $S(a)$ inherits a Poisson structure from the Lie algebra structure of $a$. This has been defined using the dual space $a^{\prime}$ to $a$. See Proposition 1.1. We now note that $a^{\prime}$ may be replaced by $a^{*} \subseteq g$.

Proposition 1.2. Let a be a Lie subalgebra of $\mathfrak{y}$. Then $S(a)$ is a Lie subalgebra of $S(g)$ with respect to the Poisson structure on $S(g)$. Furthermore the Lie algebra structure thus induced on $S(a)$ is the same as the Poisson structure on $S(a)$ which $S(a)$ would normally inherit from the Lie algebra structure in $a$. Thus if $u, v \in S(a)$ then $[u, v]$ is unambiguous and is in $S(a)$. Moreover it is given by

$$
\begin{equation*}
[u, v](x)=Q(x,[(\delta u)(x),(\delta v)(x)]) \tag{1.2.12}
\end{equation*}
$$

for $x \in g$ but is in fact determined by the restriction $[u, v] \mid a^{*}$. In particular, $[u, v]=0$ if and only if the right side of (1.2.12) vanishes for all $x \in a^{*}$.

Proof. The first statement follows from (1.1.1). The second and hence the third follows from Remark 1.2. Relation (1.2.12) is just (1.2.3). Since $[u, v] \in S(a)$ it is determined using Lemma 1.2.3, by the restriction $[u, v] \mid a^{*}$.
Q.E.D.
1.3. Now let $G=$ Aut $g$ be the adjoint group of $g$. The action of $a \in G$ on $x \in_{\mathcal{g}}$ is denoted simply by $a x \in \mathscr{g}$. The algebra $S(g)$ then becomes a $G$-module where if $a \in G, u \in S(g)$, and $x \in g$ one has $a u(a x)=u(x)$. From the invariance of $Q$ it follows that $a(i(x) u)=i(a x) a u$. It then follows easily from (1.2.2) that

$$
\begin{equation*}
a((\delta u)(x))=(\delta(a u))(a x) \tag{1.3.1}
\end{equation*}
$$

for any $u \in S(g), a \in G$, and $x \in g$.
Now let $l=\operatorname{rank} g$ and let $S(g)^{G}$ be the algebra of $G$-invariants in $S(g)$. By Chevalley's theorem one knows there exist homogeneous elements $I_{j} \in S(\mathfrak{g})^{G}$, $j=1,2, \ldots, l$, referred to as the fundamental invariants, which are algebraically independent and which generate $S(g)^{G}$. That is, by abuse of notation,

$$
\begin{equation*}
S(g)^{G}=F\left[I_{1}, \ldots, I_{l}\right] \tag{1.3.2}
\end{equation*}
$$

Now for any $x \in \mathscr{g}$ let $g^{x}$ denote the centralizer of $x$ in $g$ and let cent $\mathscr{g}^{x}$ be the center of $\mathfrak{g}^{x}$. Also let $G^{x}$ be the centralizer (or rather stabilizer) of $x$ in $G$.

Proposition 1.3. For any $x \in g$ and $I \in S(g)^{G}$ one has $(\delta I)(x) \in$ cent $g^{x}$. In particular one has

$$
\begin{equation*}
[x,(\delta I)(x)]=0 \tag{1.3.3}
\end{equation*}
$$

Proof. For any $a \in G$ one has $a I=I$. Thus $a((\delta I)(x))=(\delta I)(a x)$ by (1.3.1). But then if $a \in G^{x}$ one has $a((\delta I)(x))=(\delta I)(x)$. That is, if $g^{\left(G^{x}\right)}$ is the set of

where $g^{\left(\mathscr{g}^{x}\right)}$ is the centralizer of $\mathscr{g}^{x}$ in $\mathscr{g}$. However, since $x \in \mathscr{g}^{x}$ one has $\mathscr{g}^{\left(\mathfrak{g}^{x}\right)}=$ cent $\dot{g}^{x}$. This proves the proposition.
Q.E.D.
1.4. Now let $a$ be a Lie subalgebra of $g$. We shall regard $a$ as fixed. This will enable us to suppress $a$ in notation which in fact depends upon $a$.

Now let $f \in g$ and $u \in S(g)$. The map $x \mapsto u(f+x)$ is clearly a polynomial function on $g$. In particular its restriction to $a^{*}$ is a polynomial function. By Lemma 1.2.3 therefore there exists a unique element $u^{f} \in S(a)$ such that

$$
\begin{equation*}
u^{f}(x)=u(f+x) \tag{1.4.1}
\end{equation*}
$$

for all $x \in a^{*}$. The correspondence $u \rightarrow u^{f}$ therefore defines a homomorphism

$$
\begin{equation*}
S(g) \rightarrow S(a) \tag{1.4.2}
\end{equation*}
$$

We are particularly interested in the case where $u \in S(g)^{G}$ and especially when $a$ satisfies the following condition.

We say that $a$ is a Lie summand in case $a^{0}$ is also a Lie subalgebra of $g$. The terminology stems from the fact that $a$ is a Lie summand if and only if $a^{\perp}$ is a Lie subalgebra of $g$. This is clear since, by Lemma 1.2.2, $a^{\perp}=\left(a^{0}\right)^{*}$ and one has $[x, y]^{*}=\left[y^{*}, x^{*}\right]$.

Example. Any parabolic subalgebra of $\mathfrak{g}$ is a Lie summand. This is clear since if $a$ is parabolic then $a^{o}$ is just the nilradical of $a$.

It will be seen later that the following theorem leads to the complete integrability of the Toda lattice and to extensive generalizations of it.

Theorem 1.4. Let $g$ be a semi-simple Lie algebra over $\mathbb{R}$ or $\mathbb{C}$. Let $Q$ be a bilinear form on $g$ which on each simple component is a fixed positive multiple of the Killing form and let $Q_{*}$ be the inner product in $g$ defined by $Q$ and a Cartan decomposition of $g$. See (1.2.4).

Let $a \subseteq \mathscr{g}$ be a Lie summand. That is, a and its orthocomplement $a^{\perp}$ with respect to $Q_{*}$ are Lie subalgebras of $g$. Let $f \in g$ be any element such that

$$
\begin{equation*}
Q\left(f,[a, a]+\left[a^{\perp}, a^{\perp}\right]\right)=0 \tag{1.4.3}
\end{equation*}
$$

Then, using the notation of (1.4.1) the elements $I^{f}, J^{f}$ in $S(a)$ Poisson commute for any pair of invariants $I, J \in S(g)^{G}$.

Proof. Using Proposition 1.2 it is enough to show that

$$
\begin{equation*}
Q\left(x,\left[\left(\delta I^{f}\right)(x),\left(\delta J^{f}\right)(x)\right]\right)=0 \tag{1.4.4}
\end{equation*}
$$

for any $x \in a^{*}$. But now we assert that

$$
\begin{equation*}
(\delta I f)(x)=\left(\delta_{a} I\right)(f+x) \tag{1.4.5}
\end{equation*}
$$

for $x \in a^{*}$. Indeed for any $y, z \in g$ and $u \in S(g)$ one has $(i(y) u)(z)=(d / d t)$ $\left.u(z+t y)\right|_{t=0}$. It follows therefore that if $x, y \in a^{*}$ one has $\left(i(y) I^{f}\right)(x)=$ $(i(y) I)(f+x)$. Thus recalling (1.2.10) one has $\left(\delta_{a} I^{f}\right)(x)=\left(\delta_{a} I\right)(f+x)$ for $x \in a^{*}$. But $\delta_{a} I^{f}=\delta I^{f}$ by (1.2.11). This proves (1.4.5).

But now $\left[\left(\delta_{a} I\right)(f+x),\left(\delta_{a} J\right)(f+x)\right] \in[a, a]$. However, $f$ is $Q$-orthogonal to $[a, a]$. That is,

$$
\begin{equation*}
Q\left(f,\left[\left(\delta_{a} I\right)(f+x),\left(\delta_{a} J\right)(f+x)\right]\right)=0 \tag{1.4.6}
\end{equation*}
$$

But then adding this relation to the left side of (1.4.4) and recalling (1.4.5) it is enough to show that

$$
\begin{equation*}
Q\left(f+x,\left[\left(\delta_{a} I\right)(f+x),\left(\delta_{a} J\right)(f+x)\right]\right)=0 \tag{1.4.7}
\end{equation*}
$$

But now by (1.3.3) when $f+x$ replaces $x$ one has

$$
\begin{equation*}
[f+x,(\delta I)(f+x)]=0 \tag{1.4.8}
\end{equation*}
$$

Thus by the invariance of $Q$ one does in fact have

$$
\begin{equation*}
Q\left(f+x,\left[(\delta I)(f+x),\left(\delta_{a} J\right)(f+x)\right]\right)=0 \tag{1.4.9}
\end{equation*}
$$

But $(\delta I)(f+x)=\left(\delta_{a} I\right)(f+x)+\left(\delta_{a \perp} I\right)(f+x)$ by (1.2.9). Substituting in (1.4.9) it then suffices to prove

$$
\begin{equation*}
Q\left(f+x,\left[\left(\delta_{a^{+}} I\right)(f+x),\left(\delta_{a} J\right)(f+x)\right]\right)=0 \tag{1.4.10}
\end{equation*}
$$

But (1.4.8) is valid if $J$ replaces $I$. Thus by the invariance of $Q$ one does have (1.4.10) if $\delta J$ replaces $\delta_{a} J$. Thus it suffices to show

$$
Q\left(f+x,\left[\left(\delta_{a^{+}} I\right)(f+x),\left(\delta_{a^{+}} J\right)(f+x)\right]\right)=0
$$

But by (1.4.3) one also has (1.4.6) when $a^{\perp}$ replaces $a$. Thus it suffices only to show that

$$
\begin{equation*}
Q\left(x,\left[\left(\delta_{a^{\perp}} I\right)(f+x),\left(\delta_{a^{\perp}} J\right)(f+x)\right]\right)=0 \tag{1.4.11}
\end{equation*}
$$

But since $a^{\perp}$ is a Lie algebra the commutator in (1.4.11) is an element in $a^{\perp}$. However, $x \in a^{*}$ and $a^{*}=\left(a^{\perp}\right)^{o}$ by Lemma 1.2.2. This proves (1.4.11).
Q.E.D.
15. We consider an example now which illustrates Theorem 1.4 and which plays a major role in the remainder of the paper. No special assumptions need be made if $F=\mathbb{C}$. However, if $F=\mathbb{R}$ then for this example we assume that $g$ is split. We may then find, as one knows, using, say, Weyl's normal form (see, e.g., Section 5 in [13]), a Cartan subalgebra $h \subseteq g$ (which is split if $F=\mathbb{R}$ and root vectors $e_{\phi} \in g, \varphi \in \Delta$, where $\Delta=\Delta(g, h)$ is the corresponding set of roots such that

$$
\begin{equation*}
e_{\varphi}^{*}=e_{-\Phi} \tag{1.5.1}
\end{equation*}
$$

and

$$
\begin{equation*}
Q\left(e_{\Phi}, e_{-\varphi}\right)=1 \tag{1.5.2}
\end{equation*}
$$

Let $\Delta_{+} \subseteq \Delta$ be a fixed chosen system of positive roots and let $\left\{\alpha_{1}, \ldots, \alpha_{l}\right\}=$ $\Pi \subseteq \Delta_{+}$be the set of simple positive roots.

Remark 1.5.1. One notes that the root vectors $e_{-\alpha_{i}}, i=1, \ldots, l$, may initially be chosen arbitrarily. Then $e_{\alpha_{i}}$ are determined by (1.5.2). This then, using (1.5.1), determines the Cartan decomposition $g=k+\not p$.

Now

$$
\begin{equation*}
\bar{b}=\hbar+\sum_{\varphi=\Delta_{+}} F e_{-\varphi} \quad \text { and } \quad \ell=h+\sum_{\varphi \in \Delta_{+}} F e_{\varphi} \tag{1.5.3}
\end{equation*}
$$

are "opposing" Borel subalgebras and $\bar{n}=[\bar{b}, \bar{b}]$ and $n=[b, b]$ are the corresponding nilradicals.

Theorem 1.5. Let $a=t$ and let

$$
\begin{equation*}
f=\sum_{i=1}^{l} e_{-\alpha_{i}} \tag{1.5.4}
\end{equation*}
$$

Then a and $f$ satisfy the conditions of Theorem 1.4. That is, $\bar{t}$ is a Lie summand and $f$ satisfies (1.4.3). In particular the elements $I^{f}$ and $J^{f}$ in $S(\bar{b})$ Poisson commute for any $I, J \in S(g)^{G}$.

Proof. If $\varphi \in \Delta$ then as one knows $Q\left(x, e_{\varphi}\right)=Q\left(e_{\psi}, e_{\varphi}\right)=0$ for any $x \in \mathscr{h}$, and $\psi \in \Delta$, where $\psi \neq-\varphi$. Since

$$
\begin{equation*}
n=\sum_{m \in \Delta_{+}} F e_{\varphi} \quad \text { and } \quad \bar{n}=\sum_{\varphi \in \Delta_{+}} F e_{-\infty} \tag{1.5.5}
\end{equation*}
$$

one then has $(\bar{b})^{o}=\bar{n}$. Hence $\bar{b}$ is a Lie summand. But $[a, a]=[\bar{b}, \vec{b}]=$ $\bar{n}=(\bar{b})^{o}$. Since $f \in \bar{b}$ this implies $Q(f,[a, a])=0$. But by Lemma 1.2.2 one also has $a^{\perp}=\left(a^{0}\right)^{*}=(\bar{n})^{*}=n$. Thus [ $a^{\perp}, a^{\perp}$ ] is spanned by all $e_{\varphi}$, where $\varphi \in A_{+}$is not simple. Thus $Q\left(f,\left[a^{\perp}, a^{\perp}\right]\right)=0$. Hence $f$ satisfies (1.4.3). The remaining statement follows from Theorem 1.4.
Q.E.D.

Remark 1.5.2. If $N \subseteq G$ is the subgroup corresponding to $n$ then in Section 1.3 of [17] we defined an action of $N$ on $S(\bar{b})$ so that if $S(\bar{b})^{N}$ is the algebra of invariants then when $a=\bar{b}$

$$
\begin{equation*}
S(\bar{b})^{N}=\left\{I^{f} \mid I \in S(g)^{G}\right\} \tag{1.5.6}
\end{equation*}
$$

Theorem 1.5 then says that the elements of the algebra $S(k)^{N}$ Poisson commute with one another. This algebra is very explicit in the case where $g$ is the Lie algebra of $S l(n, F)$. In that case choices can be made so that $f+b=f+a^{*}$ is the set of all traceless matrices $y$ of the form

$$
y=\left(\begin{array}{cccc}
a_{11} & a_{12} & \cdots & a_{1 n}  \tag{1.5.7}\\
1 & & & \vdots \\
\vdots & & & a_{n-1 n} \\
0 & \cdots & 1 & a_{n n}
\end{array}\right)
$$

By abuse of notation regard $a_{i j}$ as a linear functional on $\ell$ and thereby identify it with an element in $\mathscr{Z}$. The nontrivial coefficients $u_{k} \in S(\mathscr{C}), k=1, \ldots, l=$ $n-1$, of the characteristic polynomial of $y$, as polynomials in the $a_{i j}$, are of the form $u_{k}=I_{k}{ }^{f}$, where the $I_{k}$ are the fundamental invariants. The algebra $S(\bar{b})^{N}$ is just the algebra generated by the $u_{k}$. Theorem 1.5 asserts that they Poisson commute with one another. It is this example which applies directly to the Toda lattice.
1.6. Let $a$ and $f$ be as in Theorem 1.5 so that $a$ is a Lie summand in $g$ and $f \in g$ satisfies (1.4.3). Now consider the translation $f+a^{*}$ of $a^{*}$ by $f$. For convenience put $f+a^{*}=\left(a^{*}\right)_{f}$

$$
\begin{equation*}
\tau_{f}: a^{*} \rightarrow\left(a^{*}\right)_{f}, \quad x \rightarrow f+x . \tag{1.6.1}
\end{equation*}
$$

We will understand that $\left(a^{*}\right)_{f}$ has the structure of an affine variety and that $\tau_{\boldsymbol{f}}$ is an isomorphism of affine varieties.

Lemma 1.6.1. For any $I \in S(g)^{G}$ and $y \in\left(a^{*}\right)_{f}$ one has $\left[y,\left(\delta_{a} I\right)(y)\right] \in a^{*}$.
Proof. Put $w=\left(\delta_{a} I\right)(y)$ so that $w+\left(\delta_{a} I\right)(y)=(\delta I)(y)$ by (1.2.9). But $(\delta I)(y)$ commutes with $y$ by (1.3.3). Thus

$$
\begin{equation*}
\left[y,\left(\delta_{\alpha} I\right)(y)\right]=[w, y] . \tag{1.6.2}
\end{equation*}
$$

But now write $y=f+x$, where $x \in a^{*}$. We assert that

$$
\begin{equation*}
[w, f] \in a^{*} \tag{1.6.3}
\end{equation*}
$$

Indeed since $a^{*}=\left(a^{\perp}\right)^{n}$ it suffices to show that $[w, f]$ is $Q$-orthogonal to $a^{\perp}$. But $Q\left(a^{\perp},[w, f]\right)=Q\left(\left[a^{\perp}, w\right], f\right)=0$ by (1.4.3) since $w \in a^{\perp}$. This proves (1.6.3). Thus to prove the lemma it suffices by (1.6.2) to show that $[w, x] \in a^{*}$. But to prove this it certainly suffices to show that

$$
\begin{equation*}
\left[a^{\perp}, a^{*}\right] \subseteq a^{*} \tag{1.6.4}
\end{equation*}
$$

But since $a^{*}=\left(a^{\perp}\right)^{0}$ it suffices to show that $\left[a^{\perp},\left(a^{\perp}\right)^{0}\right]$ is $Q$-orthogonal to $a^{\perp}$. However, $Q\left(a^{\perp},\left[a^{\perp},\left(a^{\perp}\right)^{\circ}\right]\right) \subseteq Q\left(\left[a^{\perp}, a^{\perp}\right],\left(a^{\perp}\right)^{\circ}\right)=0$ since $a^{\perp}$ is a Lie algebra. This proves (1.6.4).
Q.E.D.

Now $g$, whether or not $F=\mathbb{R}$ or $\mathbb{C}$, has the structure of a vector space over $\mathbb{R}$. Thus if $y \in g$ the tangent space to $g$ at $y$ is naturally identified with $g$. In this identification if $y \in\left(a^{*}\right)_{f}$ then the tangent space to $\left(a^{*}\right)_{f}$ at $y$ is clearly $a^{*}$. Let $I \in S(g)^{G}$. The map given by $y \rightarrow\left[y,\left(\delta_{a} I\right)(y)\right]$ is clearly a polynomial map on $g$ and hence by Lemma 1.6 .1 there exists a smooth vector field $\eta_{I}$ on $\left(a^{*}\right)_{f}$ such that

$$
\begin{equation*}
\left(\eta_{I}\right)_{y}=\left[y,\left(\delta_{z} I\right)(y)\right] \tag{1.6.5}
\end{equation*}
$$

for any $y \in\left(a^{*}\right)_{f}$.
Now let $\tau_{f}^{-1}:\left(a^{*}\right)_{f} \rightarrow a^{*}$ be the isomorphism which is inverse to $\tau_{f}$. See (1.6.1). Thus $\tau_{f}^{-1}$ "carries" vector fields on $\left(a^{*}\right)_{f}$ to vector fields on $a^{*}$. In particular, using the same notation for this map of vector fields, $\eta_{I} \rightarrow \tau_{f}^{-1} \eta_{I}$ the vector field $\tau_{f}^{-1} \eta_{I}$ can be applied to an element $u \in S(g)$ since by restriction $u$ defines a function on $a^{*}$. We express this action in terms of Poisson brackets.

Lemma 1.6.2. For any $I \in S(g)^{G}, u \in S(a), x \in a^{*}$ one has

$$
\begin{equation*}
\left(\left(\tau_{f}^{-1} \eta_{I}\right) u\right)(x)=\left[I^{f}, u\right](x) \tag{1.6.6}
\end{equation*}
$$

Proof. Let $y=f+x=\tau_{f} x$. We express the left side of (1.6.6) in terms of the pairing of vectors and covectors. Let $c$ be the left side of (1.6.6). Recalling the definition (see (1.2.2)) of $\delta u$ one has

$$
\begin{aligned}
c & =Q\left(\left[y,\left(\delta_{a} I\right)(y)\right],(\delta u)(x)\right) \\
& =Q\left(y,\left[\left(\delta_{a} I\right)(y),(\delta u)(x)\right]\right)
\end{aligned}
$$

But $\left(\delta_{a} I\right)(f+x)=\left(\delta I^{f}\right)(x)$ by (1.4.5). Thus

$$
c=Q\left(y,\left[\left(\delta I^{f}\right)(x),(\delta u)(x)\right]\right)
$$

But $y=f+x$. Hence $c-\left[I^{f}, u\right](x)+Q\left(f,\left[\left(\delta I^{f}\right)(x),(\delta u)(x)\right]\right)$, recalling the formula (1.2.3). It suffices therefore to show that

$$
\begin{equation*}
Q\left(f,\left[\left(\delta I^{\prime}\right)(x),(\delta u)(x)\right]=0\right. \tag{1.6.7}
\end{equation*}
$$

But since $I^{f}, u \in S(a)$ the commutator in (1.6.7) is in [a,a] by Lemma 1.2.5. Thus (1.6.7) follows from (1.4.3).
Q.E.D.

Theorems 1.6.1 and 1.6.2 below are corollaries of Theorem 1.5. They are familiar consequences in a symplectic context. However, the latter does not quite apply here.

Theorem 1.6.1. Let $g$ be a semi-simple Lie algebra over $\mathbb{R}$ or $\mathbb{C}$. Let $a \subseteq g$ and $f \in g$ satisfy the conditions of Theorem 1.5 and put $\left(a^{*}\right)_{f}=f+a^{*}$. For any invariant $I \in S(g)^{G}$ let $\eta_{I}$ be the vector field on $\left(a^{*}\right)_{f}$ defined by (1.6.5). Then

$$
\begin{equation*}
\eta_{I} J=0 \tag{1.6.8}
\end{equation*}
$$

on $\left(a^{*}\right)_{f}$ for any $I, J \in S(g)^{G}$.
Proof. Apply Lemma 1.6.2, where $u=J^{f}$. But then by Theorem 1.5 one has $\left(\left(\tau_{f}^{-1} \eta_{I}\right) u\right)(x)=0$ for any $x \in a^{*}$. However, if $y=f+x$ then $\left(\left(\tau_{f}^{-1} \eta_{I}\right) u\right)(x)=\left(\eta_{I}\left(u \circ \tau_{f}^{-1}\right)\right)(y)$. But clearly $u \circ \tau_{f}^{-1}=J \mid\left(a^{*}\right)_{f}$. This proves the theorem.
Q.E.D.

Now let $P_{a^{*}}: g \rightarrow a^{*}$ be the $Q_{*}$-orthogonal projection of $g$ onto $a^{*}$. For any $y \in\left(a^{*}\right)_{f}$ let

$$
\begin{equation*}
a^{*}(y)=\left\{P_{a *}[x, y] \mid x \in a\right\} \tag{1.6.9}
\end{equation*}
$$

Remark 1.6.1. If $A \subseteq G$ is the subgroup corresponding to $a$ then one may define an action of $A$ on $\left(a^{*}\right)_{f}$ by translating to $\left(a^{*}\right)_{f}$, using $\tau_{f}$, the coadjoint action of $A$ on $a^{*}$. The subspace $a^{*}(y)$ of $a^{*}$ then turns out to be nothing more than the tangent space at $y$ to the $A$ orbit containing $y$. By Lemma 1.6.1 one has $P_{a *}\left[y,\left(\delta_{a} I\right)(y)\right]=\left[y,\left(\delta_{a} I\right)(y)\right]$. Since $\left(\delta_{a} I\right)(y) \in a$ it follows that

$$
\begin{equation*}
\left(\eta_{t}\right)_{y} \in a^{*}(y) \quad \text { for any } \quad I \in S(g)^{G} \tag{1.6.10}
\end{equation*}
$$

Theorem 1.6.2. Let $g$ be a semi-simple Lie algebra over $\mathbb{C}$ or $\mathbb{R}$. Let $a \subseteq g$ and $f \in_{g}$ satisfy the conditions of Theorem 1.5 and put $\left(a^{*}\right)_{f}=f+a^{*}$. For any invariant $I \in S(g)^{G}$ let $\eta_{I}$ be the vector field on $\left(a^{*}\right)_{f}$ defined by (1.6.5). Then

$$
\begin{equation*}
\left[\eta_{I}, \eta_{J}\right]=0 \tag{1.6.11}
\end{equation*}
$$

for any $I, J \in S(g)^{G}$. Furthermore if $y \in\left(a^{*}\right)_{f}$ then

$$
\begin{equation*}
\left(\eta_{I}\right)_{y}=0 \Leftrightarrow Q\left(a^{*}(y),(\delta I)(y)\right)=0 \tag{1.6.12}
\end{equation*}
$$

where $a^{*}(y)$ is defined by (1.6.9) and $\delta I$ is defined by (1.2.2).
Proof. We may iterate the formula (1.6.6). Let $I, J \in S(g)^{G}$ and $u \in S(a)$. That is, if $x \in a^{*}$, then

$$
\begin{equation*}
\left(\left(\tau_{f}^{-1} \eta_{I}\right)\left(\tau_{f}^{-1} \eta_{J}\right) u\right)(x)=\left[I^{f},\left[J^{f}, u\right]\right](x) . \tag{1.6.13}
\end{equation*}
$$

But then $\left(\left(\tau_{f}^{-1}\left[\eta_{I}, \eta_{J}\right]\right) u\right)(x)=\left[\left[I^{f}, J^{f}\right] u\right](x)$. However, $\left[I^{f}, J^{\prime}\right]=0$ by Theorem 1.5. But since the set of functions $u \mid a^{*}, u \in S(a)$ contains a coordinate system (e.g., let $u \in a$ ) it follows that $\tau_{f}^{-1}\left[\eta_{I}, \eta_{J}\right]=0$ and hence one has (1.6.11).

Now if $x \in a$ and $y \in\left(a^{*}\right)_{f}$ then

$$
\begin{align*}
Q\left(x,\left(\eta_{I}\right)_{y}\right) & =Q\left(x,\left[y,\left(\delta_{a} I\right)(y)\right]\right) \\
& =Q\left([x, y],\left(\delta_{a} I\right)(y)\right) \tag{1.6.14}
\end{align*}
$$

But now since $\left(a^{*}\right)^{\perp}=a^{0}$ we may replace $[x, y]$ in (1.6.14) by $P_{a^{*}}[x, y]$. But then having done that we may replace $\left(\delta_{a} I\right)(y)$ by $(\delta I)(y)$ since $\left(a^{\perp}\right)^{o}=a^{*}$. This uses (1.2.9). Thus

$$
\begin{equation*}
Q\left(x,\left(\eta_{I}\right)_{y}\right)=Q\left(P_{a *}[x, y],(\delta I)(y)\right) \tag{1.6.15}
\end{equation*}
$$

But since $\left(\eta_{I}\right)_{y} \in a^{*}$ by Lemma 1.6 .1 and $a$ and $a^{*}$ are non-singularly paired by $Q$ it follows that $\left(\eta_{I}\right)_{y} \neq 0$ if and only if there exists $x \in a$ such that (1.6.15) is not zero. This implies (1.6.12).
Q.E.D.

Remark 1.6.2. Recalling Remark 1.6 .1 one may interpret (1.6.12) as follows: Let $O$ denote the $A$-orbit of $y$ with respect to the $f$-translated coadjoint action of $A$ on $\left(a^{*}\right)_{f}$ described in Remark 1.6.1. Then $\left(\eta_{I}\right)_{y} \neq 0$ if and only if the differential $d(I \mid O)_{y} \neq 0$, where of course $I \mid O$ is the restriction of $I$ to $O$. In fact if $F=\mathbb{R}$ then $O$ has the structure of a symplectic manifold and $\eta_{I} \mid O$ (see Section 6.4 for the case where $O=Z$ ) is just the Hamiltonian vector field corresponding to the function $I \mid O$.

## 2. The Variety $Z$ of Normalized Jacobi Elements

2.1. Henceforth we restrict our attention to the case of the example in Section 1.5. We recall that there is no restriction on $g$ if $F=\mathbb{C}$ but $g$ is split if $F=\mathbb{R}$. Furthermore $\ell$ (which is split if $F=\mathbb{R}$ ), $\Delta, \Delta_{+}, \Pi=\left\{\alpha_{1}, \ldots, \alpha_{l}\right\}$, $e_{\varphi}$, for $\varphi \in \Delta$, satisfying (1.5.1) and (1.5.2) are as in Section 1.5.

Now let $H \subseteq G$ be the subgroup corresponding to $h$. Let $\hbar^{\prime}=\operatorname{Hom}_{F}(h, F)$ and where $\mathbb{Z}$ is the set of integers let

$$
\begin{equation*}
h^{\prime}(H)=\left\{\nu \in h^{\prime} \mid\langle\nu, x\rangle \in 2 \pi i \mathbb{Z} \text { for } x \in h, \text { where } \exp x=1\right\} . \tag{2.1.1}
\end{equation*}
$$

Remark 2.1. One knows that $\ell^{\prime}(H)$ is the lattice $\sum_{i=1}^{l} \mathbb{Z} \alpha_{i}$ if $F=\mathbb{C}$ and $\ell^{\prime}(H)=\ell^{\prime}$ if $F=\mathbb{R}$.

Let $F^{*}=\exp F$ so that $F^{*}$ is the multiplicative group of nonzero complex numbers if $F=\mathbb{C}$ and $F^{*}$ is the multiplicative group of positive numbers if $F=\mathbb{R}$. Then any $\nu \in h^{\prime}(H)$ defines a homomorphism $H \rightarrow F^{*}, h \rightarrow h^{\nu}$, where $h^{v}=\exp \langle\nu, x\rangle$ for $x \in h$ such that $h=\exp x$. One of course has $\Delta \subseteq h^{\prime}(H)$ and if $h \in H, \varphi \in \Delta$ then

$$
\begin{equation*}
h e_{\varphi}=h^{\Phi} e_{\varphi} \tag{2.1.2}
\end{equation*}
$$

Now let $x_{o} \in h$ be that unique element such that $\left\langle\alpha_{i}, x_{o}\right\rangle=1$ for $i=1, \ldots, l$. The eigenvalues of ad $x_{o}$ are in $\mathbb{Z}$. One puts $d_{j}=\left\{x \in g \mid\left[x_{o}, x\right]=j x\right\}$ for $j \in \mathbb{Z}$. We shall refer to the $d_{j}$ as the diagonals of $g$. Of course $d_{0}=h$ and one has the direct sum decomposition

$$
\begin{equation*}
\mathscr{y}=\oplus d_{j} \tag{2.1.3}
\end{equation*}
$$

where the sum is over $\mathbb{Z}$.
Definition 2.1. If $O \neq x \in g$ and $x_{j} \in d_{j}$ denotes the component of $x$ in $d_{j}$ relative to (2.1.3) then the minimal $j$ such that $x_{j} \neq 0$ will be called the minimal diagonal degree of $x$. The element $x_{j}$ is called the minimal diagonal component of $x$. Similarly we will speak of the maximal diagonal degree and component of $x$.

One of course has, for $i, j \in \mathbb{Z}$,

$$
\begin{equation*}
\left[d_{i}, d_{j}\right] \subseteq d_{i+j} \tag{2.1.4}
\end{equation*}
$$

Now let $\ell, \vec{b}, n$, and $\bar{n}$ be the Borel subalgebras and their commutators defined as in Section 1.5 so that

$$
\begin{equation*}
\ell=\sum_{j \geqslant 0} d_{j}, \quad \bar{\ell}=\sum_{j \leqslant 0} d_{j} \tag{2.1.5}
\end{equation*}
$$

and

$$
\begin{equation*}
n=\sum_{j \geqslant 1} d_{j}, \quad \bar{n}=\sum_{j \leqslant 1} d_{j} . \tag{2.1.6}
\end{equation*}
$$

Henceforth unless otherwise stated as in the example of Section 1.5 the
subalgebra $a$ of Sections $1.2-1.5$ will be just $Z$. Thus all notation and results referring to $a$ now apply solely to $b$. Furthermore the element $f$ is fixed and is given by (1.5.4). Thus in particular, $u^{f} \in S(\vec{b})$ for any $u \in S(g)$. Furthermore $a^{*}=\ell$ by (1.5.1) and $\left(a^{*}\right)_{f}=f+b$ which of course we now write as $\ell_{f}$. For the case of $S l(n, F)$ we note that choices can be made so that the typical element $y \in b_{f}$ is of form (1.5.7).

Now if $x \in g$ then for the centralizer $g^{x}$ of $x$ one knows that $\operatorname{dim} g^{x} \geqslant l$ and the set $R=\left\{x \in g \mid \operatorname{dim} g^{x}=l\right\}$ is an open dense subset of $\mathfrak{g}$. The elements of $R$ are called regular.

Lemma 2.1.1. One has

$$
\begin{equation*}
b_{f} \subseteq R \tag{2.1.7}
\end{equation*}
$$

Proof. It clearly suffices, by complexification, to assume that $F=\mathbb{C}$. But then Lemma 2.1.1 is Lemma 10 in [14, p. 370]. Q.E.D.

If $y \in \mathscr{C}_{f}$ then it is not in general easy to describe the centralizer $\mathscr{g}^{y}$ of $y$. However, one has

Lemma 2.1.2. For any $y \in \mathfrak{b}_{f}$ one has

$$
\begin{equation*}
z^{y} \cap b=0 \tag{2.1.8}
\end{equation*}
$$

Proof. It clearly suffices to assume that $F=\mathbb{C}$. But then Lemma 2.1.2 is $(1.2 .4)$ in [17, p. 109].
Q.E.D.

Now

$$
\begin{equation*}
g=\bar{b} \oplus n \tag{2.1.9}
\end{equation*}
$$

is a $Q_{*}$-orthogonal direct sum since $(\bar{b})^{\perp}=\left((\bar{b})^{0}\right)^{*}=(\bar{n})^{*}=n$. Thus the decomposition (1.2.9) becomes $\delta u=\delta_{\bar{\beta}} u+\delta_{n} u$ for any $u \in S(g)$.

We recall that $I_{j}, j=1, \ldots, l$, are the generators of $S(g)^{G}$.
Proposition 2.1. Let $y \in \mathscr{E}_{f}$ be arbitrary and put

$$
\left(\delta_{b}\left(S(g)^{G}\right)\right)(y)=\left\{\left(\delta_{b} I\right)(y) \mid I \in S(\mathscr{g})^{G}\right\}
$$

Then the elements $\left(\delta_{\delta} I_{j}\right)(y), j=1, \ldots, l$, are linearly independent and are a basis of $\left(\delta_{\ell}^{-}\left(S(g)^{G}\right)\right)(y)$. Furthermore $[n, y] \subseteq \ell$ and one has

$$
\begin{equation*}
\left.Q\left(\left(\delta_{\bar{d}}(S(g))^{G}\right)\right)(y),[n, y]\right)=0 . \tag{2.1.10}
\end{equation*}
$$

In fact $[n, y]$ has codimension $l$ in $\&$ and is the Q-orthocomplement of $\left(\delta_{\bar{\delta}}\left(S(g)^{G}\right)\right)(y)$ in $\ell$.

Proof. Assume that there is a nontrivial relation $\sum c_{i}\left(\delta_{\bar{\beta}} I_{i}\right)(y)=0$. This implies that if $w=\sum c_{i}\left(\delta I_{i}\right)(y)$ then $w \in n$ by (2.1.9). But $w \in g^{y}$ by (1.3.3). However, $g^{y} \cap n=0$ by (2.1.8). Thus $w=0$. But $y \in R$ by (2.1.7). This contradicts Theorem 9 in [14, p. 382], which asserts that $\left(\delta I_{i}\right)(y), i=1, \ldots, l$, are linearly independent for any $y \in R$. Thus the $\left(\delta_{\bar{G}} I_{j}\right)(y)$ are linearly independent. From the differentiation properties of the map $I \rightarrow\left(\delta_{\bar{I}} I\right)(y)$ it is clear that they span $\left(\delta_{\bar{\sigma}}\left(S(g)^{G}\right)\right)(y)$. This proves the first statement of the proposition.

Now the minimal diagonal degree (see Definition 2.1) of $y$ is -1 . Thus $[n, y] \subseteq b$ by (2.1.4), (2.1.5), and (2.1.6). Thus if $I \in S(g)$ one has $Q\left(\left(\delta_{n} I\right)(y)\right.$, $[n, y])=0$ since $f^{\circ}=n$. However, $Q((\delta I)(y),[n, y])=0$ by the invariance of $Q$ and the fact that $[(\delta I)(y), y]=0$. See (1.3.3). Thus, subtracting, $Q((\delta \sigma I)(y),[n, y])=0$. This proves (2.1.10). But now (2.1.8) implies that $\operatorname{dim}[n, y]=\operatorname{dim} n$. However, $\operatorname{dim} b=l+\operatorname{dim} n$. Also, $b$ and $b$ are nonsingularly paired by $Q$ and we have shown that $\operatorname{dim}\left(\delta_{\bar{A}}(S(g))^{G}\right)(y)=l$. This proves the last statement of the proposition. Q.E.D.
2.2. Now it is clear that the diagonal $d_{1}$ is the $l$-dimensional subspace given by

$$
\begin{equation*}
d_{1}=\sum_{i=1}^{l} F e_{\alpha_{i}} \tag{2.2.1}
\end{equation*}
$$

Let

$$
\begin{equation*}
\tilde{d_{1}}=\left\{\sum_{i=1}^{l} a_{i} e_{\alpha_{i}} \mid \text { all } a_{i} \neq 0 \text { if } F=\mathbb{C} \text {; all } a_{i}>0 \text { if } F=\mathbb{R}\right\} \tag{2.2.2}
\end{equation*}
$$

Remark 2.2.1. If $F=\mathbb{C}$ note that $\tilde{d}_{1}$ is the set of all principal nilpotent elements in $d_{1}$ by Theorem 5.3 in [12].

Now put

$$
\begin{equation*}
Z_{o}=h+\tilde{d_{1}} \tag{2.2.3}
\end{equation*}
$$

so that $Z_{o}$ is a $2 l$-dimensional submanifold of $\ell$. Our primary concern is with the $2 l$-dimensional submanifold $Z$ of $\ell_{p}$ obtained by translating $Z_{o}$ by $f$. That is,

$$
\begin{equation*}
Z=f+Z_{o} \subseteq \mathscr{b}_{f} \tag{2.2.4}
\end{equation*}
$$

Let $\bar{B} \subseteq G$ be the subgroup corresponding to $\mathscr{C}$.
Remark 2.2.2. Note that $\bar{B}$ in the notation of Remark 1.6 .2 plays the role
of $A$. We shall not be concerned with this at present but it is easy to prove (see Section 6.4) that $Z$, in the sense of Remark 1.6.2, is an orbit of $\bar{B}$ in $\ell_{f}$.

For the case of $S l(n, F)$ note that choices can be made so that $Z$ is the set of all traceless Jacobi matrices of the form

$$
y=\left(\begin{array}{ccccc}
b_{1} & a_{1} & & \ldots &  \tag{2.2.5}\\
1 & b_{2} & a_{2} & & \\
\vdots & 1 & \cdot & \cdot & \vdots \\
\vdots & \cdot & \cdot & \cdot & \cdot \\
0 & \ldots & \cdot & a_{n-1} \\
0 & b_{n}
\end{array}\right)
$$

where $a_{i}, b_{i} \in F$ and $a_{i} \neq 0$ if $F=\mathbb{C}, a_{i}>0$ if $F=\mathbb{R}$.
The notion of Jacobi matrix generalizes to $g$. Let

$$
\begin{align*}
& \mathrm{Jac}_{g}=\left\{x+\sum_{i=1}^{l}\left(a_{-i} e_{-\alpha_{i}}+a_{i} e_{\alpha_{i}}\right) \mid x \in \ell, a_{-i} a_{i} \neq 0\right. \\
&\text { if } \left.F=\mathbb{C} ; a_{-i}, a_{i}>0 \text { if } F=\mathbb{R}\right\} . \tag{2.2.6}
\end{align*}
$$

Clearly one has

$$
\begin{equation*}
Z \subseteq \mathrm{Jac} g \tag{2.2.7}
\end{equation*}
$$

Of course the definition of $\mathrm{Jac} g$ depends upon the choice of positive and negative simple root vectors. Also, it appears at first in considering $Z$ that we would be restricting ourselves in any investigation of Jac $g$. This, however, is not the case. We are in fact just factorizing out the trivial action of $H$ on Jac $g$. That is, one has

Proposition 2.2. The map

$$
\begin{equation*}
H \times Z \rightarrow \mathrm{Jac} g, \quad(a, x) \rightarrow a x \tag{2.2.8}
\end{equation*}
$$

is bijective.
Proof. It is clear that Jac $g$ is stable under the action of $H$. On the other hand since $G$ is the adjoint group the map

$$
\begin{equation*}
H \rightarrow\left(F^{*}\right)^{l}, \quad h \rightarrow\left(h^{\alpha_{1}}, \ldots, h^{\alpha_{l}}\right) \tag{2.2.9}
\end{equation*}
$$

is bijective. It follows that $h$ leaves no element of Jac $g$ fixed in case $h \neq 1$. It also follows (replacing $\alpha_{i}$ by $-\alpha_{i}$ ) that every $H$-orbit in Jac $g$ contains a unique element in $Z$. This proves the proposition.
Q.E.D.

We refer to $Z$ as the manifold of normalized Jacobi elements.

Now if $y \in Z$ the tangent space to $Z$ at $y$ is clearly given by

$$
\begin{equation*}
T_{y}(Z)=h+d_{1} \tag{2.2.10}
\end{equation*}
$$

On the other hand for any $y \in \mathscr{\ell}=a^{*}$ we have defined a subspace $a^{*}(y)$ now written as $\ell(y)$, by (1.6.9).

Lemma 2.2.1. For any $y \in Z$ one also has

$$
\begin{equation*}
\ell(y)=h+d_{1} \tag{2.2.11}
\end{equation*}
$$

Proof. The projection $P_{a^{*}}$ of Section 1.6 is just $P_{\delta}$ and the kernel of $P_{\phi}$ is $\bar{n}$ since $\ell^{\perp}=\left(\ell^{0}\right)^{*}=n^{*}=\bar{n}$. By definition, $\ell(y)=P_{6}[\mathscr{b}, y]$. But since the maximal diagonal degree of $y$ is 1 (see Definition 2.1) one has $\ell(y) \subseteq h+d_{1}$ by (2.1.4) and (2.1.5). On the other hand since $\left[e_{-\alpha_{i}}, e_{\alpha_{j}}\right]=0$ if $i \neq j$ and the elements $\left[e_{-\alpha_{i}}, e_{\alpha_{i}}\right]$ are a basis of $h$ it follows that $P_{\delta}\left[d_{-1}, y\right]=d_{0}=h$. But also since $\left[d_{0}, d_{0}\right]=0$ one clearly has $P_{\delta}\left[d_{0}, y\right]=d_{1}$. Thus $h+d_{1} \subseteq \mathscr{G}(y)$. This proves the lemma.
Q.E.D.

Remark 2.2.3. Recalling Remarks 1.6.1 and 2.2.2 note that Lemma 2.2.1 is anticipated by the fact that $Z$ is an orbit of $\bar{B}$ in $\ell$.

Now for any $y \in Z$ note that

$$
\begin{equation*}
[y,[\bar{n}, \bar{n}]] \subseteq \bar{n} . \tag{2.2.12}
\end{equation*}
$$

This is clear since the maximal diagonal degree of $y$ is 1 and the maximal diagonal degree of any nonzero element in $[\bar{n}, \bar{n}]$ is -2 . Now for $y \in Z$ let

$$
\begin{equation*}
n_{(y)}=\{x \in n \mid Q(x,[y,[\bar{n}, \bar{n}]])=0\} . \tag{2.2.13}
\end{equation*}
$$

Lemma 2.2.2. For any $y \in Z$ one has $\operatorname{dim} n_{(y)}=$ l. Furthermore $\operatorname{dim}\left(h+d_{1}\right) \cap$ $[n, y]=l$ and one has

$$
\begin{equation*}
\left[n_{(y)}, y\right]=\left(h+d_{1}\right) \cap[n, y] . \tag{2.2.14}
\end{equation*}
$$

Proof. One first has that

$$
\begin{equation*}
g^{y} \cap \vec{b}=0 . \tag{2.2.15}
\end{equation*}
$$

Indeed if $e \in \tilde{d_{1}}$ is the maximal diagonal component of $y$ then $y \in e+\tilde{\ell}$. On the other hand no normalization was made with regard to the simple root vector so that we can interchange the roles of positive and negative roots where
$\mathscr{f}$ replaces $f$ and $e$ replaces $f$. But then (2.2.15) follows from (2.1.8). But now one knows

$$
\begin{equation*}
[x, \bar{x}]=\sum_{j \leqslant-2} d_{j} \tag{2.2.16}
\end{equation*}
$$

so that $[\bar{n}, \bar{n}]$ has codimension $l$ in $\bar{x}$. However, ad $y$ has no kernel in $[\bar{n}, \bar{n}] \subseteq \bar{b}$ by (2.2.15) and hence $[y,[\bar{x}, \bar{x}]]$ also has codimensional $l$ in $\bar{x}$, recalling (2.2.12). This proves $\operatorname{dim} n_{(y)}=l$. But then by (2.1.8) one also has

$$
\begin{equation*}
\operatorname{dim}[n(v), y]=l \tag{2.2.17}
\end{equation*}
$$

Again by (2.1.8) any element in $[n, y]$ is uniquely of the form $[x, y]$ for $x \in n$. On the other hand $h+d_{1}$ by (2.2.16) is clearly the $Q$-orthocomplement of $[\bar{n}, \bar{x}]$ in $\ell$. Thus $[x, y]$ is in the right side of (2.2.14) if and only if $Q([x, y]$, $[\bar{n}, \bar{n}])=0$. Therefore by the invariance of $Q$ the element $[x, y]$ is in the right side of (2.2.14) if and only if $x \in n_{(y)}$, using the definition (2.2.13) of $n_{(y)}$. This proves (2.2.14) and the lemma follows from (2.2.17). Q.E.D.

Now recall (see (1.6.5)) that $\eta_{I}$ for any invariant $\left.I \in S(g)\right)^{G}$ is the vector field on $f_{f}$ such that

$$
\begin{equation*}
\left(\eta_{I}\right)_{y}=\left[y,\left(\delta_{\Delta} I\right)(y)\right] . \tag{2.2.18}
\end{equation*}
$$

By (2.2.10), (2.2.11), and (1.6.10), $\eta_{I}$ is tangent to the submanifold $Z$. Let

$$
\begin{equation*}
\xi_{I}=\eta_{I} \mid Z \tag{2.2.19}
\end{equation*}
$$

and let $x$ be the space of vector fields on $Z$ defined by

$$
\begin{equation*}
z=\left\{\xi_{I} \mid I \in S(g)^{\sigma}\right\} \tag{2.2.20}
\end{equation*}
$$

Also for any $y \in Z$, let

$$
\begin{equation*}
z_{y}=\left\{\left(\xi_{l}\right)_{y} \mid I \in S(g)^{G}\right\} . \tag{2.2.21}
\end{equation*}
$$

Theorem 2.2. Let $g$ be any complex or real split semi-simple Lie algebra. Let $l=\operatorname{rank}_{g}$ and let $Z$ be the $2 l$-dimensional manifold of normalized Jacobi elements in g. See (2.2.4). For any invariant $I \in S(g)^{G}$ let $\xi_{I}$ be the vector field on $Z$ defined by putting

$$
\begin{equation*}
\left(\xi_{I}\right)_{v}=\left[y,\left(\delta_{6} I\right)(y)\right] \tag{2.2.22}
\end{equation*}
$$

for any $y \in Z$. The element $\left(\delta_{\bar{\beta}} I\right)(y)$ is given by (1.2.10).

Now let $z$ be the space spanned by $\xi_{I}$ for all $I \in S(\mathscr{g})^{G}$. Then $z$ is a commutative Lie algebra of vector fields. Furthermore for any $y \in Z$ the subspace $z_{y} \subseteq T_{y}(Z)$ defined by (2.2.21) is $l$-dimensional and the elements $\left(\xi_{I_{j}}\right)_{y}, j=1, \ldots, l$, are a basis of $z_{y}$, where the $I_{j} \in S(g)^{G}$ are the fundamental invariants. Moreover if $n_{(y)} \subseteq n$ is defined by (2.2.13) then the elements $\left(\delta_{n} I_{j}\right)(y), j=1, \ldots, l$, are a basis of $n(y)$ and

$$
\begin{equation*}
\check{n}_{y}=\left[n_{(y)}, y\right] . \tag{2.2.23}
\end{equation*}
$$

Proof. It follows immediately from Theorem 1.6.2 that $\approx$ is a commutative Lie algebra of vector fields. Now let $y \in Z$. By Proposition 2.1, $\left(\delta_{\bar{\delta}} I\right)(y)$ is in the span of $\left(\delta_{\bar{\sigma}} I_{j}\right)(y), j=1, \ldots, l$, for any $I \in S(g)^{G}$. Thus the elements $\left(\delta_{I_{i}}\right)_{v}$ span $z_{y}$. But also by Proposition 2.1 the elements $\left(\delta_{\bar{a}} I_{j}\right)(y)$ are linearly independent. Furthermore $g^{\nu} \cap \delta=0$ by (2.2.15). Thus the elements $\left[y,\left(\delta_{\overline{\mathcal{A}}} I_{j}\right)(y)\right]=$ $\left(\xi_{I}\right)_{y}$ are also linearly independent and hence are a basis of $z_{y}$. . $\operatorname{But}[y,(\delta I)(y)]=$ 0 by (1.3.3). Thus for any $I \in S(g)^{G}$

$$
\begin{align*}
\left(\xi_{I}\right)_{y} & =\left[y,\left(\delta_{\bar{z}} I\right)(y)\right] \\
& =\left[\left(\delta_{n} I\right)(y), y\right] . \tag{2.2.24}
\end{align*}
$$

But $\left(\xi_{I}\right)_{y} \in T_{y}(Z)=\hbar+d_{1}$. Thus by (2.2.24), $\left(\xi_{1}\right)_{y} \in\left(h+d_{1}\right) \cap[n, y]$. Thus $\left[\left(\delta_{n} I\right)(y), y\right] \in\left[n_{(v)}, y\right]$ by (2.2.14). But since $g^{y} \cap_{n}=0$ by (2.1.8) this implies $\left(\delta_{n} I\right)(y) \in n_{(y)}$. However, by (2.2.24) the elements $\left(\delta_{n} I_{j}\right)(y)$ must be linearly independent. But then they are a basis of $n(y)$ by Lemma 2.2.2. Furthermore (2.2.24) then immediately implies (2.2.23).
Q.E.D.
2.3. Now the fundamental invariants $I_{j}$ define a differentiable map

$$
\begin{equation*}
\mathscr{I}: g \rightarrow F^{l}, \tag{2.3.1}
\end{equation*}
$$

where $\mathscr{I}(x)=\left(I_{1}(x), \ldots, I_{l}(x)\right)$.
For each $\gamma$ in the image, $\mathscr{I}(Z)$, of $Z$ under $\mathscr{I}$ one defines a nonempty closed subset $Z(\gamma)$ of $Z$ by putting

$$
\begin{equation*}
Z(\gamma)=\mathscr{I}^{-1}(\gamma) \cap Z \tag{2.3.2}
\end{equation*}
$$

One of course then has the disjoint union

$$
\begin{equation*}
Z=\bigcup_{\gamma \in \mathscr{F}(Z)} Z(\gamma) . \tag{2.3.3}
\end{equation*}
$$

Remark 2.3. For the case of $S l(n, F)$ note that the equivalence relation in $Z$ defined by regarding (2.3.3) as a union of equivalence classes is the relation
$x \sim y$ when $x$ and $y$ are normalized Jacobi matrices with the same characteristic polynomial.

Consider the functions $I_{i} \mid Z$ on $Z$ defined by restricting the fundamental invariants to $Z$. The following guarantees that the $Z(\gamma)$ are submanifolds, although at this point it is not clear whether or not they are connected.

Proposition 2.3.1. For any $y \in Z$ the differentials $\left(d\left(I_{i} \mid Z\right)\right)_{y}, i=1, \ldots, l$, are linearly independent.

Proof. Now as in (2.2.10) the tangent space to $Z$ at $y$ is identified with $h+d_{1}$. With respect to this identification if $w \in h+d_{1}$ then by (1.1.2) and definition (1.2.2) one has

$$
\begin{equation*}
\left\langle w,(d(I \mid Z))_{v}\right\rangle=Q(w,(\delta I)(y)) . \tag{2.3.4}
\end{equation*}
$$

Thus $(d(I \mid Z))_{y}=0$ if and only if $(\delta I)(y)$ is $Q$-orthogonal to $h+d_{1}$. But $\hbar+d_{1}=\ell(y)$ by (2.2.11). Thus $(d(I \mid Z))_{y}=0$ if and only if $\left(\xi_{I}\right)_{y}=0$ by (1.6.12). However, the vectors $\left(\xi_{I_{i}}\right)_{y}$ are linearly independent by Theorem 2.2. Thus the differentials $\left(d\left(I_{i} \mid Z\right)\right)_{y}$ are then also linearly independent. Q.E.D.

Let $\bar{N}, N \subseteq G$ be the subgroups respectively corresponding to $\bar{n}$ and $n$. We return briefly to $\ell_{f}=f+\ell$.

Proposition 2.3.2. Let $y_{1}, y_{2} \in \mathfrak{b}_{f}$. Then there exists $n \in N$ such that $n y_{1}=y_{2}$ if and only if $\mathscr{I}\left(y_{1}\right)=\mathscr{I}\left(y_{2}\right)$. Furthermore in such a case $n$ is unique.

Proof. If such an element $n \in N$ exists then by invariance one has $\mathscr{I}\left(y_{1}\right)=$ $\mathscr{I}\left(y_{2}\right)$. Conversely assume $\mathscr{I}\left(y_{1}\right)=\mathscr{I}\left(y_{2}\right)$. We must show there exists a unique $n \in N$ such that $n y_{1}=y_{2}$. Assume first that $F=\mathbb{C}$.

Let $x_{o} \in \ell$ be as in Section 2.1. It is clear from (2.1.4) that $[f, n] \subseteq \mathscr{b}$ and $[f, n]$ is stable under ad $x_{o}$. Let $\varsigma \subseteq \mathscr{b}$ be any ad $x_{o}$ stable subspace such that

$$
b=[f, n]+0
$$

is a direct sum. Now by Theorem 1.2 in [17] there exist elements $n_{i} \in N$, $x_{i} \in f+\sigma$ such that $n_{i} x_{i}=y_{i}, i=1,2$. Clearly $\mathscr{I}\left(x_{i}\right)=\mathscr{I}\left(y_{i}\right)$. However, since $\mathscr{I}\left(y_{1}\right)=\mathscr{I}\left(y_{2}\right)$ one then has $\mathscr{I}\left(x_{1}\right)=\mathscr{I}\left(x_{2}\right)$. But then $x_{1}=x_{2}$ by Theorem 7 in [14, p. 381]. (See also Remark 19' in [14, p. 375]. This proves there exists $n \in N$ such that $n y_{1}=y_{2}$. The uniqueness of $n$ follows from Theorem 1.2 in [17].

Now assume $F=\mathbb{R}$. Let $g_{\mathbb{C}}=g+i g$ and let $G_{\mathbb{C}}$ be the adjoint group of $g_{\mathbb{C}}$. Let $n_{\mathbb{C}}=n+i n$ and let $N_{\mathbb{C}} \subseteq G_{\mathbb{C}}$ be the subgroup corresponding to $n_{\mathbb{C}}$. By the result above there exists a unique $n \in N_{\mathbb{C}}$ such that $n y_{1}=y_{2}$. Now regarding $N \subseteq N_{\mathbb{C}}$ we have only to show that $n \in N$. Write $n=\exp z$,
where $z \in n_{\mathbb{C}}$, and write $z=u+i v$, where $u, v \in n$. It is enough to show that $v=0$. Assume not and let $j$ and $v_{j}$ be respectively the minimal diagonal degree of $v$ and the minimal diagonal component of $v$. Then $j \geqslant 1$ since $v \in u$ and hence

$$
\begin{equation*}
\left[v_{j}, f\right] \in d_{j-1} \subseteq \ell \tag{2.3.5}
\end{equation*}
$$

But now

$$
\begin{equation*}
\left(\sum_{m} \frac{\left(\operatorname{ad}(u+i v)^{m}\right)}{m!}\right)\left(f+x_{1}\right)=y_{2} \tag{2.3.6}
\end{equation*}
$$

where we have written $y_{1}=f+x_{1}$ so that $x_{1} \in \ell$.
Now the minimal diagonal degree of any nonzero element in $[\tau, 0]$ is at least $j$ and any further bracketing with $u$ or $v$ would only increase the degree. Thus upon writing $g_{\mathbb{C}}=\sum d_{k}+i \sum d_{k}$ as a real direct sum it follows that the component of the left side of (2.3.6) in $i d_{j-1}$ is just $i\left[v_{j}, f\right]$. But $\left[v_{j}, f\right] \neq 0$ by (2.1.8). This contradicts the fact that $y_{2} \in g$. Thus $v=0$. Q.E.D.
2.4. Here and in Sections 2.5 and 2.6 we assume that $F=\mathbb{C}$. Now the Weyl group $W=W(g, h)$ operates on $\hbar$ and on its dual $\ell^{\prime}$. For each $\sigma \in W$ let $s(\sigma) \in G$ be the unique element in the normalizer of $h$, corresponding to $\sigma$, such that

$$
\begin{equation*}
s(\sigma) e_{\alpha_{i}}=e_{o \alpha_{i}}, \quad i=1, \ldots, l . \tag{2.4.1}
\end{equation*}
$$

The Bruhat-Gelfand decomposition of $G$ asserts that

$$
\begin{equation*}
G=\bigcup_{\sigma \in W} N s(\sigma) H N \tag{2.4.2}
\end{equation*}
$$

is a disjoint union. Now let $\kappa \in W$ be the unique element such that $\kappa \Delta_{+}=-\Delta_{+}$. Then $s(\kappa) N s(\kappa)^{-1}=\bar{N}$. Multiplying the components of (2.4.2) on the left by $s(\kappa)$ one also has the disjoint union

$$
\begin{equation*}
G=\bigcup_{\tau \in W} \bar{N}_{s}(\tau) H N \tag{2.4.3}
\end{equation*}
$$

Now if $\tau$ is the identity then $\overline{N s}(\tau) H N$ is just $\bar{N} H N$ Now put

$$
\begin{equation*}
G_{*}=\bar{N} H N \tag{2.4.4}
\end{equation*}
$$

One knows that $G_{*}$ is a Zariski open subset of $G$ and the map

$$
\begin{equation*}
\bar{N} \times H \times N \rightarrow G_{*} \tag{2.4.5}
\end{equation*}
$$

where ( $\bar{n}, h, n$ ) $\rightarrow \bar{n} h n$ is an isomorphism of algebraic varieties. Thus for any $d \in G_{*}$ there exist unique elements $\bar{n}_{d} \in \bar{N}, h_{d} \in H$, and $n_{d} \in N$ such that

$$
\begin{equation*}
d=\bar{n}_{d} h_{d} n_{d} \tag{2.4.6}
\end{equation*}
$$

We will retain this notation throughout. One notes in particular that by inverting (2.4.5) the map

$$
\begin{equation*}
G_{*} \rightarrow N, \quad d \mapsto n_{d} \tag{2.4.7}
\end{equation*}
$$

is a morphism of algebraic varieties.
Now for any $y \in G$ we recall that $G^{y}$ is the centralizer of $y$ in $G$. Thus $G^{v}$ is a Zariski closed algebraic subgroup of $G$. In case $y$ is regular, i.e., $y \in R$, this group is of a particularly simple form.

Proposition 2.4. If $y \in R$ then $G^{y}$ is an Abelian connected algebraic group of complex dimension l. That is, as algebraic groups one has the isomorphism

$$
\begin{equation*}
G^{y} \cong\left(\mathbb{C}^{*}\right)^{p} \times \mathbb{C}^{q} \tag{2.4.8}
\end{equation*}
$$

for some $p, q$ where $p+q=l$.
Proof. The first statement is contained in Proposition 14 in [14, p. 362]. The second statement follows then from the well-known classification of connected Abelian complex algebraic groups.
Q.E.D.

Now let $y \in Z$. 'I'hen $y \in R$ so that Proposition 2.4 applies to $y$. Let

$$
\begin{equation*}
G_{*}^{y}=G^{y} \cap G_{*} \tag{2.4.9}
\end{equation*}
$$

so that $G_{*}^{y}$ is a Zariski open subset of $G^{v}$. Since $1 \in G_{*}^{y}$ it follows that $G_{*}^{y}$ is not empty and hence is dense in $G^{y}$. Furthermore since $G^{y}$ is connected and nonsingular it follows that $G_{\star}^{y}$ is a connected, nonsingular algebraic variety of dimension $l$. We recall that a variety of $\operatorname{dim} l$ is called rational if its function field is isomorphic to $\mathbb{C}\left(X_{1}, \ldots, X_{l}\right)$, where the $X_{i}$ are indeterminates. Since $G^{y}$ is an algebraic group it is rational. (This is of course obvious from (2.4.8).) Since $G_{*}^{y}$ is Zariski dense in $G^{y}$ it has the same function field so that it too is rational.

Now for any $d \in G_{*}^{y}$ one has the decomposition (2.4.6). That is, $d=n_{d} h_{d} \bar{n}_{d}$ and by restriction (2.4.7) the map

$$
\begin{equation*}
G_{*}^{y} \rightarrow N, \quad d \rightarrow n_{d} \tag{2.4.10}
\end{equation*}
$$

is a morphism of algebraic varieties.
Now if $\gamma$ is the image $\mathscr{I}(Z)$ of the map $\mathscr{I}$ it is clear that the "isospectral" set $Z(\gamma)$ is a closed subvariety of the variety $Z$ of normalized Jacobi elements. However, it is not clear just what variety it is and in fact whether or not it is connected. These questions are settled by

Theorem 2.4. Let $g$ be a complex semi-simple Lie algebra and let $Z$ be the
set of normalized Jacobi elements defined by (2.2.4). Let $y \in Z, \gamma=\mathscr{F}(y)$, and $Z(\gamma) \subseteq Z$ be defined by (2.3.2) (so that $Z(\gamma)$ is the set of all elements in $Z$ which are conjugate to $y$ ).

Now let $G$ be the adjoint group of g. Let $G^{y}$ be the centralizer of $y$ in $G$ and let $G_{*}^{y}$ be the Zariski open subset of $G^{y}$ defined by (2.4.9). Then for any $d \in G_{*}^{y}$ one has $n_{d} y \in Z(\gamma)$. (See (2.4.6).) Furthermore if $\beta_{y}(d)=n_{d} y$ then

$$
\begin{equation*}
\beta_{y}: G_{*}^{y} \rightarrow Z(\gamma) \tag{2.4.11}
\end{equation*}
$$

is an isomorphism of algebraic varieties. In particular $Z(\gamma)$ is an irreducible, nonsingular, rational algebraic variety of dimension $l$ where $l=\mathrm{rank} g$ (so that $Z(\gamma)$ is a connected complex manifold of dimension $l$.

Proof. Now since $y \in Z \subseteq b_{f}$ and $l_{f}$ is stable under the action of $N$ it is clear that $\beta_{y}(d) \in \mathscr{b}_{f}$. Let $w=\beta_{y}(d)$ and let $j$ and $w_{j}$ be respectively the maximal diagonal degree and component of $w$. See Definition 2.1. To prove that $w \in Z$ it suffices then to show that $j=1$ and $w_{j} \in \tilde{d_{1}}$. Now clearly $n_{d}=h^{-1}(\bar{n})^{-1} d$, where we have written $\bar{n}=\bar{n}_{d}, h=h_{d}$. Thus $w=h^{-1}(\bar{n})^{-1} d y=h^{-1}(\bar{n})^{-1} y$ since $d y=y$. That is,

$$
\begin{equation*}
\beta_{y}(d)=h_{d}^{-1}\left(\bar{n}_{d}\right)^{-1} y . \tag{2.4.12}
\end{equation*}
$$

But if $y_{1}$ is the maximal diagonal component of $y$ then $y_{1} \in \tilde{d_{1}}$. However, clearly the maximal diagonal component is unchanged by the action of any element in $\bar{N}$. Thus $y_{1}$ is also the maximal diagonal component of $(\bar{n})^{-1} y$. But now obviously $h^{-1} y_{1}$ is the maximal diagonal component of $h^{-1}(\bar{n})^{-1} y=w$. That is, there exists $x \in h$ so that

$$
\begin{equation*}
\beta_{y}(d)=f+x+h_{d}^{-1} y_{\mathbf{1}} \tag{2.4.13}
\end{equation*}
$$

Since $\tilde{d_{1}}$ is stable under $H$ this proves that $\beta_{y}(d) \in Z$. However, then certainly $\beta_{y}(d)=n_{d} y \in Z(\gamma)$ by the invariance of $\mathscr{I}$. This defines the map (2.4.11). Assume now that $\beta_{y}\left(d_{1}\right)=\beta_{y}\left(d_{2}\right)$, where $d_{1}, d_{2} \in G_{*}^{y}$. That is, $n_{1} y=n_{2} y$, where $n_{i}=n_{d_{i}}, i=1,2$. But then by the uniqueness in Proposition 2.3.2 one has $n_{1}=n_{2}$. But then $h_{1}^{-1} y_{1}=h_{2}^{-1} y_{1}$ by (2.4.13), where $h_{i}=h_{d_{i}}$. Since $H$ operates in a simple transitive way on ${\tilde{d_{1}}}_{1}$ it follows that $h_{1}=h_{2}$. Finally if $\bar{n}_{i}=\bar{n}_{d_{i}}$ one has $\left(\bar{n}_{1}\right)^{-1} y=\left(\bar{n}_{2}\right)^{-1} y$ by (2.4.12). But then by Proposition 2.3 .2 with the roles of $\ell$ and $\mathscr{\ell}, N$ and $\bar{N}, f$ and $y_{1}$ reversed, one has $\bar{n}_{1}=\bar{n}_{2}$ by uniqueness since $y \in y_{1}+\ell$. This proves $d_{1}=d_{2}$ so that $\beta_{y}$ is injective.

Now let $w \in Z(\gamma)$. Since $w, y \in \mathscr{b}_{f}$ this implies there exists $n \in N$ by Proposition 2.3.2 such that

$$
\begin{equation*}
n y=w \tag{2.4.14}
\end{equation*}
$$

Now let $y_{1}$ and $w_{1}$ be respectively the maximal diagonal components of $y$ and $w$.

Then $y_{1}, w_{1} \in \tilde{d_{1}}$ and hence by the transitivity of $H$ on $\tilde{d}_{1}$ there exists $h \in H$ such that $h w_{1}=y_{1}$. Thus $h w, y \in y_{1}+\zeta$ and $\mathscr{I}(h w)=\mathscr{I}(y)=\gamma$ by invariance of $\mathscr{I}$ under conjugation. But then by Proposition 2.3.2 with the roles reversed as above (in the proof of injectivity) there exists $\bar{n} \in N$ such that

$$
\begin{equation*}
\bar{n} h w=y . \tag{2.4.15}
\end{equation*}
$$

Substituting for $w$ using (2.4.14) one has $\bar{n} h n y=y$. Thus $\bar{n} h n=d \in G^{y}$ and $n=n_{d}$. Thus $w=\beta_{y}(d)$ by (2.4.14). This proves the surjectivity of $\beta_{y}$ and hence $\beta_{y}$ is bijective. On the other hand $\beta_{y}$ is clearly a morphism of algebraic varieties since (2.4.10), as noted, is such a morphism. Thus $Z(\gamma)$ is an irreducible variety. But, by Proposition 2.3.1, $Z(\gamma)$ is certainly nonsingular, using, e.g., Zariski's criterion. Thus $Z(\gamma)$ is a nonsingular irreducible variety. In particular it is normal. But then since $\beta_{y}$ is bijective it is an isomorphism of algebraic varieties by the Zariski main theorem.
Q.E.D.

Remark 2.4.1. We thank D. Mumford and D. Kazhdan for pointing out that our isomorphism (2.4.11) implies the rationality of $Z(\gamma)$. For the special case of $S l(n)$ and where $y$ has distinct eigenvalues this result is due to J. Moser. The rationality in that case was established by Moser, using continued fractions.

Remark 2.4.2. Since $Z(\gamma)$ is a nonsingular variety it also has the structure of a smooth manifold. Moreover as such it is a submanifold of $Z$. Furthermore since $\beta_{y}$ is an isomorphism of algebraic varieties one knows then that $\beta_{y}$ is also a diffeomorphism of smooth manifolds. See, e.g., Chapter VII, Section 1, in [22].
2.5. As in Section 2.4 we assume that $F=\mathbb{C}$. Now one can easily determine the image $\mathscr{I}(Z)$ of $Z$ under the map $\mathscr{I}$. See (2.3.1). Let $e \in d_{1}$ and put $g=f+e$. Let

$$
\begin{equation*}
\mathscr{I}_{g}: h \rightarrow \mathbb{C}^{l} \tag{2.5.1}
\end{equation*}
$$

be the map defined by putting $\mathscr{g}_{g}(x)=\mathscr{I}(g+x)$ for $x \in h$. Recall that $W$ is the Weyl group $W(g, h)$. For any finite set $S$ let $|S|$ denote its cardinality.

Proposition 2.5.1. Let $\gamma \in \mathbb{C}^{l}$ be arbitrary. Then $\mathscr{I}_{g}^{-1}(\gamma)$ is finite and in fact

$$
\begin{equation*}
1 \leqslant\left|\mathscr{J}_{g}^{-1}(\gamma)\right| \leqslant|W| \tag{2.5.2}
\end{equation*}
$$

In particular the map $\mathscr{I}_{g}$ is surjective.
Proof. Let $J_{i} \in S(h)$ be defined by putting $J_{i}=I_{i} \mid h$. (We are of course identifying $S(h)$ with its dual using $Q \mid h$.) If $\operatorname{deg} I_{i}=m_{i}+1$ then of course $J_{i}$ is also homogeneous of degree $m_{i}+1$. On the other hand by Chevalley
the $J_{i}$ are algebraically independent and the algebra $\left.S(\ell)\right)^{W}$ of Weyl group invariants is exactly the polynomial algebra, $\mathbb{C}\left[J_{1}, \ldots, J_{l}\right]$, generated by the $J_{i}$. Furthermore it is also due to Chevalley that $S(h)$ is a free $S(h)^{W}$ with $|W|$ generators. (See, e.g., Theorem 4.15 .28 in [26, p. 386].) That is, there exists a graded subspace $D \subseteq S(h)$ of dimension $|W|$ such that $S(h)=S(h)^{W} \otimes D$. But then since $D$ is graded it is easy to see that

$$
\begin{equation*}
S(h)=A_{g} \otimes D \tag{2.5.3}
\end{equation*}
$$

where if $J_{i}{ }^{g} \in S(h)$ is defined by $J_{i}{ }^{g}(x)=J_{i}(g+x)$ for $x \in h$ then $A=$ $\mathbb{C}\left[J_{1}{ }^{g}, \ldots, J_{l}{ }^{g}\right]$. One needs only the obvious fact that $J_{i}-J_{i}{ }^{g}$ is a polynomial of degree at most $m_{i}$.

Thus $S(h)$ is a free $A_{g}$ module with $|W|$ generators. In particular $S(h)$ is integral over $A_{g}$. On the other hand $A_{g}$ is just the pullback of the affine algebra on $\mathbb{C}^{l}$ with respect to the map $\mathscr{I}_{g}$. Thus $\mathscr{I}_{g}$ is a finite map. See, e.g., [22, p. 48]. But then $\mathscr{I}_{g}$ is surjective (see, e.g., Theorem 4 in [22, p. 48]) so that $1 \leqslant\left|\mathscr{I}_{g}^{-1}(\gamma)\right|$ for any $\gamma \in \mathbb{C}^{l}$. On the other hand if $(A)$ denotes the quotient field of an integral domain $A$ then (2.5.3) easily implies $(S(\ell))=\left(A_{g}\right) \otimes D$ so that $|W|$ is the degree of the map $\mathscr{I}_{g}$. See, e.g., [22, p. 116]. But then one knows $\left|\mathscr{I}_{g}^{-1}(\gamma)\right| \leqslant W$ for any $\gamma \in \mathbb{C}^{l}$. See, e.g., Theorem 6 in [22, p. 116].
Q.E.D.

The set of unramified points is of course nonempty and Zariski open. That is,
Proposition 2.5.2. Let the notation be as in Proposition 2.5.1. Then the set

$$
\begin{equation*}
\mathbb{C}_{g}^{l}=\left\{\gamma \in \mathbb{C}^{l}| | \mathscr{F}_{g}^{-1}(\gamma)|=|W|\}\right. \tag{2.5.4}
\end{equation*}
$$

is a nonempty open subset of $\mathbb{C}^{l}$.
Proof. Since $|W|$ is the degree of the map $\mathscr{I}_{g}$ as noted in the proof of Proposition 2.5.1 the result is just Theorem 7 in [22, p. 117]. Q.E.D.

Summarizing one has
Theorem 2.5. Let $g$ be a complex semi-simple Lie algebra. Let $l=$ rank $g$ and let $Z$ be the 2l-dimensional variety of normalized Jacobi elements in $g$ defined by (2.2.4). For any $\gamma \in \mathbb{C}^{l}$ let $Z(\gamma) \subseteq Z$ be the subvariety defined by (2.3.2) in terms of the fundamental invariants $I_{j}, j=1, \ldots, l$. Then $Z(\gamma)$ is an l-dimensional nonsingular, irreducible closed rational subvariety of $Z$ and hence

$$
\begin{equation*}
Z=\bigcup_{\gamma \in \mathbb{C}^{\mathbf{l}}} Z(\gamma) \tag{2.5.5}
\end{equation*}
$$

is a foliation of $Z$ by such subvarieties.

Proof. Recall the notation and statement of Proposition 2.5.1. Since the map $\mathscr{I}_{g}$ is surjective it of course follows that $\mathscr{I}$ is surjective. That is,

$$
\begin{equation*}
\mathscr{I}(Z)=\mathbb{C}^{l} \tag{2.5.6}
\end{equation*}
$$

The result then follows from Theorem 2.4 and (2.3.3).
Q.E.D.

Henceforth we will speak of the sets $Z(\gamma)$ as the isospectral leaves of $Z$.
Remark 2.5. Although the elements of $Z$ are regular one notes that by (2.5.6) the elements of $Z$ are not in general semi-simple. In fact since two regular elements $x, y \in \mathscr{g}$ are $G$-conjugate if and only if $\mathscr{I}(x)=\mathscr{I}(y)$ (see (3.8.4) in [14]) it follows from (2.5.6) that every regular element is $G$-conjugate to a normalized Jacobi element. In particular taking $\gamma=0$ in $\mathbb{C}^{l}$ the isospectral leaf $Z(0)$ is the set of all nilpotent elements in $Z$. Since $Z \subseteq R$ any element in $Z(0)$ is in fact a principal nilpotent element.

Here again as in Sections 2.4 and 2.5 we assume $F=\mathbb{C}$.
2.6. Now if $\gamma \in \mathbb{C}^{l}$ then the isospectral leaf $Z(\gamma)$ of $Z$ has been described by Theorem 2.4 in terms of a Zariski open subset of the centralizer $G^{y}$ of an element $y \in Z$. However, the centralizer $G^{y}$ is itself fairly complicated. We wish now to describe $Z(\gamma)$ in terms of the centralizer $G^{w}$ of an element $w$ which is more readily accessible.

Now let $w_{0} \in h$ and put $w=f+w_{0}$. Thus $w \in b_{f}$ and hence $w \in R$ by (2.1.7). Thus $G^{w}$ by Proposition 2.4 is a connected Abelian subgroup of complex dimension $l$. Now since $G_{*}=\bar{N} H N$ is Zariski open in $G$ the same is true for its translate $s(\kappa) G_{*}$. Thus if we put

$$
\begin{equation*}
G_{(*)}=s(\kappa) G_{*} \tag{2.6.1}
\end{equation*}
$$

then $G_{(*)}$ is a Zariski open subset of $G$ and for any $g \in G_{(*)}$ there exist unique elements $\bar{n}(g) \in \bar{N}, h(g) \in H$, and $n(g) \in N$ such that

$$
\begin{equation*}
g=s(\kappa) \bar{n}(g) h(g) n(g) \tag{2.6.2}
\end{equation*}
$$

This notation as with (2.4.6) will be retained throughout. By inverting (2.4.5) it is clear that the map

$$
\begin{equation*}
G_{(*)} \rightarrow N, \quad g \rightarrow n(g) \tag{2.6.3}
\end{equation*}
$$

is a morphism of algebraic varieties

$$
\begin{equation*}
G_{(*)}^{w}=G^{w} \cap G_{(*)} \tag{2.6.4}
\end{equation*}
$$

so that $G_{(*)}^{w}$ is Zariski open in $G^{w}$.

Remark 2.6.1. Unlike the case of $G_{*}^{v}$ for $y \in Z$ it is not immediately obvious (at least to us) whether $G_{(*)}^{w}$ is empty or not. It is in fact nonempty as stated in Theorem 2.6 below.

Now for any $g \in G_{(*)}^{w}$ one has the decomposition (2.6.2). That is, $s(\kappa)^{-1} g=$ $\bar{n}(g) h(g) n(g)$. Furthermore by restricting (2.6.3) the map

$$
\begin{equation*}
G_{(*)}^{w} \rightarrow N, \quad g \rightarrow n(g) \tag{2.6.5}
\end{equation*}
$$

is a morphism of algebraic varieties.

Lemma 2.6.1. Let $\gamma=\mathscr{I}(w)$ and let $g \in G_{(*)}^{w}$. Then $n(g) w \in Z(\gamma)$.
Proof. Let $y=n(g) w$. Since $\ell_{f}$ is stable under the action of $N$ it is clear that $y \in b_{f}$. Let $j$ and $y_{j}$ be respectively the maximal diagonal degree and component of $y$. We have to show that $j=1$ and $y_{1} \in \tilde{d_{1}}$. But the maximal diagonal component of $\bar{n}(g) h(g) y$ is clearly just $h(g) y_{j}$. However, $\bar{n}(g) h(g) y=$ $\bar{n}(g) h(g) n(g) w=s(\kappa)^{-1} g w$. But $g w=w$ and $s(\kappa)^{-1} w=\kappa^{-1} w_{o}+e$, where, by (2.4.1) and (1.5.4), $e=\sum_{i=1}^{l} e_{\alpha_{i}}$. Thus $h(g) y_{j}=e$. Hence $j=1$ and $y_{1}=$ $h(g)^{-1} e$. That is,

$$
\begin{equation*}
y_{1}=\sum_{i=1}^{l} h(g)^{-\alpha_{i}} e_{\alpha_{i}} \tag{2.6.6}
\end{equation*}
$$

Note also that one has

$$
\begin{equation*}
\bar{n}(g) h(g) y=\kappa^{-1} w_{o}+e \tag{2.6.7}
\end{equation*}
$$

But now (2.6.6) proves $y_{1} \in \tilde{d}_{1}$ and hence $y \in Z$. Also, since $y$ and $w$ are $G$-conjugate one has $\mathscr{I}(y)=\mathscr{J}(w)=\gamma$ so that $y \in Z(\gamma)$. $\quad$ Q.E.D.

Now let the notation be as in Lemma 2.6.1. Put $y=n(g) w$ so that $y \in Z(\gamma)$. Since $n(g)$ carries $w$ to $y$ one clearly has

$$
\begin{equation*}
n(g) G^{w} n(g)^{-1}=G^{y} \tag{2.6.8}
\end{equation*}
$$

But now for any $a \in G^{w}$ one has $g^{-1} a \in G^{w}$ and hence

$$
\begin{equation*}
\psi_{g}: G^{w} \rightarrow G^{y} \tag{2.6.9}
\end{equation*}
$$

is an isomorphism of varieties (not of groups), where

$$
\begin{equation*}
\psi_{g}(a)=n(g) g^{-1} a n(g)^{-1} \tag{2.6.10}
\end{equation*}
$$

Lemma 2.6.2. Recalling that $G_{*}^{y}=G^{y} \cap \bar{N} H N$ and $G_{(*)}^{w}=G^{w} \cap s(\kappa) \bar{N} H N$ one has

$$
\begin{equation*}
\psi_{g}\left(G_{(*)}^{w}\right)=G_{*}^{y} \tag{2.6.11}
\end{equation*}
$$

Proof. Let $a \in G^{w}$ and put $c=g^{-1} a$ so that $s(\kappa)^{-1} a=s(\kappa)^{-1} g c=\bar{n} h n c$, where we have put $\bar{n}=\bar{n}(g), h=h(g)$, and $n=n(g)$. But $n c=n c n^{-1} n=$ $\psi_{g}(a) n$. Thus if we put $d=\psi_{g}(a)$ then $d \in G^{y}$ by (2.6.9) and

$$
\begin{equation*}
s(\kappa)^{-1} a-\bar{n} h d n \tag{2.6.12}
\end{equation*}
$$

But now by the decomposition (2.4.3) there exists a unique $\sigma \in W$ such that $d \in \bar{N} s(\sigma) H N$. Since $s(\sigma)$ normalizes $H$ this implies $s(\kappa)^{-1} a \in \bar{N} s(\sigma) H N$. Thus $d=\psi_{g}(a) \in G_{*}^{y}$ if and only if $s(\kappa)^{-1} a \in \bar{N} H N$ or if and only if $a \in G_{(*)}^{w}$. This proves the lemma.
Q.E.D.

Remark 2.6.2. Let the notation be as in the proof above. If $a \in G_{(*)}^{w}$ so that $d \in \bar{N} H N$, note that (2.6.12) implies $n(a)=n_{d} n$ using the notation of (2.4.6). That is, for any $g, a \in G_{(*)}^{w}$ one has

$$
\begin{equation*}
n(a)=n_{d} n(g) \tag{2.6.13}
\end{equation*}
$$

where $d=\psi_{g}(a)$.
The element $w=f+w_{o}$ in our considerations here can be chosen so that $\mathscr{F}(w)=\gamma$ where $\gamma$ is any given element in $\mathbb{C}^{l}$. That is, one has

Proposition 2.6. Given any $\gamma \in \mathbb{C}^{l}$ there exists a $w_{o} \in h$, unique up to conjugacy by $W$ such that $\mathscr{I}\left(w_{o}\right)=\gamma$. Furthermore if $w=f+w_{o}$ then one also has

$$
\begin{equation*}
\mathscr{I}(w)=\mathscr{I}\left(w_{o}\right) \tag{2.6.14}
\end{equation*}
$$

so that $\mathscr{I}(w)=\gamma$.
Proof. The first statement is a well-known result of Chevalley. It follows, for example, from the surjectivity part of the proof of Proposition 10 in [14, p. 355] together with Lemma 9.2 in [12]. The second statement is part of Lemma 11 in [14, p. 371].
Q.E.D.

Theorem 2.6. Let $y$ be a complex semi-simple Lie algebra and let $Z$ be the set of normalized Jacobi elementsing. This set is defined by (2.2.4). Let $\gamma \in \mathbb{C}^{l}$ and let $Z(\gamma) \subseteq Z$ be the isospectral leaf in $Z$ defined by $\gamma$. See (2.3.2). Let $w_{o}$ be any element in the Cartan subalgebra $\ell$ such that $\mathscr{I}\left(w_{o}\right)=\gamma$. Such an element $w_{o}$ exists. Let $w=f+w_{o}$, where we recall $f$ is defined by (1.5.4). Let $G_{(*)}^{w}$ be the Zariski open subset of the centralizer of $w$ defined by (2.6.4). Then $G_{(*)}^{w}$ is not empty so that it is a complex, connected l-dimensional, nonsingular rational variety.

Furthermore $Z(\gamma)$ is also such a variety. Moreover if $a \in G_{(*)}^{w}$ and $n(a) \in N$ is defined by (2.6.2) then $n(a) w \in Z(\gamma)$ and if $\beta_{(w)}(a)=n(a) w$ then

$$
\begin{equation*}
\beta_{(w)}: G_{(*)}^{w} \rightarrow Z(\gamma) \tag{2.6.15}
\end{equation*}
$$

is an isomorphism of algebraic varieties.
Proof. By Theorem 2.5, $Z(\gamma)$ is not empty and by Proposition 2.6 there exists an element $w_{o} \in \mathscr{h}$ such that $\mathscr{I}\left(w_{o}\right)=\gamma$.

Now let $y \in Z(\gamma)$. Then $y, w \in \mathscr{f}_{f}$ and $\mathscr{\mathscr { F }}(y)=\mathscr{I}(w)=\gamma$ by (2.6.14). But then by Proposition 2.3.2 there exists $n \in N$ such that

$$
\begin{equation*}
n w=y . \tag{2.6.16}
\end{equation*}
$$

But now if $y_{1} \in \tilde{d_{1}}$ is the maximal diagonal component of $y$ there exists $h \in H$ such that

$$
\begin{equation*}
e=h y_{1} \tag{2.6.17}
\end{equation*}
$$

where

$$
\begin{equation*}
e=\sum_{i=1}^{l} e_{\alpha_{i}} \tag{2.6.18}
\end{equation*}
$$

Thus $h y \in e+\ell$. On the other hand $s(\kappa)^{-1} w=\kappa^{-1} w_{o}+e$ since it is clear from (2.4.1) that

$$
\begin{equation*}
s(\kappa) e=f \tag{2.6.19}
\end{equation*}
$$

Hence $h y$ and $s(\kappa)^{-1} w$ are elements in $e+\mathscr{Z}$ and $\mathscr{I}(h y)=\mathscr{\mathscr { C }}\left(s(\kappa)^{-1} w\right)=\gamma$. Thus, as argued in the proof of Theorem 2.4, Proposition 2.3.2 applies with the roles of $\ell$ and $\mathscr{Z}, N$ and $\bar{N}$, and now $f$ and $e$ reversed. That is, there exists $\bar{n} \in \bar{N}$ such that $\bar{n} h y=s(\kappa)^{-1} w$ or $s(\kappa) \bar{n} h y=w$. But since $n w=y$ this implies that $g \in G^{w}$, where $g=s(\kappa) \bar{n} h n$. Hence $g \in G_{(*)}^{w}$ so that $G_{(*)}^{w}$ is not empty. Moreover note that $n=\boldsymbol{n}(g)$. But now by Lemma 2.6.1 the map $\beta_{(w)}$ exists and is a morphism of algebraic varieties since (2.6.5) is such a morphism. Furthermore by (2.6.16) one has

$$
\begin{equation*}
\beta_{(w)}(g)=n(g) y=w . \tag{2.6.20}
\end{equation*}
$$

Now recall (see (2.4.11)) the map $\beta_{y}: G_{*}^{y} \rightarrow Z(\gamma)$, where if $d \in G_{*}^{y}$ then $\beta_{y}(d)=$ $n_{d} y$, and the map $G_{(*)}^{w} \rightarrow G_{*}^{y}$, defined by restricting $\psi_{d}$ to $G_{(*)}^{w}$. See (2.6.11), where we recall $\psi_{g}(a)=n(g) g^{-1} a n(g)^{-1}$. We assert the diagram

is commutative. Indeed if $a \in G_{(*)}^{w}$ and $d=\psi_{g}(a)$ then $\beta_{v}\left(\psi_{g}(a)\right)=\beta_{y}(d)=$ $n_{d} y$. However, $y=n(g) w$ by (2.6.16). But $n_{d} n(g)=n(a)$ by (2.6.13). Thus $\beta_{y}\left(\psi_{g}(a)\right)=n(a) w=\beta_{(w)}(a)$. This proves that diagram (2.6.21) is commutative. But $\beta_{y}$ is an isomorphism by Theorem 2.4. Also since $\psi_{g}$ is an isomorphism the vertical map in (2.6.11) is an isomorphism by (2.6.11). Thus $\beta_{(w)}$ is an isomorphism of algebraic varieties. By Proposition 2.4, $G_{(*)}^{w}$ and hence $Z(\gamma)$ is isomorphic to a nonempty Zariski open subset of $\left(\mathbb{C}^{*}\right)^{p} \times \mathbb{C}^{q}$ for some $p, q$ where $p+q=l$ and hence has the properties described in the statement of the theorem.
Q.E.D.

## 3. The Parametrization $Z \cong H \times h_{+}$in the Real Case

3.1. Throughout Section 3 we assume that $F=\mathbb{R}$. We will see that the situation is more complicated but perhaps more interesting. In order to make use of the results of Sections $2.4-2.6$ we put $g_{\mathbb{C}}=g+i g$ and let $G_{\mathbb{C}}$ be the complex adjoint group. Of course we may assume that $G \subset G_{\mathbb{C}}$. Also let $n_{\mathbb{C}}, h_{\mathbb{C}}$, and $\bar{n}_{\mathbb{C}}$ be the complexifications of $n, h$, and $\bar{n}$ and let $N_{\mathbb{C}}, H_{\mathbb{C}}$, and $\bar{N}_{\mathbb{C}}$ be the corresponding subgroups of $G_{\mathbb{C}}$.

Now let $K \subseteq G$ be the subgroup corresponding to $k$ where we recall $g=$ $k+\not p$ is the given Cartan decomposition of $g$. See Section 1.2. Thus $K$ is a maximal compact subgroup. Put $P=\exp \not p$ so that as one knows, any $g \in G$ can be uniquely written

$$
\begin{equation*}
g=k p \tag{3.1.1}
\end{equation*}
$$

where $k \in K, p \in P$. One refers to (3.1.1) as the polar decomposition of $g$.
One also knows that the map

$$
\begin{equation*}
\not p \rightarrow P, \quad x \mapsto \exp x \tag{3.1.2}
\end{equation*}
$$

is a diffeomorphism. This clearly implies the familiar fact that the map

$$
\begin{equation*}
P \rightarrow P, \quad p \rightarrow p^{2} \tag{3.1.3}
\end{equation*}
$$

is also a diffeomorphism. For any $p \in P$ its unique square root in $P$ will be denoted by $p^{1 / 2}$.

Now the real form $g$ of $g_{\mathbb{C}}$ defines a conjugation operation on $g_{\mathbb{C}}$ and $G_{\mathbb{C}}$. That is, if $z \in \mathscr{g} \mathbb{C}$ and we write $z=x+i y$ for $x, y \in g$ then $z^{c} \in g_{\mathbb{C}}$ is defined by putting

$$
\begin{equation*}
z^{0}=x-i y \tag{3.1.4}
\end{equation*}
$$

Since $G_{\mathbb{C}}$ is the adjoint group the operation induces an automorphism $g \rightarrow g^{c}$ of $G_{\mathbb{C}}$, where of course for any $x \in g_{\mathbb{C}}$

$$
\begin{equation*}
(\exp x)^{c}=\exp x^{c} . \tag{3.1.5}
\end{equation*}
$$

Now the Cartan involution $\theta$ (see Section 1.2) extends by linearity to an automorphism of $g_{\mathbb{C}}$ and, since $G_{\mathbb{C}}$ is the adjoint group, it also defines an automorphism of $G_{\mathbb{C}}$ which we continue to denote by $\theta$. Now recall that the *-operation on $g$ was defined by putting $x^{*}=\theta(-x)$. We now extend the operation to $g_{\mathbb{C}}$ as a conjugate linear map by putting $x^{*}=\theta\left(-x^{c}\right)$ for any $x \in g_{\mathbb{C}}$. Thus if

$$
\begin{equation*}
g_{u}=k+i \not{ }_{h} \tag{3.1.6}
\end{equation*}
$$

then $g_{u}$ is a compact form of $g_{\mathbb{C}}$ and

$$
\begin{equation*}
x^{*}=-x \quad \text { if and only if } x \in g_{u} . \tag{3.1.7}
\end{equation*}
$$

One also defines a $*$-operation on $G_{\subset}$ by putting $g^{*}=\theta\left(\left(g^{c}\right)^{-1}\right)$ for any $g \in G_{\mathbb{C}}$. Clearly if $a, b \in G_{\mathbb{C}}, x, y \in \mathscr{g}_{\mathbb{C}}$ one has

$$
\begin{equation*}
(a b)^{*}=b^{*} a^{*} \quad \text { and } \quad[x, y]^{*}=\left[y^{*}, x^{*}\right] \tag{3.1.8}
\end{equation*}
$$

and also

$$
\begin{equation*}
(\exp x)^{*}=\exp x^{*} \tag{3.1.9}
\end{equation*}
$$

Now let $G_{u} \subseteq G_{\mathbb{C}}$ be the (maximal compact) subgroup of $G_{\mathbb{C}}$ corresponding to $g_{u}$ (see (3.1.6)). Also let $P_{u}=\exp i g_{u}$ so that as one knows $P_{u}$ is a closed submanifold of $G_{\mathrm{C}}$ and the map

$$
\begin{equation*}
G_{u} \times P_{u} \rightarrow G_{\mathbb{C}}, \quad(w, q) \mapsto w q \tag{3.1.10}
\end{equation*}
$$

is a diffeomorphism. In fact the decomposition of $G_{\mathbb{C}}$ defined by (3.1.10) is just a polar decomposition. One has $K \subseteq G_{u}, P \subseteq P_{u}$ and restricting (3.1.10) one has a diffeomorphism

$$
\begin{equation*}
K \times P \rightarrow G \tag{3.1.11}
\end{equation*}
$$

which of course defines our given polar decomposition of $G$.
Now it is clear that $G$ is stable under the $*$-operation and in fact if $g \in G$ is given by (3.1.1) then by (3.1.8) one sees that

$$
\begin{equation*}
g^{*}=p k^{-1} . \tag{3.1.12}
\end{equation*}
$$

This also uses (3.1.7) and (3.1.9), which imply $a=a^{*}$ if $a \in P_{u}$ and $a^{*}=a^{-1}$ if $a \in G_{u}$.

Now note that

$$
\begin{equation*}
h \subseteq p h \tag{3.1.13}
\end{equation*}
$$

Indccd for any $\varphi \in \Delta$ let

$$
\begin{equation*}
h_{\varphi}=\left[e_{\varphi}, e_{-\varphi}\right] \in h . \tag{3.1.14}
\end{equation*}
$$

But then since $e_{-\infty}=e_{\Phi}^{*}(\operatorname{see}(1.5 .1))$ one has $h_{\infty}^{*}=h_{\varphi}$ by (3.1.8) so that $h_{\infty} \in \not p$. On the other hand the root normals $h_{\varphi}$ span $h$ so that one has (3.1.13). Since $\operatorname{dim} h=l$ the Cartan subalgebra $h$ is of course a Cartan subspace of $\not p$ and by exponentiating (3.1.13) one has

$$
\begin{equation*}
H \subseteq P . \tag{3.1.15}
\end{equation*}
$$

Now let $T$ be the subgroup of $G_{\mathbb{C}}$ corresponding to $i \hbar$. Then by (3.1.6) and (3.1.13) one has $T \subseteq G_{u}$ (in fact $T$ is clearly a maximal torus in $G_{u}$ ) and clearly the restriction of (3.1.10) to $T \times H$ defines a diffeomorphism

$$
\begin{equation*}
T \times H \rightarrow H_{\mathbb{C}} \tag{3.1.16}
\end{equation*}
$$

One notes then that if $a \in H_{\mathbb{C}}$ and we write $a=t h$, where $t \in T$ and $h \in H$ then

$$
\begin{equation*}
a^{*}=h t^{-1} . \tag{3.1.17}
\end{equation*}
$$

Now let $\bar{M}$ be the set of all elements of order 2 or 1 in $H_{\mathbb{C}}$. Obviously

$$
\begin{equation*}
\tilde{M} \subseteq T \tag{3.1.18}
\end{equation*}
$$

by (3.1.16) and since $T$ is a torus of rank $l$ it is clear that $\tilde{M}$ is a finite Abelian group and in fact its cardinality $|\tilde{M}|$ is given by

$$
\begin{equation*}
|\tilde{M}|=2^{l} \tag{3.1.19}
\end{equation*}
$$

Lemma 3.1.1. Let $a \in H_{\mathbb{C}}$. Then $a=a^{*}$ if and only if $a \in \tilde{M} H$.
Proof. This is immediate from (3.1.16) and (3.1.17) since $H_{\mathbb{C}}$ is commutative.
Q.E.D.

Let $M=\tilde{M} \cap G$.

Lemma 3.1.2. Let $g \in G$ and assume $g=g^{*}$. Then $g \in K$ if and only if $g^{2}=1$. Furthermore one has

$$
\begin{equation*}
M=K \cap T \tag{3.1.20}
\end{equation*}
$$

Proof. If $g=k p$ is the polar decomposition of $g$ then by (3.1.12)

$$
\begin{equation*}
k p=p k^{-1} \tag{3.1.21}
\end{equation*}
$$

Thus if $g \in K$ one has $g=k=k^{-1}$ so that $g^{2}=1$. On the other hand by (3.1.21) one has $g^{2}=p^{2} \in P$. Thus if $g^{2}=1$ then $p^{2}=1$ which implies $p=1$. See (3.1.3). Hence $g \in K$. This proves the first statement. Now $M \subseteq K \cap T$ by the first statement and Lemma 3.1.1. On the other hand if $g \in K \subseteq G$ then $g^{c}=g$ recalling (3.1.4) and (3.1.5). But if $g \in T=\exp$ ih then also $g^{c}=g^{-1}$. Thus $K \cap T \subseteq \tilde{M}$. However, $K \cap T \subseteq G$ also, so that one has (3.1.20). Q.E.D.

For any $x \in g$ let $G_{e}^{x}$ be the identity component of the centralizer $G^{x}$ of $x$ in $G$.
Proposition 3.1. Assume $x \in \notin$ is regular. Then

$$
\begin{equation*}
g=g^{*} \quad \text { for any } g \in G^{x} \tag{3.1.22}
\end{equation*}
$$

On the other hand one has

$$
\begin{equation*}
G^{x} \cap P=G_{e}^{x} \tag{3.1.23}
\end{equation*}
$$

Furthermore as Lie groups one has the isomorphism

$$
\begin{equation*}
G_{e}^{x} \cong \mathbb{R}^{\boldsymbol{l}} \tag{3.1.24}
\end{equation*}
$$

Proof. From the conjugation theory of Cartan subspaces of $\not p$ one knows that $x$ is $K$-conjugate to an element in $\ell$. But since the $*$-operation commutes with the conjugacy action of $K$ and $I$ is stable under this action it then suffices to assume, as we shall, that $x \in h$. But now clearly $G^{x} \subseteq G_{\mathbb{C}}{ }^{x}$, where $G_{\mathbb{C}}{ }^{x}$ is the centralizer of $x$ in $G_{\mathbb{C}}$. On the other hand $G_{\mathbb{C}}{ }^{x}=H_{\mathbb{C}}$ since $G_{\mathbb{C}}{ }^{x}$ is connected by Proposition 2.4 and $\mathscr{g}_{\mathbb{C}}{ }^{x}=h_{\mathbb{C}}$ by regularity. Thus

$$
\begin{equation*}
G^{x}=H_{\mathbb{C}} \cap G \tag{3.1.25}
\end{equation*}
$$

But now comparing the diffeomorphisms (3.1.11) and (3.1.16) one has $G^{x}=$ ( $K \cap T$ ) $H$ by (3.1.25). But $K \cap T=M$ by (3.1.20) so that $G^{x}=M H$. But then in fact restricting (3.1.16) one has a diffeomorphism

$$
\begin{equation*}
M \times H \rightarrow G^{x} \tag{3.1.26}
\end{equation*}
$$

But then by the uniqueness of the polar decomposition this implies

$$
\begin{equation*}
G^{x} \cap P=H \tag{3.1.27}
\end{equation*}
$$

However, one also clearly has $G_{e}{ }^{x}=H$ by (3.1.26). This proves (3.1.23). Furthermore one also has (3.1.22) by Lemma 3.1.1 and (3.1.26).

Now the restriction of the bijection (3.1.2) to $\hbar$ is clearly an isomorphism $h \rightarrow H$ of Abelian Lie groups. This proves (3.1.24).
Q.E.D.

### 3.2. Now put

$$
\begin{equation*}
G_{*}=\bar{N} H N \tag{3.2.1}
\end{equation*}
$$

Proposition 3.2. Assume $x \in \neq$ is regular. Then

$$
\begin{equation*}
G^{x} \cap G_{*}=G_{e}^{x} \tag{3.2.2}
\end{equation*}
$$

where we recall that $G_{e}{ }^{x}$ is the identity component of $G^{x}$.
Proof. Assume that $g$ is in the left side of (3.2.2). Write $g=\bar{n} h n$ for $\bar{n} \in \bar{N}$, $h \in H$ and $n \in N$. Then $g^{*}=n^{*} h \bar{n}^{*}$. But now by (1.5.1) and (3.1.9) one has

$$
\begin{equation*}
N^{*}=\bar{N} \quad \text { and } \quad(\bar{N})^{*}=N \tag{3.2.3}
\end{equation*}
$$

Thus $n^{*} \in \bar{N}$ and $(\bar{n})^{*} \in N$. However, $g=g^{*}$ by (3.1.22). But then by the isomorphism (2.4.5) applied to $N_{\mathrm{C}}, H_{\mathrm{C}}$, and $\bar{N}_{\mathrm{C}}$ one has $n^{*}=\bar{n}$. Hence $g=n^{*} h n$. Thus if we put $h^{1 / 2} n=a \in G$ one has $g=a^{*} a$. However, for any $b \in G$ one has

$$
\begin{equation*}
b^{*} b \in P \tag{3.2.4}
\end{equation*}
$$

Indeed if $b=k p$ is the polar decomposition of $b$ then $b^{*} b=p^{2} \in P$, proving (3.2.4). Thus $g \in P$ and hence $g \in G^{x} \cap P$. But then $g \in G_{e}{ }^{x}$ by (3.1.23).

Now the Iwasawa decomposition asserts that the map

$$
\begin{equation*}
K \times H \times N \rightarrow G, \quad(k, a, n) \rightarrow k a n \tag{3.2.5}
\end{equation*}
$$

is a diffeomorphism. Thus if $g \in G$ we can uniquely write $g=k a n$, where $k \in K, a \in H$, and $n \in N$, and we refer to this as the Iwasawa decomposition of $g$.

Now recalling that $G_{*}=\bar{N} H N$ we assert that

$$
\begin{equation*}
P \subseteq G_{*} \tag{3.2.6}
\end{equation*}
$$

Indeed let $p \in P$ and let $p^{1 / 2}=k a n$ be the Iwasawa decomposition of the square root $p^{1 / 2}$ of $p$. Then $\left(p^{1 / 2}\right)^{*}=n^{*} a k^{-1}$. Since $\left(p^{1 / 2}\right)^{*}=p^{1 / 2}$ this implies

$$
\begin{equation*}
p=\left(p^{1 / 2}\right)^{*}\left(p^{1 / 2}\right)=n^{*} a^{2} n \tag{3.2.7}
\end{equation*}
$$

But $n^{*} \in \bar{N}$ and $a^{2}=h \in H$. This proves (3.2.6).
Now to prove the proposition we have only to show that the right side of (3.2.2) is contained in the left side. Let $g \in G_{e}^{x}$. Then $g \in P$ by (3.1.23). Thus $g \in G_{*}$ by (3.2.6). Hence $g$ is contained in the left side of (3.2.2). Q.E.D.

Now recalling (2.3.1) let $\gamma \in \mathscr{I}(Z)$. We wish to obtain results about the closed subset $Z(\gamma)$ in the present case (i.e., where $F=\mathbb{R}$ ) which are analogous to those given in Theorem 2.6. If $y \in Z$ then even though $y$ is not necessarily in $\nless$ it is $H$-conjugate to an element in $\not p$ so, as we now observe, the statement of Proposition 3.2 holds for $y$.

Lemma 3.2. Let $y \in Z$. Then there exists a unique $c \in H$ such that

$$
\begin{equation*}
c y \in \not x \tag{3.2.8}
\end{equation*}
$$

Also, one has

$$
\begin{equation*}
G^{y} \cap G_{*}=G_{e}{ }^{y} \tag{3.2.9}
\end{equation*}
$$

Proof. Let $x \in h$ and let $a_{i} \in \mathbb{R}^{*}$ (i.e., $a_{i}>0$ ), $i=1, \ldots, l$, be such that - $y=f+x+\sum_{i} a_{i} e_{\alpha_{i}}$. Now for any $c \in H$ one has

$$
\begin{equation*}
c y=x+\sum_{i}\left(c^{-\alpha_{i}} e_{-\alpha_{i}}+c^{\alpha_{i}} a_{i} e_{\alpha_{i}}\right) \tag{3.2.10}
\end{equation*}
$$

But now recalling the isomorphism (2.2.8) it follows that there is a unique $c \in H$ such that $c^{-\alpha_{i}}=c^{\alpha_{i}} a_{i}$, namely, that element $c$ such that $c^{\alpha_{i}}=a_{i}^{-1 / 2}$. However, $x^{*}=x$ by (3.1.13) and $e_{\alpha_{i}}^{*}=e_{-\alpha_{i}}$ by (1.5.1). Thus only for this value of $c$ does one have $(c y)^{*}=c y$ or equivalently $c y \in \mu$. This proves the first statement of the lemma.

Now fix $c \in H$ so that one has (3.2.8). But $y \in R$ by (2.1.7). Thus $c y$ is a regular element of $p$ and hence by (3.2.2) one has $G^{c y} \cap G_{*}=G_{e}^{c y}$. On the other hand one certainly has $c^{-1} G^{c v} c=G^{y}$ and a similar relation for the identity components. Since $\bar{N}, H$, and $N$ are stable under conjugation by $c$ it follows that $G_{*}$ is stable under conjugation by $c$. This proves (3.2.9).
Q.E.D.

Now we may apply the results of Sections 2.4-2.6 to the complexification $G_{\mathbb{C}}$ of $G$. Thus recalling (2.4.5) the map $\bar{N}_{\mathbb{C}} \times H_{\mathbb{C}} \times N_{\mathbb{C}} \rightarrow \bar{N}_{\mathbb{C}} H_{\mathbb{C}} N_{\mathbb{C}}$ defined by multiplication is an algebraic isomorphism of nonsingular affine varieties.

In particular then it is a diffeomorphism. Thus, by dimension, $G_{*}$ is an open connected subset of $G$ and, by restriction the map

$$
\begin{equation*}
\bar{N} \times H \times N \rightarrow G_{*}, \quad(\bar{n}, h, n) \rightarrow \bar{n} h n \tag{3.2.11}
\end{equation*}
$$

is a diffeomorphism. Thus using the notation of (2.4.6) for the complexification $G_{C}$ one has the decomposition

$$
\begin{equation*}
d=\bar{n}_{d} h_{d} n_{d} \tag{3.2.12}
\end{equation*}
$$

for any $d \in G_{*}$, where now $\bar{n}_{d} \in \bar{N}, h_{d} \in H$, and $n_{d} \in N$. Furthermore upon restriction of (2.4.7) the map

$$
\begin{equation*}
G_{*} \rightarrow N, \quad d \rightarrow n_{d} \tag{3.2.13}
\end{equation*}
$$

is smooth.

Theorem 3.2. Let $g$ be a real split semi-simple Lie algebra. Let $l=$ rank $g$ and let $G$ be the adjoint group of $g$. Let $Z \subseteq g$ be the $2 l$ dimensional manifold of normalized Jacobi elements in $g$. The manifold $Z$ is defined by (2.2.3) and (2.2.4). Now let $\mathscr{I}: g \rightarrow \mathbb{R}^{l}$ be the map (see (2.3.1)) defined by the fundamental $G$-invariant polynomials on $g$ and for any $\gamma \in \mathscr{I}(Z)$ let $Z(\gamma)=\mathscr{J}^{-1}(\gamma) \subseteq Z$. Then the isospectral set $Z(\gamma)$ is a closed connected submanifold of dimension $l$ in $Z$. In fact one has a diffeomorphism

$$
\begin{equation*}
Z(\gamma) \cong \mathbb{R}^{l} \tag{3.2.14}
\end{equation*}
$$

Furthermore for any $y \in Z(\gamma)$ the identity component $G_{e}{ }^{y}$ of the centralizer of $y$ in $G$ is also isomorphic (as Lie groups) to $\mathbb{R}^{l}$. Moreover using the notation of (3.2.1) one has $G_{e}^{y} \subseteq G_{*}$ and if for any $d \in G_{e}^{y}$ the element $n_{d} \in N$ is defined by (3.2.12) then $n_{d} y \in Z(\gamma)$ and the map

$$
\begin{equation*}
\beta_{y}: G_{e}^{y} \rightarrow Z(\gamma), \quad d \mapsto n_{d} y \tag{3.2.15}
\end{equation*}
$$

is a diffeomorphism.
Proof. Let $\left(d_{1}\right)_{\mathbb{C}}$ be the complexification of the diagonal $d_{1}$ so that $\left(d_{1}\right)_{\mathbb{C}}$ is stable under the $H_{\mathbb{C}}$ and let $\left(\tilde{d_{1}}\right)_{\mathbb{C}}$ be defined by (2.2.2), where $F=\mathbb{C}$. Put $Z_{\mathbb{C}}=f+h_{\mathbb{C}}+\left(\tilde{d}_{1}\right)_{\mathbb{C}}$ so that we can apply Theorem 2.4 , where $Z_{\mathbb{C}}$ replaces $Z$. But then regarding $S\left(g_{\mathrm{C}}\right)^{G_{\mathrm{C}}}$ as the complexification of $S(g)^{G}$ one has $y \in Z_{\mathbb{C}}(\gamma)$. Furthermore if we let $\left(\beta_{\mathbb{C}}\right)_{y}$ now denote the $\beta_{y}$ of Theorem 2.4 then by Theorem 2.4 one has an isomorphism

$$
\begin{equation*}
\left(\beta_{\mathbb{C}}\right)_{y}:\left(G_{\mathbb{C}}\right)_{*} \rightarrow Z_{\mathbb{C}}(\gamma) \tag{3.2.16}
\end{equation*}
$$

But now by (3.2.9) one has $G_{e}{ }^{y} \subseteq\left(G_{\mathbb{C}}{ }^{y}\right)_{*}$ and clearly for $d \in G_{e}{ }^{y}$ the definition of $h_{d}$ and $n_{d}$ whether given by (3.2.12) or (2.4.0) is the same. But now $Z(\gamma) \subseteq$ $Z_{\mathbb{C}}(\gamma)$ and if $d \in G_{e}{ }^{y}$ we assert that $\left(\beta_{\mathbb{C}}\right)_{y}(d) \in Z(\gamma)$. Indeed $\left(\beta_{\mathbb{C}}\right)_{y}(d)=n_{d} y$ so that $\left(\beta_{\mathrm{C}}\right)_{y}(d) \in g$. But then recalling (2.4.13) one must have $x \in h$. Furthermore again using the notation of (2.4.13) if we write $y_{1}=\sum a_{i} e_{\alpha_{i}}$ one has $a_{i}>0$ since $y \in Z$. But then

$$
\begin{equation*}
h_{d}^{-1} y_{1}=\sum\left(h_{d}^{-\alpha_{i}} a_{i}\right) e_{\alpha_{i}} \tag{3.2.17}
\end{equation*}
$$

Since $h_{d}^{-\alpha_{i}}>0$ this implies by (2.4.13) that $\left(\beta_{\mathbb{c}}\right)_{y}(d) \in Z(\gamma)$ and hence one has a map (3.2.15) and the map (3.2.15) is just the restriction of (2.4.11) to $G_{e}{ }^{y}$. We next assert that (3.2.15) is surjective. Let $w \in Z(\gamma)$. Recalling the argument of the surjectivity of $\beta_{y}$ in the proof of Theorem 2.4 we first observe that the element $n \in N_{\mathrm{C}}$ of (2.4.14) is in fact in $N$. This is clear from Proposition 2.3.2 since $y$ and $w$ are in $b_{f}$ and not just $f+b_{c}$. Similarly the unique element $h \in H_{\mathrm{C}}$ such that $h w_{1}=y_{1}$ is in fact in $H$ since $w_{1}, y_{1} \in \tilde{d_{1}}$. Finally the element $\bar{n} \in \bar{N}_{C}$ is in $\bar{N}$ again by Proposition 2.3.2 and (2.4.15). Thus $d=\bar{n} h n$ is in $G$ and we recall $\left(\beta_{\mathrm{C}}\right)_{y}(d)-w$. But clearly $d \in G^{y} \cap G_{*}$. Thus $d \in G_{e}^{y}$ by (3.2.9) and hence (3.2.15) is surjective.

Now $Z(\gamma)$, by Proposition 2.3.1, is a (not necessarily connected) submanifold of dimension $l$ in $Z$. But in fact (2.4.11) is an isomorphism of nonsingular varieties and hence as noted in Remark 2.4.2 is a diffeomorphism of manifolds. Since (3.2.15) is surjective and is the restriction of (2.4.11) to $G_{e}{ }^{y}$ it follows that (3.2.15) is a diffeomorphism. Now, using the notation of Lemma 3.2, $G_{e}{ }^{y}$ is isomorphic to $G_{e}^{c y}$. But $G_{e}^{c y}$ is isomorphic to $\mathbb{R}^{l}$ as Lie groups by (3.1.24). Thus $G_{e}{ }^{y}$ and hence $Z(\gamma)$ are isomorphic to $\mathbb{R}^{l}$ as manifolds. This proves the theorem.
Q.E.D.

Remark 3.2. One aspect of Theorem 3.2 which we think ought to be emphasized is the connectivity of $Z(\gamma)$. If $Z$ was defined by (2.2.3) and (2.2.4) except that, say, $a_{i}<0$ instead of $a_{i}>0$, then it would not be necessarily true that $Z(\gamma)$ is connected. An example of this disconnectivity is easily constructed for the case where $g$ is the Lie algebra of $S l(2, \mathbb{R})$.
3.3. Now let $h_{+}$be the open Weyl chamber in $h$ defined by putting

$$
\begin{equation*}
\ell_{+}=\left\{x \in h \mid\langle\varphi, x\rangle>0 \text { for all } \varphi \in \Delta_{+}\right\} . \tag{3.3.1}
\end{equation*}
$$

Also let $\bar{h}_{+}$be the closure of $h_{+}$in $h$ so that $\bar{h}_{+}$is a closed Weyl chamber. As one knows and easily sees, $h \subseteq g^{x}$ for any $x \in \bar{h}_{+}$and equality holds if and only if $x \in h_{+}$. That is,

$$
\begin{equation*}
h_{+}=\bar{h}_{+} \cap R \tag{3.3.2}
\end{equation*}
$$

Now recalling the map (2.3.1) one has

Proposition 3.3.1. $\mathscr{I}\left(\ell_{+}\right)$is a connected open subset of $\mathbb{R}^{l}$ and the map

$$
\begin{equation*}
h_{+} \rightarrow \mathscr{I}\left(h_{+}\right), \quad x \rightarrow \mathscr{I}(x) \tag{3.3.3}
\end{equation*}
$$

is a diffeomorphism.
Proof. The bilinear form $Q$ is positive definite on $\hbar$ and hence we can find a basis $x_{i}$ of $g$ such that $x_{i}$ for $i \leqslant l$ is an orthonormal basis of $h$ and $Q\left(x_{i}, x_{j}\right)=0$ for $i \leqslant l<j$. But now if $x \in h_{+}$then $g^{x}=h$. But then for any $I \in S(g)^{G}$ one has $(\delta I)(x) \in h$ by (1.3.1). But then using the basis $x_{i}$ of $g$ and its dual basis the summands in (1.2.2) must vanish for $i>l$. That is, one has $(\delta I)(x)=\sum_{j=1}^{l}\left(i\left(x_{j}\right) I\right)(x) x_{j}$. On the other hand as already noted in the proof of Proposition 2.1, since $x \in R$, the elements $\left(\delta I_{k}\right)(x), k=1, \ldots, l$, are linearly independent (see Theorem 9 in [14, p. 382]). Thus the $l \times l$ matrix $M_{j k}=$ $\left(i\left(x_{j}\right) I_{k}\right)(x)$ is nonsingular. But this matrix is just the Jacobian of the map (3.3.3) at $x$ with respect to the coordinate system on $h_{+}$defined by the $x_{j}$, $1 \leqslant j \leqslant l$. Thus $\mathscr{I}\left(h_{+}\right)$is open and (3.3.3) is a local diffeomorphism. It is connected since $h_{+}$is clearly connected. It suffices then only to show that (3.3.3) is injective. However, if $x, y \in h_{+}$and $\mathscr{I}(x)=\mathscr{I}(y)$ then by Chevalley's theorem $x$ and $y$ are $W$-conjugate. But, as one knows, $\tilde{h}_{+}$is a fundamental domain for the action of $W$ on $h$. Thus $x=y$.
Q.E.D.

An element $x \in \mathscr{g}$ is called semi-simple if ad $x$ is diagonalizable over $\mathbb{C}$. It is called real semi-simple if ad $x$ is diagonalizable over $\mathbb{R}$. The following is well known.

Lemma 3.3.1. Let $x \in g$. Then $x$ is real semi-simple if and only if $x$ is $G$ conjugate to an element in $h$.

Proof. See, e.g., Proposition 2.4 in [16] and Theorem 2(2) in [23, p. 383].

> Q.E.D.

Now let $R_{+}$denote the set of all regular real semi-simple elements in $g$. It is clear that $R_{+}$is stable under the action of $G$.

Lemma 3.3.2. Let $x \in g$. The following conditions are all equivalent.
(1) $x \in R_{+}$.
(2) $x$ is $G$-conjugate to an element in $h_{+}$.
(3) $\mathscr{I}(x) \in \mathscr{I}\left(h_{+}\right)$.

Proof. Using the fact that $\bar{h}_{+}$is a fundamental domain for the action of the Weyl group $W$ on $h$ the equivalence of (1) and (2) follows from Lemma 3.3.1 and (3.3.2). Obviously (2) implies (3). It suffices to show that (3) implies (1).

Assume $\mathscr{I}(x) \in \mathscr{I}\left(h_{+}\right)$. Let $z \in \ell_{+}$be such that $\mathscr{F}(x)=\mathscr{I}(z)$. Then since $z \in R$ one has $x \in \overline{G_{\mathbb{C}} z}$ by Theorem 3 in [14, p. 365] (see in particular (3.8.7) in [14]), where the bar denotes closure. However, since $z$ is semi-simple $G_{\mathbb{C}} z$ is closed by Lemma 5 in [14, p. 353]. Thus $x \in G_{\mathbb{C}} z$. But then ad $x$ is semisimple with real eigenvalues. But then $x$ is real semi-simple. Furthermore $x \in R$ since $z \in R$. Thus $x \in R_{+}$.
Q.E.D.

Proposition 3.3.2. One has

$$
\begin{equation*}
Z \subseteq R_{+} \tag{3.3.4}
\end{equation*}
$$

Furthermore

$$
\begin{equation*}
\mathscr{I}(Z) \subseteq \mathscr{I}\left(f_{+}\right) . \tag{3.3.5}
\end{equation*}
$$

Proof. The second statement follows from (3.3.4) and Lemma 3.3.2. Thus it suffices to prove (3.3.4). Let $y \in Z$ then $y \in R$ by (2.1.7). But now $y$ is conjugate to an element in $p$ by Lemma 3.2. However, the elements of $p$ are real semisimple. In fact any element in $\nsim$ is $K$-conjugate to an element in $\hbar$ by the conjugacy theorem of Cartan subspaces. Thus $y$ is real semi-simple by Lemma 3.1. Hence $y \in R_{+}$.
Q.E.D.

Remark 3.3. In the complex case we found (see (2.5.6)) that $\mathscr{F}(Z)=\mathbb{C}^{l}$. From the nonconjugacy of Cartan subalgebras in $g$ in the present case it is clear from Proposition 3.3 that we cannot now expect an analogous result.
3.4. Recall that (see (3.1.1)) $z \rightarrow z^{c}$ is the conjugate linear automorphism of $g_{\mathbb{C}}$ given by putting $z^{c}=x-i y$, where we have written $z=$ $x+i y$ for $x, y \in g$. We recall also that this automorphism induces an automorphism

$$
\begin{equation*}
G_{\mathbb{C}} \rightarrow G_{\mathbb{C}}, \quad a \rightarrow a^{c} \tag{3.4.1}
\end{equation*}
$$

of $G_{\mathbb{C}}$. One of course has

$$
\begin{equation*}
a^{c} z^{c}=(a z)^{c} \tag{3.4.2}
\end{equation*}
$$

for $a \in G_{\mathbb{C}}, z \in g_{\mathbb{C}}$. One also has

$$
\begin{equation*}
(\exp z)^{c}=\exp \left(z^{c}\right) \tag{3.4.3}
\end{equation*}
$$

Now let

$$
\begin{equation*}
\tilde{G}=\left\{a \in G_{\mathbb{C}} \mid a^{c}=a\right\} \tag{3.4.4}
\end{equation*}
$$

One easily has

Proposition 3.4. $\tilde{G}$ is the set of all elements $g \in G_{\mathbb{C}}$ which stabilize g. In particular one has

$$
\begin{equation*}
G \subseteq G^{\star} . \tag{3.4.5}
\end{equation*}
$$

Proof. Clearly an element in $G_{\mathbb{C}}$ is determined (since it is $\mathbb{C}$-linear) by its restriction to $g$. The first statement and hence the second of the proposition then follows from (3.4.2) since, in the notation of (3.4.2), $z^{c}=z$ if and only if $z \in g$.
Q.E.D.

Recall (see (3.1.18)) that $\tilde{M}$ is the set of all $a \in H_{\mathrm{C}}$ such that $a^{2}=1$.

Lemma 3.4.1. One has

$$
\begin{equation*}
\tilde{G} \cap H_{C}=\tilde{M} H \tag{3.4.6}
\end{equation*}
$$

Proof. Now recalling (3.1.16) and the notation of (3.1.16) any element $a \in H_{\mathbb{C}}$ can be uniquely written $a=s h$ when $s \in T$ and $h \in H$. But then one clearly has $a^{c}=s^{-1} h$. Thus $a=a^{c}$ if and only if $s=s^{-1}$, that is, if and only if $s \in \tilde{M}$.
Q.E.D.

For any $x \in \mathcal{g}$ let $\vec{G}^{x}$ denote the centralizer of $x$ in $\vec{G}$.

Lemma 3.4.2. For any $x \in h_{+}$one has

$$
\begin{equation*}
\tilde{G}^{x}=\tilde{M} H \tag{3.4.7}
\end{equation*}
$$

Proof. One has $\tilde{G}^{x}=\tilde{G} \cap G_{\mathbb{C}}{ }^{x}$. But since $x$ is a regular element of $g_{\mathbb{C}}$ and $x \in h_{+}$then

$$
\begin{equation*}
G_{\mathbb{C}}{ }^{x}=H_{\mathbb{C}} \tag{3.4.8}
\end{equation*}
$$

by Proposition 2.4. The result then follows from (3.4.6).
Remark 3.4. We note then by (3.1.16) and (3.4.7) that $\mathcal{G}^{x}$ has $2^{l}$ connected components for any $x \in \hbar_{+}$. Recalling Lemma 3.3.2 the same statement is then true for any $x \in R_{+}$. Furthermore since $H=G_{e}^{x}=G_{e}^{x}$ for $x \in h_{+}$ one has

$$
\begin{equation*}
G_{e}^{x}=G_{e}^{x} \quad \text { for any } \quad x \in R_{+} . \tag{3.4.9}
\end{equation*}
$$

Lemma 3.4.3. One has

$$
\begin{equation*}
\tilde{G} \cap \bar{N}_{\mathbb{C}} H_{\mathbb{C}} N_{\mathbb{C}}=\bar{N} \tilde{M} H N \tag{3.4.10}
\end{equation*}
$$

and any element $g$ in this set can be uniquely written

$$
\begin{equation*}
g=\bar{n} m h n, \tag{3.4.11}
\end{equation*}
$$

where $\bar{n} \in \bar{N}, m \in \tilde{M}, h \in H$, and $n \in N$.
Proof. Clearly the right side of (3.4.10) is contained in the left side. Let $g$ be in the left side so that we can uniquely write $g=\bar{n} a n$, where $\bar{n} \in \bar{N}_{\mathbb{C}}$, $a \in H_{\mathbb{C}}$, and $n \in N_{\mathbb{C}}$. But $g^{c}=(\bar{n})^{c} a^{c} n^{c}$. However, $\bar{N}_{\mathbb{C}}, H_{\mathbb{C}}$, and $N_{\mathbb{C}}$ are clearly stable under (3.4.1) by (3.4.3). Thus $\bar{n}^{c}=\bar{n}, a^{e}=a$, and $n^{c}=n$ by uniqueness. But the map exp: $n_{\mathbb{C}} \rightarrow N_{\mathbb{C}}$ is a bijection. Thus $n \in N$ by (3.4.3). Similarly $\bar{n} \in \bar{N}$. One also has $a \in \tilde{M} H$ by (3.4.6). Thus $g$ is in the right side of (3.4.10). The uniqueness of (3.4.11) follows from the injectivity of (2.4.5) and (3.1.16).
Q.E.D.

Now for any $\sigma \in W$ let $s(\sigma) \in G_{\mathbb{C}}$ be the unique element in the normalizer of $\ell$ in $G_{\mathbb{C}}$ such that

$$
\begin{equation*}
s(\sigma) e_{\alpha_{i}}=e_{\sigma \alpha_{i}}, \quad i=1, \ldots, l . \tag{3.4.12}
\end{equation*}
$$

Lemma 3.4.4. One has $s(\sigma) \in \tilde{G}$ for any $\sigma \in W$.
Proof. By Proposition 3.4 it suffices to show that $s(\sigma)$ stabilizes $g$. Since $g$ is generated by $e_{\alpha_{i}}, e_{-\alpha_{i}}, i=1, \ldots, l$, it suffices to show that $s(\sigma) e_{-\alpha_{i}} \in g$ for all $i$. Thus if we let $\lambda_{i} \in \mathbb{C}$ be defined by $s(\sigma) e_{-\alpha_{i}}=\lambda_{i} e_{-\sigma \alpha_{i}}$ it suffices to show $\lambda_{i} \in \mathbb{R}$. But now $Q\left(e_{\Phi}, e_{-\varphi}\right)$ is a nonvanishing real number for any $\varphi \in \Delta$. But if we extend $Q$ to $\mathfrak{g c}$ by $\mathbb{C}$-linearity one has $Q\left(s(\sigma) e_{-\alpha_{i}}, e_{\sigma \alpha_{i}}\right)=\lambda_{i} Q\left(e_{\sigma \alpha_{i}}, e_{-\sigma x_{i}}\right)$ so that it suffices to show $Q\left(s(\sigma) e_{-\alpha_{i}}, e_{\sigma \alpha_{i}}\right) \in \mathbb{R}$. But by the invariance of $Q$ one has $Q\left(s(\sigma) e_{-\alpha_{i}}, e_{\sigma \alpha_{i}}\right)=Q\left(e_{-\alpha_{i}}, s(\sigma)^{-1} e_{\sigma \alpha_{i}}\right)=Q\left(e_{-\alpha_{i}}, e_{\alpha_{i}}\right) \in \mathbb{R}$.
Q.E.D.

Now it is clear from Proposition 3.4 that $G$ is the identity component of $G$. On the other hand $G_{*}=\bar{N} H N$ is an open connected subset of $\tilde{G}$. But $s(\kappa) \in \tilde{G}$ by Lemma 3.4.4. Thus if we put

$$
\begin{equation*}
G_{(*)}=s(\kappa) G_{*} \tag{3.4.13}
\end{equation*}
$$

then $G_{(*)}$ is an open connected subset of $\tilde{G}$. Furthermore using the notation of (2.6.2) for the complexification $G_{\mathbb{C}}$ of $G$ one has the decomposition

$$
\begin{equation*}
s(\kappa)^{-1} g=\bar{n}(g) h(g) n(g) \tag{3.4.14}
\end{equation*}
$$

for any $g \in G_{(*)}$, where now $\bar{n}(g) \in \bar{N}, h(g) \in H$, and $n(g) \in N$. Furthermore by restriction the map

$$
\begin{equation*}
G_{(*)} \rightarrow N, \quad g \rightarrow n(g) \tag{3.4.15}
\end{equation*}
$$

is smooth.
3.5. The following lemma is no doubt known. Because of the overriding importance for us of the element $\bar{n}_{f}(w)$ in Lemma 3.5.2 we prove Lemma 3.5.1 for completeness.

Lemma 3.5.1. For any $\bar{n} \in \bar{N}$ and $x \in h_{+}$one has $\bar{n} x \in x+\bar{n}$ so that one has a map

$$
\begin{equation*}
\bar{N} \times h_{+} \rightarrow h_{+}+\bar{n}, \quad(\bar{n}, x) \rightarrow \bar{n} x \tag{3.5.1}
\end{equation*}
$$

Furthermore the map (3.5.1) is a diffeomorphism.
Proof. Let $x \in h_{+}$and $\bar{n} \in \bar{N}$. Since $x$ and $\bar{n} x$ have the same maximal diagonal component (see Definition 2.1) one has $\bar{n} x \in x+\bar{n}$. This proves the first statement of the proposition and establishes the map (3.5.1). But now by (3.4.8) one has $G^{x} \cap N=(1)$ so that (3.5.1) is injective. Now let $y \in h_{+}+\bar{n}$ so that $y=z+v$, where $z \in h_{+}$and $v \in \bar{n}$. But $\mathscr{I}(y)=\mathscr{I}(z)$ by Proposition 17 in [14, p. 369]. Thus $\mathscr{I}(y) \in \mathscr{I}\left(h_{+}\right)$and hence by Proposition 3.3.1 and Lemma 3.3.2 there exists $g \in G$ such that $g z=y$. But now applying the decomposition (2.4.3) to $G_{\mathbb{C}}$ there exist $\sigma \in W, n \in N_{\mathbb{C}}, h \in H_{\mathbb{C}}, \bar{n} \in \bar{N}_{\mathbb{C}}$ such that $g=\bar{n} s(\sigma) h n$. Put $w=(\bar{n})^{-1} y$ so that if $w_{0}$ is the component of $w$ in $d_{0}(=h)$ relative to (2.1.3) then $w_{o}=z$. However, $w=s(\sigma) h n z$ since $g z=y$. But the component of $h n z$ in $d_{0}$ relative to (2.1.3) is also $z$ and hence the component of $w=s(\sigma) h n z$ in $d_{0}$ relative to (2.1.3) is $\sigma z$. Thus $\sigma z=z$. However, $z$ is regular and hence as one knows this implies $\sigma$ is the identity. Thus $s(\sigma)=1$ so that $g=\bar{n} h n$. But then $\bar{n} \in \bar{N}, h \in \tilde{M} H$, and $n \in N$ by (3.4.10). On the other hand the equation becomes $w=h n z$. But the maximal diagonal degree of $w$, and hence of $n z$, is zero. Thus $n z=z$. This implies $n=1$ since $G^{z} \cap N=(1)$ by (3.4.8). But then $y=g z=\bar{n} h z=\bar{n} z$. Thus $y$ is in the image of the map (3.5.1) so that the map is bijective.

Obviously the map (3.4.16) is smooth. If $x \in h_{+}$then the tangent space to $\bar{N} \times \ell_{+}$at $(1, x)$ may be identified with $\left(\bar{n}, h_{)}\right)$and the tangent space to $h_{+}+\bar{n}$ at $x$ may be identified with $h+\bar{n}$. The differential of the map at $(1, x)$ carries $(\bar{n}, h)$ to $h+[\bar{n}, x]$. But $[\bar{n}, x]=\bar{n}$ since $x$ is regular. Thus (3.5.1) is a local diffeomorphism at $(1, x)$. Translation by $\bar{N}$ then implies that there is a local diffeomorphism at all points. Since the map (3.5.1) is bijective it then follows that it is a diffeomorphism.
Q.E.D.

Lemma 3.5.2. Let $w_{o} \in \ell_{+}$and (using the notation of (1.5.4)) put $w=f+w_{o}$. Then there exists a unique element $n_{f}(w) \in \bar{N}$ such that

$$
\begin{equation*}
\bar{n}_{f}(w) w_{o}=w \tag{3.5.2}
\end{equation*}
$$

Furthermore the map

$$
\begin{equation*}
h_{+} \rightarrow \bar{N}, \quad w_{o} \mapsto \bar{n}_{f}(w) \tag{3.5.3}
\end{equation*}
$$

is smooth.
Proof. The first statement follows from Lemma 3.5.1. Consider the map

$$
\begin{equation*}
\ell_{+} \rightarrow h_{+}+\bar{n}, \quad w_{o} \rightarrow w=f+w . \tag{3.5.4}
\end{equation*}
$$

Obviously (3.5.4) is smooth. On the other hand (3.5.3) is just the composite of (3.5.4) with the inverse to (3.5.1) and then with the projection onto $\bar{N}$. This proves (3.5.3) is smooth.
Q.E.D.

Remark 3.5.1. The map (3.5.3) will play an important role in Section 5. We give an explicit formula for $\bar{n}_{f}(w)^{-\mathbf{1}}$ in Section 5.8.

Now by (2.2.9) (for the case where $F=\mathbb{C}$ ) there exists a unique element $m \in H_{\mathbb{C}}$ such that

$$
\begin{equation*}
m e_{\alpha_{i}}=-e_{\alpha_{i}}, \quad i=1, \ldots, l \tag{3.5.5}
\end{equation*}
$$

Clearly $m \in \tilde{M}$. Now let $\tilde{M}=\left\{m_{j}\right\}, j=0, \ldots, 2^{l}-1$, be some ordering of the set $\tilde{M}$, where $m_{0}=m$ and $m_{1}=1$, the identity element of $G$. Also put

$$
\begin{equation*}
H_{j}=m_{j} H, \quad j=0, \ldots, 2^{l}-1 \tag{3.5.6}
\end{equation*}
$$

Lemma 3.5.3. Let $w_{o} \in h_{+}$and as usual put $w=f+w_{o}$. Then if $\bar{n}_{f}(w) \in \bar{N}$ is defined by (3.5.2) one has $\bar{n}_{f}(w) H_{j}\left(\bar{n}_{f}(w)\right)^{-1} \subseteq G^{w}$. In fact if we put

$$
\begin{equation*}
G_{j}^{w}=\bar{n}_{f}(w) H_{j}\left(\bar{n}_{f}(w)\right)^{-1} \tag{3.5.7}
\end{equation*}
$$

then the $G_{j}{ }^{w}$ for $j=0,1, \ldots, 2^{l}-1$ are the cosets of the identity component of $\tilde{G}^{w}$ so that

$$
\begin{equation*}
\tilde{G}^{w}=\bigcup_{j=0} G_{j}^{w} \tag{3.5.8}
\end{equation*}
$$

is a disjoint union.
Proof. Since $w_{o} \in h_{+}$one has $\tilde{G}^{w_{0}}=\tilde{M} H$ by (3.4.7). Thus the $H_{j}$ are the cosets of the identity component of $\tilde{G}^{w_{0}}$. But since $\bar{n}_{f}(w) w_{o}=w$ it follows immediately that the $G^{w}$, defined by (3.5.7), are the cosets of the identity component of $\mathcal{G}^{w}$.
Q.E.D.

Remark 3.5.2. One notes that $G_{1}{ }^{w}$ in the notation of (3.5.8) is in fact the identity component of $\vec{G}^{w}$.

Our interest is not in $G_{1}{ }^{w}$ but in $G_{0}{ }^{w}$. The following result is somewhat surprising (at least to us).

Proposition 3.5. If $w_{0} \in h_{+}$and $w=f-\mid w_{o}$ then in the notation of (3.4.13) and (3.5.8) one has

$$
\begin{equation*}
\tilde{G}^{w} \cap G_{(*)}=G_{\mathbf{0}}{ }^{w} . \tag{3.5.9}
\end{equation*}
$$

Proof. If $w_{o} \in h_{+}$and $w=f+w_{o}$ we may apply the results of Section 2.6 to $G_{\mathbb{C}}{ }^{w}$. That is, if $\left(G_{\mathbb{C}}{ }^{w}\right)_{(*)}=G_{\mathbb{C}}{ }^{w} \cap\left(G_{\mathbb{C}}\right)_{(*)}$, where $\left(G_{\mathbb{C}}\right)_{*}=s(\kappa) \bar{N}_{\mathbb{C}} H_{\mathbb{C}} N_{\mathbb{C}}$, then, by Theorem $2.6,\left(G_{\mathbb{C}}{ }^{w}\right)_{(*)}$ is a nonempty Zariski open subset of $G_{\mathbb{C}}{ }^{w}$.

Now for $0 \leqslant j \leqslant 2^{l}-1$ let

$$
\begin{equation*}
\left(\hbar_{+}\right)_{j}=\left\{w_{o} \in h_{+} \mid G_{j}^{w} \cap G_{(*)} \text { is not empty }\right\} . \tag{3.5.10}
\end{equation*}
$$

If $w_{o} \in\left(\ell_{+}\right)_{j}$ we assert that

$$
\begin{equation*}
\tilde{G}^{w} \cap G_{(*)}=G_{j}^{w} . \tag{3.5.11}
\end{equation*}
$$

In order to prove (3.5.11) we first show that the right side is contained in the left side. Let $a \in G_{j}{ }^{w}$. We must show $a \in G_{(*)}$ since of course $a \in \mathcal{G}^{w}$ by (3.5.8). Now by assumption there exists $g \in G_{j}^{w} \cap G_{(*)}$. But $G_{(*)} \subseteq\left(G_{\mathbb{C}}\right)_{(*)}$. On the other hand if we let $Z_{\mathbb{C}}(\gamma)$, where $\gamma=\mathscr{I}\left(w_{o}\right)$, be defined as in the proof of Theorem 3.2 and let $\left(\beta_{\mathbb{C}}\right)_{(w)}$ be the map $\beta_{(w)}$ of Theorem 2.6 then by Theorem 2.6 one has an algebraic isomorphism

$$
\begin{equation*}
\left(\beta_{\mathbb{C}}\right)_{(w)}:\left(G_{\mathbb{C}}\right)_{(*)} \rightarrow Z_{\mathbb{C}}(\gamma), \tag{3.5.12}
\end{equation*}
$$

where $\left(\beta_{\mathbb{C}}\right)_{(w)}(c)=n(c) w$ for any $c \in\left(G_{\mathbb{C}}\right)_{(*)}$ and $n(c) \in N_{\mathbb{C}}$ is defined by (2.6.2). But now since $g \in G_{(*)}$ one has $n(g) \in N$ and hence if $y=\left(\beta_{c}\right)_{(w)}(g)=n(g) w$ then $y \in g$. Furthermore the maximal diagonal component $y_{1}$ of $y$ is given by (2.6.6). But $h(g) \in H$ so that $h(g)^{-\alpha_{i}}>0$. Thus

$$
\begin{equation*}
\mathscr{I}\left(w_{o}\right)=\gamma \in \mathscr{I}(Z) \quad \text { and } \quad y=\left(\beta_{\mathbb{C}}\right)_{(w)}(g) \in Z(\gamma) . \tag{3.5.13}
\end{equation*}
$$

Now since $g$ and $a$ are in the same connected component of $G^{y}$ one has $g^{-1} a \in G_{e}{ }^{w}$ and hence if $d=n(g)\left(g^{-1} a\right) n(g)^{-1}$ then $d \in G_{e}{ }^{y}=G_{e}{ }^{y}$ (recalling (3.3.4) and (3.4.9)). But $G_{e} \subseteq G_{*}$ by (3.2.9). Thus $d \in\left(G_{\mathbb{C}}\right)_{*}$. But then by (2.6.11) there exists $a^{\prime} \in\left(G_{\mathbb{C}}{ }^{w}\right)_{(*)}$ such that $\psi_{g}\left(a^{\prime}\right)=n(g) g^{-1} a^{\prime} n(g)^{-1}=d$. But then from the injectivity of (2.6.9) one has $a=a^{\prime}$ so that $a \in\left(G_{\mathbb{C}}{ }^{w}\right)_{(*)}$. Thus in particular $a \in s(\kappa) \bar{N}_{\mathbb{C}} H_{\mathbb{C}} N_{\mathbb{C}}$. But then, by Lemma 3.4.4, $s(\kappa)^{-1} a \in \mathcal{G}^{\mathscr{\prime}} \cap$ $\bar{N}_{\mathbb{C}} H_{\mathbb{C}} N_{\mathbb{C}}$. Thus $s(\kappa)^{-1} a \in \bar{N} \tilde{M} H N$ by (3.4.10) so that we can write, using (2.6.2),

$$
\begin{equation*}
h(a)=m^{\prime} h \tag{3.5.14}
\end{equation*}
$$

for $m^{\prime} \in \tilde{M}, h \in H$.

But if $z=\left(\beta_{\mathrm{C}}\right)_{(w)}(a)$ then by the commutative diagram (2.6.21) one has $z=\left(\beta_{\mathrm{C}}\right)_{y}(d)=\beta_{y}(d)$. But then $z \in Z(\gamma)$ by (3.2.15) since $d \in G_{e}{ }^{y}$. Thus if $z_{1}$ is the maximal diagonal component of $z$ then $z_{1}=\sum r_{i} e_{\alpha_{i}}$, where $r_{i}>0$. But by (2.6.6) one has $r_{i}=h(a)^{-\alpha_{i}}$. However, if $m^{\prime}$ in (3.5.14) is not the identity in $\tilde{M}$ there clearly exists $1 \leqslant j \leqslant l$ such that $m^{\prime} e_{\alpha_{j}}=-e_{\alpha_{j}}$. This contradicts the positivity of $r_{j}$. Hence $m^{\prime}=1$ so that $a \in s(\kappa) \widetilde{N} H N=G_{(*)}$. Thus we have proved the right side of (3.5.11) is contained in the left side. We have also shown that if $w_{o} \in\left(h_{+}\right)_{j}$ then

$$
\begin{equation*}
\left(\beta_{\mathbb{C}}\right)_{(w)}\left(G_{j}^{w}\right) \subseteq Z(\gamma) \quad \text { for } \quad \gamma=\mathscr{I}(w) \tag{3.5.15}
\end{equation*}
$$

Now assume that $b$ is in the left side of (3.5.11). Thus $b \in G_{k}{ }^{w}$ for some $k$. Let $g$ be as above. We have to show that $k=j$. But now $b \in G^{w} \cap G_{(*)} \subseteq\left(G_{\mathbb{C}}{ }^{w}\right)_{(*)}$ and $\left(\beta_{\mathbb{C}}\right)_{(w)}(b)=n(b) w \in Z_{\mathbb{C}}(\gamma)$, where $\gamma=\mathscr{I}(w)$. But $n(b) \in N$ so that if $z=\left(\beta_{\mathbb{C}}\right)_{(w)}(b)$ then $z \in g$. On the other hand if $z_{1}$ is the maximal diagonal component of $z$ then $z_{1}=\sum h(b)^{-\alpha_{i}} e_{\alpha_{i}}$ by (2.6.6). But $h(b) \in H$ so that $h(b)^{-\alpha_{i}}>0$. Thus $z \in Z(\gamma)$ and hence

$$
\begin{equation*}
\left(\beta_{\mathbb{C}}\right)_{(w)}(b) \in Z(\gamma) . \tag{3.5.16}
\end{equation*}
$$

Now let $a \in\left(G_{\mathbb{C}}{ }^{w}\right)_{(*)}$ be any element such that $\left(\beta_{\mathbb{C}}\right)_{(w)}(a) \in Z(\gamma)$. Put $v=$ $\left(\beta_{\mathbb{C}}\right)_{(w)}(a)$. Then by the commutative diagram (2.6.21) one has $v=\left(\beta_{\mathbb{C}}\right)_{y}(d)$ for $d=\psi_{g}(a)$ using the notation of (3.5.13), where of course $d \in\left(G_{\mathbb{C}}{ }^{y}\right)_{*}$ and $y=n(g)$ w. But by Theorem 3.2 (see (3.2.15)) there exists $d^{\prime} \in G_{e}{ }^{y}$ such that $v=\beta_{y}\left(d^{\prime}\right)$. But since $\left(\beta_{\mathfrak{c}}\right)_{y}$ is an isomorphism one has $d^{\prime}=d$ or $d=\psi_{g}(a)=$ $n(g) g^{-1} a n(g)^{-1} \in G_{e}{ }^{y}$. Thus $g^{-1} a \in \mathcal{G}_{e}{ }^{w}$. Hence $a \in G_{j}^{w}$. Thus we have proved

$$
\begin{equation*}
\left(\beta_{\mathbb{C}}\right)_{(w)}^{-1}(Z(\gamma)) \subseteq G_{j}^{w} \quad \text { in case } \quad w_{o} \in\left(\ell_{+}\right)_{j} \tag{3.5.17}
\end{equation*}
$$

Thus $b \in G_{j}{ }^{w}$ by (3.5.16) so that $j=k$. This proves the assertion.
Now for any $w_{o} \in h_{+}$and $0 \leqslant j \leqslant 2^{l}-1$ let

$$
\begin{equation*}
m_{j}(w)=\bar{n}_{f}(w) m_{j} \bar{n}_{f}(w)^{-1} . \tag{3.5.18}
\end{equation*}
$$

Since $m_{j} \in H_{j}$ one has

$$
\begin{equation*}
m_{j}(w) \in G_{j}^{w} \tag{3.5.19}
\end{equation*}
$$

by (3.5.7). Now let $\sigma_{j}$ be the map

$$
\begin{equation*}
\sigma_{j}: \ell_{+} \rightarrow \hat{G}, \quad w_{o} \rightarrow m_{j}(w) \tag{3.5.20}
\end{equation*}
$$

It follows from the smoothness of (3.5.3) that $\sigma_{j}$ is smooth. We assert

$$
\begin{equation*}
\left(\kappa_{+}\right)_{j}=\sigma_{j}^{-1}\left(G_{(*)}\right) . \tag{3.5.21}
\end{equation*}
$$

Indeed if $w_{o} \in\left(h_{+}\right)_{j}$ then $G_{j}{ }^{w} \subseteq G_{(*)}$ by (3.5.11). Thus $\sigma_{j}\left(w_{o}\right) \in G_{(*)}$. Conversely if $\sigma_{j}\left(w_{o}\right) \in G_{(*)}$ then $G_{i}{ }^{w} \cap G_{(*)}$ is not empty. Hence $w_{o} \in\left(h_{+}\right)_{j}$ by definition. This proves (3.5.21). We next assert that

$$
\begin{equation*}
\bigcup_{j=0} \mathscr{I}\left(\left(\hbar_{+}\right)_{j}\right)=\mathscr{I}(Z) . \tag{3.5.22}
\end{equation*}
$$

We first observe that the left side of (3.5.22) is contained in the right side of (3.5.13). Now let $\gamma \in \mathscr{I}(Z)$. Then $\gamma=\mathscr{I}\left(w_{o}\right)$ for some $w \in h_{+}$by (3.3.5). Now let $y \in Z(\gamma)$. By Theorem 2.6 there exists $g \in\left(G_{\mathbb{C}}{ }^{w}\right)_{(*)}$ such that $\left(\beta_{\mathbb{C}}\right)_{(w)}(g)=$ $n(g) w=y$, where, of course, $n(g) \in N_{\mathrm{C}}$. But $w, y \in \mathcal{l}_{f}$ so that by Proposition 2.3.2 one has $n(g) \in N$. Dut now using the notation of (2.6.0) one has $h(g)^{-\alpha_{i}}>0$ since $y_{1} \in \tilde{d}_{1}$. Thus $h(g) \in H$. Furthermore recalling (2.6.7) and its notation one must have $h(g) y, \kappa^{-1} w_{o}+e \in e+b$ and by (2.6.7) these elements are conjugate under $\bar{n}(g)$. But then $\bar{n}(g) \in \bar{N}$ by Proposition 2.3.2 with a reversal of positive and negative roots and a reversal of $e$ and $f$. Thus $s(\kappa)^{-1} g=$ $\bar{n}(g) h(g) n(g) \in G$. Consequently $g \in \mathcal{G}^{w} \cap G_{(*)}$. But then $g \in G_{j}{ }^{w} \cap G_{(*)}$ for some $j$ by (3.5.8). But then $w_{o} \in\left(h_{+}\right)$, by definition and hence $\gamma=\mathscr{I}\left(w_{o}\right) \in$ $\mathscr{I}\left(\left(\kappa_{+}\right)_{j}\right)$. This proves equality (3.5.22).
We next assert that $\left(\ell_{+}\right)_{j}$, for any $j$ is closed in $\ell_{+}$. Indeed let $w_{0} \in\left(\ell_{+}\right)_{j} \cap \hbar_{+}$. We first show that

$$
\begin{equation*}
\left(G_{\mathrm{C}}^{*}\right)_{(*)} \cap G_{j}^{* *} \quad \text { is not empty. } \tag{3.5.23}
\end{equation*}
$$

Indeed $\left(G_{\mathbb{C}}{ }^{w}\right)_{(*)}$ is not an empty Zariski open subset of $G_{\mathbb{C}}{ }^{w}$ by Theorem 2.6. On the other hand $g_{\mathbb{C}}{ }^{w}$ is the complexification of $g^{w}$ so that for any $g \in G_{j}{ }^{w}$ the tangent space to $G_{j}{ }^{w}$ at $g$ is a real form of the tangent space to $G_{\mathbb{C}}{ }^{w}$ at $g$. Thus any regular function on $G_{\mathrm{C}}{ }^{w}$ which vanishes on $G_{j}{ }^{w}$ must vanish on $G_{\mathrm{C}}{ }^{w}$. Thus $G_{j}{ }^{w}$ is Zariski dense in $G_{\mathbb{C}}{ }^{w}$. This proves (3.5.23).
Let $g \in\left(G_{\mathbb{C}}{ }^{w}\right)_{(*)} \cap G_{j}{ }^{w}$. But then $s(\kappa)^{-1} g \in G \in \bar{N}_{\mathbb{C}} H_{\mathbb{C}} N_{\mathbb{C}}$ and hence by Lemma 3.4.3 one has $\bar{n}(g) \in \bar{N}, n(g) \in N$ and we can write $h(g)=m^{\prime} h$, where $m^{\prime} \in \tilde{M}$ and $h \in H$. Now if $m^{\prime}=1$ then $g \in G_{(*)}$ and hence $G_{j}{ }^{w} \cap G_{(*)}$ is not empty. But then $w_{o} \in\left(h_{+}\right)_{j}$. Thus in order to prove that $\left(h_{+}\right)_{j}$ is closed in $h_{+}$ it suffices to prove that $m^{\prime}=1$. Assume $m^{\prime} \neq 1$. Then there exists $1 \leqslant k \leqslant l$ such that $m^{\prime} e_{\alpha_{k}}=-e_{\alpha_{k}}$. But then if $y=n(g) w$ and we use the notation of (2.6.6) one has $y_{1}=\sum r_{i} e_{\alpha_{i}}$, where $r_{k}<0$. Now let $w_{o}^{(m)} \in\left(\hbar_{+}\right)_{j}, m=1,2, \ldots$, be a sequence such that $w_{o}^{(m)}$ converges to $w_{o}$. Put $w^{(m)}=f+w_{0}^{(m)}$. But by (3.5.7) there exists $g_{0} \in H_{j}$ so that $g=\bar{n}_{f}(w) g_{o} \bar{n}_{f}(w)^{-1}$. Put $g^{(m)}=$ $\bar{n}_{f}\left(w^{(m)}\right) g_{0} \bar{n}_{f}\left(w^{(m)}\right)^{-1}$ so that $g^{(m)}$ converges to $g$. On the other hand $g^{(m)} \in G_{j}^{w(m)}$ by (3.5.7). But since $w_{o}^{(m)} \in\left(h_{+}\right)_{j}$ it follows that $g^{(m)} \in \mathcal{G}^{w^{(m)}} \cap G_{(*)} \subseteq\left(G_{\mathrm{C}}^{(t m}\right)(*)$
by (3.5.11). Furthermore if $y^{(m)}=\left(\beta_{\mathbb{C}}\right)_{\left(w^{(m)}\right)}\left(g^{(m)}\right)-n\left(g^{(m)}\right) w^{(m)}$ then $y^{(m)} \in Z$ by Theorem 2.6. But clearly $y^{(m)}$ converges to $y$ since $g^{(m)}$ and $w^{(m)}$ respectively converge to $g$ and $w$. Thus if $y_{i}^{(m)}=\sum_{i} r_{i}^{(m)} e_{\alpha_{i}}$ is the maximal diagonal component of $y^{(m)}$ then $r_{k}^{(m)}$ converges to $r_{k}<0$. But $r_{k}^{(m)}>0$ since $y^{(m)} \in Z$. This is a contradiction and hence $\left(h_{+}\right)_{j}$ is closed in $\ell_{+}$.

But now by (3.5.7), and recalling Lemma 3.5.2, $\left(f_{+}\right)_{j}$ is open in $\ell_{+}$since $G_{(*)}$ is open in $G$. But clearly $h_{+}$is connected. Hence for any $1 \leqslant j \leqslant 2^{l}-1$ either $\left(h_{+}\right)_{j}$ is empty or $\left(h_{+}\right)_{j}=h_{+}$. Furthermore clearly by (3.5.11) the sets $\left(h_{+}\right)_{j}$, over all $j$, are mutually disjoint. But by (3.5.22), since $Z$ is not empty, it follows that there exists a unique $0 \leqslant k \leqslant 2^{l}-1$ such that

$$
\begin{equation*}
\left(h_{+}\right)_{k}=h_{+} \quad \text { and } \quad\left(h_{+}\right)_{j} \text { is empty if } j \neq k \tag{3.5.24}
\end{equation*}
$$

It follows then from (3.5.11) that

$$
\begin{equation*}
G^{w} \cap G_{(*)}=G_{k}^{w} \quad \text { for any } \quad w_{o} \in h_{+} \tag{3.5.25}
\end{equation*}
$$

One also notes that (3.5.22) and (3.5.24) imply that

$$
\begin{equation*}
\mathscr{I}\left(h_{+}\right)=\mathscr{I}(Z) \tag{3.5.26}
\end{equation*}
$$

It remains only to prove that $k=0$.
Now if $h_{\Phi}$ is given by (3.1.14) then for any $x \in h_{\mathbb{C}}$ one has $Q\left(x, h_{\varphi}\right)=$ $\langle\varphi, x\rangle Q\left(e_{\varphi}, e_{-\varphi}\right)$ by the invariance of $Q$. However, $Q\left(e_{\varphi}, e_{-\varphi}\right)=1$ by (1.5.2). Thus if we consider the isomorphism

$$
\begin{equation*}
h_{\mathbb{C}} \rightarrow \ell_{\mathbb{C}}^{\prime} \tag{3.5.27}
\end{equation*}
$$

defined by $Q$ then

$$
\begin{equation*}
h_{\Phi} \rightarrow \varphi \tag{3.5.28}
\end{equation*}
$$

with respect to (3.5.27). Now one knows that the root normals $h_{\alpha_{i}}$ are a basis of $h$. Let $x_{o} \in h$ be as in Section 2.1 so that $\left\langle\alpha_{i}, x_{o}\right\rangle=1$. Let $s_{i} \in \mathbb{R}$ be such that $x_{o}=\sum s_{i} h_{\alpha_{i}}$. We assert that

$$
\begin{equation*}
s_{i}>0 \tag{3.5.29}
\end{equation*}
$$

Indeed one identifies $h_{\mathbb{C}}$ with $h_{\mathbb{C}}^{\prime}$ by (3.5.27) and if we write $x_{o}=\sum t_{i} \alpha_{i}$ then as one knows (see, e.g., Lemma 8.3 in [6, p. 166]), $t_{i}>0$ since $Q\left(x_{o}, \alpha_{i}\right)>0$. However, one has $s_{i}=t_{i}$ by (3.5.28). But this implies (3.5.29). (See also the argument in Lemma 15 in [18, p. 790.])

Now put

$$
\begin{equation*}
e_{o}=\sum s_{i} e_{\alpha_{i}} \tag{3.5.30}
\end{equation*}
$$

Since, as one knows, $\left[e_{\alpha_{i}}, e_{-\alpha_{j}}\right]=0$ if $i \neq j$ it follows that $\left[e_{o}, f\right]=x_{0}$ so that if $a$ is the $\mathbb{R}$-span of $x_{o}, e_{o}$, and $f$ then $a$ is a Lie subalgebra whose complexification $a_{\mathbb{C}}$ is a principal TDS of $g_{\mathbb{C}}$ in the notation of [12]. See Section 5 in [12]. Furthermore we note that $e_{o} \in \tilde{d_{1}}$ by (3.5.30) and hence there exists $a \in H$ such that

$$
\begin{equation*}
a e=e_{o} \tag{3.5.31}
\end{equation*}
$$

where $e$ is given by (2.6.18).
Now let $A$ and $A_{\mathbb{C}}$ be the subgroups of $G_{\mathbb{C}}$ corresponding to $a$ and $a_{\mathbb{C}}$. Since all the eigenvalues of ad $x_{0}$ are integers it follows that $A_{\mathbb{C}}$ is isomorphic to the adjoint group of $a_{\mathbb{C}}$. In particular all of our previous results and notation with respect to $g, g_{\mathbb{C}}, G$, and $G_{\mathbb{C}}$ may be applied to $a, a_{\mathbb{C}}, A$, and $A_{\mathbb{C}}$. To indicate that such an application is being made we will use $L(A)$ for any previously defined notation $L$ (e.g., $L=\bar{B}, N$, etc.).

Now recall that $m \in \hat{M}$ is defined by the relations $m e_{\alpha_{j}}=-e_{\alpha_{j}}$ for all $j$. But then clearly $m=\exp i \pi x_{o}$. Thus $m \in A_{\mathbb{C}}$. Now $x_{o} \in h_{+}$and hence using our previous notation, where $w_{a} \in h_{+}$, we now fix $w_{o}$ so that $w_{o}=x_{a}$. Note then that $w_{o}$ and $w=f+w_{o}$ are both in $a_{0}$. If we now put $\bar{n}(A)=\mathbb{R} f$ and recall that if $\bar{n}_{f}(w) \in \bar{N}$ is the unique element given by (3.5.2) (so that $\bar{n}_{f}(w) w_{o}=w$ ) one clearly has

$$
\begin{equation*}
1 \neq \bar{n}_{f}(w) \in \bar{N}(A) . \tag{3.5.32}
\end{equation*}
$$

But now since obviously $m f=-f$ it follows that $\bar{n}_{f}(w) m=m\left(\bar{n}_{f}(w)\right)^{-1}$. Thus if we put

$$
\begin{equation*}
g=\bar{n}_{f}(w) m \bar{n}_{f}(w)^{-1} \tag{3.5.33}
\end{equation*}
$$

one has $g=\left(\bar{n}_{f}(w)\right)^{2} m$ and hence $g \in \bar{N}(A) \tilde{M}(A)$ where $h(A)=\mathbb{R} x_{o}$ so that

$$
\begin{equation*}
\{1, m\}=\tilde{M}(A) \tag{3.5.34}
\end{equation*}
$$

Now $g \in \tilde{G}(A)$ and hence there exist $c, c^{\prime}, c^{\prime \prime} \in \mathbb{R}$ such that $g e_{o}=c^{\prime \prime} e_{o}+$ $c^{\prime} x_{0}+c f$. We assert that

$$
\begin{equation*}
c>0 \tag{3.5.35}
\end{equation*}
$$

In fact for any $\bar{n} \in \bar{N}(A)$ it is clear that we may find $d(\bar{n}), d^{\prime} \in \mathbb{R}$ such that $\bar{n} e_{o}=e_{o}+d^{\prime} x_{o}+d(\bar{n}) f$. We first prove that

$$
\begin{equation*}
d(\bar{n}) \leqslant 0 \quad \text { and } \quad d(\bar{n})<0 \quad \text { if and only if } \quad \bar{n} \neq 1 \tag{3.5.36}
\end{equation*}
$$

Indeed write $\bar{n}=\exp r f$, where $r \in \mathbb{R}$. However, $\left[f,\left[f, e_{0}\right]\right]=\left[x_{0}, f\right]=-f$.

Thus $d(\bar{n})=-r^{2} / 2$, which proves (3.5.36). But now $g=\left(\bar{n}_{f}(w)\right)^{2} m$, where $\bar{n}_{f}(w)^{2} \neq 1$ by (3.5.32). However, $m e_{o}=-e_{o}$. Thus (3.5.36) implies (3.5.35).

Now let $n(A)=\mathbb{R} e_{o}$. Furthermore from the bracket relations of $f, x_{o}$, and $e_{o}$ it is clear that we can introduce a compact form of $a_{\mathbb{C}}$ and hence a *-operation, written $x \rightarrow x^{(*)}$ (to avoid confusion with $*$-operation in $g c$ ), so that $e_{o}^{(*)}=f$. But then we can fix $e(A)=e_{o}$ and $f(A)=f$. Put $s=s(\kappa)(A)$ so that $s x_{o}=-x_{o}, s e_{o}=f$. If $a \in H$ is defined by (3.5.31) we assert that

$$
\begin{equation*}
s a=s(\kappa) . \tag{3.5.37}
\end{equation*}
$$

Indeed if $h=s(\kappa)^{-1}$ sa then clearly $h x_{o}=x_{o}$ since one of course has $\kappa x_{o}=-x_{o}$. But then $h \in G_{\mathbb{C}^{x_{0}}}=H_{\mathbb{C}}$ (recalling (3.4.8)). However, he $=e$ so that $h^{\alpha_{i}}=1$ for $i$. Thus $h=1$, proving (3.5.37).

But now the Gelfand-Bruhat decomposition of $A_{\mathbb{C}}$ becomes the disjoint union

$$
\begin{equation*}
A_{\mathbb{C}}=N_{\mathbb{C}}(A) s B_{\mathbb{C}}(A) \cup B_{\mathbb{C}}(A) \tag{3.5.38}
\end{equation*}
$$

However, $B_{\mathbb{C}}(A)$ normalizes $\mathbb{C} e$ so that by (3.5.38) one must have $g \in N_{\mathbb{C}}(A) s B_{\mathbb{C}}(A)$. Thus there exist $n^{\prime}, n \in N_{\mathbb{C}}(A)$ and $h^{\prime} \in H_{\mathbb{C}}(A)$ so that $g=n^{\prime} s h^{\prime} n$. Thus if $\bar{n}_{1}=s^{-1} n^{\prime} s$ one has $\bar{n}_{1} \in \bar{N}_{\mathbb{C}}$ and $s^{-1} g=\bar{n}_{1} h^{\prime} n$. But now $s^{-1} g \in \tilde{G}(A)$ so that by (3.4.11) one has $\bar{n}_{1} \in \bar{N}(A) n \in N(A)$ and $h^{\prime}=m^{\prime} h_{1}$, where $h_{1} \in H(A)$ and $m^{\prime}$ is either $m$ or 1 by (3.5.34). Put $b=s^{-1} g$ and write $b e_{o}=r e_{o}+r^{\prime} x_{o}+r^{\prime \prime} f$, where $r, r^{\prime}, r^{\prime \prime} \in \mathbb{R}$. Since $s^{-1} f=e_{o}$ it then follows from (3.5.35) that $r>0$. However, $b=\bar{n}_{1} h^{\prime} n$. Hence $b e_{o}=\bar{n}_{1} h^{\prime} e_{o}=\left(h^{\prime}\right)^{\alpha} \bar{n}_{1} e_{o}$, where $\alpha=\alpha_{1}(A)$. But then $r=\left(h^{\prime}\right)^{\alpha}$ since $\bar{n}_{1} e \in e+\bar{\theta}(A)$. Thus $\left(h^{\prime}\right)^{\alpha}>0$ so that $m^{\prime}=1$. But then $s^{-1} g=\bar{n}_{1} h_{1} n$. However, $s^{-1}=a s(\kappa)^{-1}$ by (3.5.37). Thus if we put $\bar{n}=a^{-1} \bar{n}_{1} a$ and $h=a^{-1} h_{1}$ then $\bar{n} \in \bar{N}$ and $h \in H$ since of course $\bar{N}(A) \subseteq \bar{N}$ and $H(A) \subseteq H$ and one has $s(\kappa)^{-1} g=\bar{n} h n$. But then $g \in G_{(*)}$ by definition (see (3.2.1)). However, recalling (3.5.6) one has $m=m_{0} \in H_{0}$. But then $g \in G_{0}{ }^{w}$ by (3.5.7) and (3.5.33). Thus $G_{0}{ }^{w} \cap G_{(*)}$ is not empty. Hence by definition $x_{0}=w_{o} \in\left(h_{+}\right)_{0}$ so that $\left(h_{+}\right)_{0}$ is not empty. Thus $k=0$ by (3.5.24). The result then follows, as noted, by (3.5.25).
Q.E.D.
3.6. The following is our main result in Section 3. Among other things it provides the structure for the integration of Hamilton's equation in Section 7.

Theorem 3.6. Let $g$ be a split real semi-simple Lie algebra. Let $\mathbb{G}$ be the normalizer of $g$ in the complexified adjoint group of $g$. Let $l=\operatorname{rank} g$ and let $\mathscr{I}: g \rightarrow \mathbb{R}^{l}$ be the map defined by the l fundamental polynomial invariants on $g$. (See (2.3.1).)

Now let $h_{+}$be the open Weyl chamber of a split Cartan subalgebra $h_{h}$ defined
as in (3.3.1). Let $\Pi=\left\{\alpha_{1}, \ldots, \alpha_{l}\right\}$ be the corresponding set of simple positive roots and let $e_{\alpha_{i}}, e_{-\alpha_{i}}$ be corresponding root vectors, normalized as in Section 1.5. Let $\kappa$ be that element of the Weyl group $W(g, h)$ such that $\kappa \Pi=-\Pi$. Let $s(\kappa) \in G$ be that unique element in the normalizer of $h$ such that $s(\kappa) e_{\alpha_{i}}=e_{\kappa \alpha_{i}}$. Let $n$ and $\bar{n}$ be the nilpotent subalgebras generated respectively by the $e_{\alpha_{i}}$ and the $e_{-\alpha_{i}}$ and let $\bar{N}, H$, and $N$ be the subgroup of $\bar{G}$ corresponding respectively to $\bar{N}, H$, and $N$. Also let $f=\sum e_{-\alpha_{i}}$ and let $Z$, as in Theorem 3.2, be the $2 l$-dimensional manifold of all normalized Jacobi elements. That is, all elements $y \in g$ of the form $y=f+x+\sum a_{i} e_{\alpha_{i}}$, where $x \in h$ and $a_{i}>0$.

Now let $w_{o} \in h_{+}$and put $w=f+w_{0}$. Then for any $g$ in the centralizer $\hat{G}^{w}$ of $w$ in $G$ there exist uniquely nonzero real numbers $r_{i}(g), i=1, \ldots, l$, such that $g e_{\alpha_{i}}-r_{i}(g) e_{\alpha_{i}} \in h_{h}+\bar{n} . P u t$

$$
\begin{equation*}
G_{0}{ }^{w}=\left\{g \in \mathbf{G}^{w} \mid \text { all } r_{i}(g)<0\right\} . \tag{3.6.1}
\end{equation*}
$$

This set may also be given by (3.5.7). Then for any $g \in G_{0}{ }^{w}$ there exist uniquely $\bar{n}(g) \in N, h(g) \in H$, and $n(g) \in N$ such that $g=s(\kappa) \bar{n}(g) h(g) n(g)$. Furthermore if $\gamma=\mathscr{I}\left(w_{o}\right) \in \mathbb{R}^{l}$ and $Z(\gamma)=\mathscr{I}^{-1}(\gamma) \cap Z$ then $n(g) w \in Z(\gamma)$ defining a map

$$
\begin{equation*}
\beta_{(w)}: G_{0}{ }^{w} \rightarrow Z(\gamma) . \tag{3.6.2}
\end{equation*}
$$

Moreover $G_{0}{ }^{w}$ and $Z(\gamma)$ are manifolds and (3.6.2) is a diffeomorphism. In fact both manifolds are diffeomorphic to $\mathbb{R}^{l}$. Finally

$$
\begin{equation*}
Z=\bigcup Z(\gamma) \tag{3.6.3}
\end{equation*}
$$

is a disjoint union where, writing $\gamma=\mathscr{I}\left(w_{o}\right)$, the union is over the open Weyl chamber $h_{+}$as an index set.

Proof. Now by (3.5.6) and (3.5.7) one has $\mathcal{G}^{w}=\bar{n}_{f}(w) \tilde{M} H \bar{n}_{f}(w)^{-1}$. If $g \in G^{w}$ then by (3.5.8) there exist $0 \leqslant j \leqslant 2^{l}-1$ and $h \in H$ such that $g=$ $\bar{n}_{f}(w) m_{j} h \bar{n}_{f}(w)^{-1}$. Let $d_{i j}=\{0,1\}, i=1, \ldots, l$, be such that $m_{j} e_{\alpha_{i}}=(-1)^{d_{i j}} e_{\alpha_{i}}$. But since $\bar{n} e_{\alpha_{i}}-e_{\alpha_{i}} \in h+\bar{n}=\bar{b}$ it follows that for

$$
\begin{equation*}
r_{i}(g)=(-1)^{\alpha_{i j}} h^{\alpha_{i}} \tag{3.6.4}
\end{equation*}
$$

one has $g e_{\alpha_{i}}-r_{i}(g) e_{\alpha_{i}} \in h+\bar{n}$. This proves the first statement. Furthermore by (3.6.4) it follows that $r_{i}(g)<0$ if and only if $d_{i j}=1$ for $i=1, \ldots, l$. But then by definition of $m=m_{0}$ this is the case if and only if $j=0$. Thus $G_{0}{ }^{w}$ can be given by (3.6.1). Next the existence and uniqueness of the decomposition $g=s(\kappa) \bar{n}(g) h(g) n(g)$ for $g \in G_{0}{ }^{w}$ are given by (3.5.9) and (3.4.14). But then $G_{0}{ }^{w} \subseteq\left(G_{\mathbb{C}}{ }^{w}\right)_{(*)}$, recalling (2.6.4), and hence (see (2.6.15)) the map $\left(\beta_{\subsetneq}\right)_{(w)}$
is defined on $G_{0}{ }^{\text {m }}$. By (3.5.15), recalling (3.5.24), one has $\left(\beta_{\mathbb{C}}\right)_{(w)}\left(G_{0}{ }^{w}\right) \subseteq Z(\gamma)$ so that (3.6.2) is defined where

$$
\begin{equation*}
\beta_{(w)}=\left(\beta_{\mathcal{C}}\right)_{(w)} \mid G_{0}{ }^{w} . \tag{3.6.5}
\end{equation*}
$$

But then recalling that (3.5.12) is bijective it follows from (3.5.17), for $j=0$, that (3.6.2) is also bijective. However, $Z(\gamma)$ is clearly an $l$-dimensional submanifold of $Z_{\mathbb{c}}(\gamma)$ by Theorem 3.2 and furthermore one has a diffeomorphism $Z(\gamma) \cong \mathbb{R}^{l}$ by Theorem 3.2. Also, $G_{0}{ }^{* w}$ is clearly a submanifold of $\left(G_{\mathbb{C}}{ }^{w}\right)(*)$ and $G_{0}{ }^{u} \cong \mathbb{R}^{l}$ by (3.5.7). Thus (3.6.2) is a diffeomorphism by (3.6.5) since $\left(\beta_{\mathbb{C}}\right)_{(w)}$ is an algebraic isomorphism, recalling (2.6.15), and hence a diffeomorphism. The relation (3.6.3) now follows from (2.3.3), (3.5.26), and Proposition 3.3.1.
Q.E.D.

As in the complex case, we will refer to the submanifolds $Z(\gamma)$ of $Z$ (now in the real case) for $\gamma \in \mathscr{I}\left(h_{+}\right)$, as the isospectral leaves of $Z$.

Remark 3.6. Let $w_{o} \in h_{+}, w=f+w_{o}, \gamma=\mathscr{I}\left(w_{o}\right)$, and $g \in G_{0}{ }^{w}$. Since $G_{0}{ }^{w} \subseteq\left(G_{\mathbb{C}}{ }^{w}\right)_{(*)}$ the diffeomorphism $\psi_{g}:\left(G_{\mathbb{C}}{ }^{w}\right)_{(*)} \rightarrow\left(G_{\mathbb{C}}{ }^{y}\right)_{*}$, where $y=n(g) w \in$ $Z(\gamma)$ is defined on $G_{0}{ }^{w}$. If $a \in\left(G_{\mathbb{C}}{ }^{w}\right)_{(*)}$ we recall that $\psi_{g}(a)=n(g) g^{-1} a n(g)^{-1}$. It follows then from Remark 3.5 .2 that $\psi_{g}\left(G_{0}{ }^{w}\right)=G_{e}{ }^{y}$. Furthermore since, recalling Theorem 3.2 and its proof,

$$
\begin{equation*}
\beta_{y}=\left(\beta_{\mathbb{C}}\right)_{y} \mid G_{e}{ }^{y} \tag{3.6.6}
\end{equation*}
$$

the commutative diagram (2.6.21) of algebraic isomorphisms, upon restriction, becomes a commutative diagram

where all three maps are diffeomorphisms. If $g, a \in G_{0}{ }^{w}$ and $\psi_{g}(a)=$ $n(g) g^{1} a n(g){ }^{1}=d \in G_{e}{ }^{y}$ then by (2.6.13) one has

$$
\begin{equation*}
n(a)=n_{d} n(g) \tag{3.6.8}
\end{equation*}
$$

where now $n(a), n(g)$, and $n_{d}$ are in $N$.
One consequence of (3.6.3), or rather its proof, is the answer to the following question (which perhaps is of independent interest): What is a necessary and sufficient condition that an element $x \in g$ "can be put" in Jacobi form? The answer is: The element $x$ should be a regular real semi-simple element. That is,

Prorosition 3.6. Let $x \in g . A$ necessary and sufficient condition that $x$ be $G$-conjugate to an element in $Z$ is that $x \in R_{+}$. (See Section 3.3.)

Proof. One has $Z \subseteq R_{+}$by (3.3.4). Thus the condition is necessary. But if $x \in R_{+}$then $x$ is $G$-conjugate to an element $w_{o} \in h_{+}$by Lemma 3.3.2. However, recalling (3.5.26) one has

$$
\mathscr{I}(Z)=\mathscr{I}\left(h_{+}\right) .
$$

Thus there exists $y \in Z$ such that $\mathscr{I}(y)-\mathscr{I}\left(w_{o}\right)$. Hence $y$ is $G$-conjugate to an element in $h_{+}$by Lemma 3.3.2. This element must be $w_{o}$ by Proposition 3.3.1. Thus $x$ and $y$ are $G$-conjugate.
Q.E.D.
3.7. Finally, in this section we show that the space $Z$ of normalized Jacobi elements may be smoothly parametrized by the product $H \times h_{+}$of the split Cartan subgroup $H$ and open Weyl chamber $\hbar_{+}$.

Theorem 3.7. Let $h_{+}, \bar{N}, H, N, \widehat{G}$, and $f$ be as in Theorem 3.6. Let $\left(g_{o}, w_{o}\right) \in$ $H \times h_{+}$and let $w=f+w_{o}$. Let $\bar{n}_{f}(w) \in \bar{N}$ be as in (3.5.2) so that $\bar{n}_{f}(w) w_{o}=w$. Let $m \in G$ be as in Section 3.5 so that $m$ fixes $H$ and $m e_{\alpha_{i}}=-e_{\alpha_{i}}, i=1, \ldots, l$. Put

$$
\begin{equation*}
g=\bar{n}_{f}(w) m g_{o} \bar{n}_{f}(w)^{-1} \tag{3.7.1}
\end{equation*}
$$

so that $g \in G_{0}{ }^{w}$ by (3.5.7). Let $n(g) \in N$ be as in Theorem 3.6 so that $n(g) w \in Z$ by Theorem 3.6. Now let

$$
\begin{equation*}
H \times n_{+} \rightarrow Z \tag{3.7.2}
\end{equation*}
$$

be the map defined by $\left(g_{o}, w_{o}\right) \mapsto n(g) w$. Then (3.7.2) is a diffeomorphism.
Proof. It is obvious using Lemma 3.5.2 and (3.7.1) that the map

$$
\begin{equation*}
H \times \ell_{+} \rightarrow G_{0}^{w}, \quad\left(g_{o}, w_{o}\right) \mapsto g \tag{3.7.3}
\end{equation*}
$$

is smooth. See (3.5.3). But then the smoothness of (3.4.15) implies that the map $I I \times \ell_{+} \rightarrow N,\left(g_{o}, w_{o}\right) \mapsto n(g)$ is smooth. It follows immediately then that (3.7.2) is smooth. But now if $\mathscr{I}\left(w_{o}\right)=\gamma$ then, by Theorem 3.6, the restriction of (3.7.2) to $H \times\left\{w_{0}\right\}$ maps $H \times\left\{w_{o}\right\}$ bijectively onto $Z(\gamma)$. But then (3.7.2) is a bijection since (3.6.3) is, by Theorem 3.6, a disjoint union. To prove Theorem 3.7 it therefore suffices to show that (3.7.2) is a local diffeomorphism. Let $T, U$, and $V$ be the tangent space to $H \times \ell_{+}, H \times\left\{w_{o}\right\}$, and $\left\{g_{o}\right\} \times h_{+}$, respectively, at $\left(g_{o}, w_{o}\right)$ so that $T=U \oplus V$. If $\sigma$ denotes the map (3.7.2) it suffices then, by dimension, to show that its differential $\sigma_{*}$ is injective on $T$. But now clearly the restriction $\sigma \mid H \times\left\{w_{o}\right\}=\beta_{(w)} \circ \tau$, where $\tau$ is the
restriction of (3.7.3) to $H \times\left\{w_{0}\right\}$ and $\beta_{(w)}$ is given by (3.6.2). However, $\beta_{(w)}$ is a diffeomorphism by Theorem 3.6 and $\tau$, given by group multiplication (see (3.7.1)), is obviously a diffeomorphism. Thus $\sigma$ induces a diffeomorphism $H \times\left\{w_{o}\right\} \rightarrow Z(\gamma)$. However, $Z(\gamma)$ is an $l$-dimensional submanifold of $Z$ by Theorem 3.2. Thus $\sigma_{*} \mid U$ is injective. Now clearly we may identify $V$ with the tangent space to $h_{+}$at $w_{o}$. Let $0 \neq v \in V$. Then by Proposition 3.3.1 there exists an invariant $I \in S(g)^{G}$ such that $\langle v, d I\rangle \neq 0$. But since (3.6.2) is defined by the action of the adjoint group one has $\left\langle\sigma_{*} v, d I\right\rangle=\langle v, d I\rangle \neq 0$. This implies $\sigma_{*} v \notin \sigma_{*}(U)$ since clearly $\sigma_{*}(U)$ is the tangent space to $Z(\gamma)$ at $\sigma\left(g_{o}, w_{o}\right)$ and $I \mid Z(\gamma)$ is a constant function. Thus not only is $\sigma_{*} \mid V$ injective but $\sigma_{*}(V) \cap \sigma_{*}(U)=0$. This implies that $\sigma_{*}$ is injective on $T$. Q.E.D.

Given an element $y \in Z$ we will later refer to $g_{0} \in H$ and $w_{0} \in h_{+}$as its ( $H \times h_{+}$) parameters in case ( $g_{o}, w_{o}$ ) corresponds to $y$ by (3.7.2).

Remark 3.7. If $y \in Z$ and $\left(g_{o}, w_{o}\right)$ are its $H \times \hbar_{+}$parameters note that

$$
\begin{equation*}
y=f+z+\sum_{i=1}^{l} h(g)^{-\alpha_{i}} e_{\alpha_{i}}, \tag{3.7.4}
\end{equation*}
$$

where $z \in \hbar$ and $g$ is given by (3.7.1). This is clear from (3.6.5) and the application of (2.6.6) to $g_{\mathrm{C}}$.

## 4. The Isospectral Leaf $Z(\gamma)$ as a Complete, Flat, Affinely Connected Manifold

4.1. In Section $4, F$ is either $\mathbb{R}$ or $\mathbb{C}$. As in Section 2.1 let $g$ be a semisimple Lie algebra over $F$ which is split if $F=\mathbb{R}$ and where $l=\operatorname{rank} g$ let $Z$ be the $2 l$-dimensional manifold of normalized Jacobi elements defined as in Section 2.2. For any invariant $I \in S(g)^{G}$ we recall that $\xi_{I}$ is the vector field on $Z$ defined so that $\left(\xi_{I}\right)_{y}=\left[y,\left(\delta_{\bar{G}} I\right)(y)\right]$ for any $y \in Z$. See (2.2.19). Also, $z_{y}$ is the subspace of the tangent space to $Z$ at $y$ whose elements are all vectors of the form $\left(\xi_{I}\right)_{y}, I \in S(g)^{G}$. See (2.2.21). Furthermore we note that by Theorem 2.2 the map $y \rightarrow k_{y}$ defines a smooth distribution (in the sense of E. Cartan) on $Z$ of dimension $l$, which we denote by $\mathscr{Z}$. Of course any smooth involutory distribution on a manifold defines a foliation of the manifold by the family of maximal integral submanifolds. We refer to this foliation as the corresponding foliation. On the other hand for the case of $Z$ one already has the disjoint union

$$
\begin{equation*}
Z=\bigcup_{\gamma \in \mathscr{\mathcal { S }}(Z)} Z(\gamma), \tag{4.1.1}
\end{equation*}
$$

where the isospectral leaves $Z(\gamma)$ by Theorems 2.5 and 3.6 are connected submanifolds of dimension $l$.

Proposition 4.1. The distribution $\mathscr{Z}$ on $Z$ is involutory and (4.1.1) is the corresponding foliation.

Proof. As in Section 2.2 let

$$
\begin{equation*}
z=\left\{\xi_{I} \mid I \in S(g)^{G}\right\} \tag{4.1.2}
\end{equation*}
$$

By Theorem 2.6, $z$ is a Lie algebra (in fact, a commutative Lie algebra). This proves that $\mathscr{Z}$ is involutory.

Since $Z(\gamma)$ is a closed connected submanifold of dimension $l$ for any $\gamma \in \mathscr{I}(Z)$ it then suffices only to prove

$$
\begin{equation*}
z_{y}=T_{y}(Z(\gamma)) \tag{4.1.3}
\end{equation*}
$$

for any $y \in Z(\gamma)$. Here of course the right side of (4.1.3) is the tangent space to $Z(\gamma)$ at $y$. However, since $\operatorname{dim} Z=2 l$ this is an immediate consequence of (1.6.8), Proposition 2.3.1, and Theorem 2.2.
Q.E.D.
4.2. Now let $\Lambda$ be any smooth manifold. If $\Lambda$ has an affine connection we can speak of covariant constant vector fields $\xi$ (either global or local) on $\Lambda$. If the affine connection is flat (i.e., the curvature and torsion vanish) then in a sufficiently small neighborhood $V$ of any point the space $a$ of such vector fields spans the tangent space at all points of $V$ and is a commutative Lie algebra. On the other hand one easily sees that if $a$ is any commutative Lie algebra of vector fields on $\Lambda$ which spans the tangent space at each point of $\Lambda$ then there exists a unique flat affine connection such that the elements of $a$ are covariant constant.

Remark 4.2. One notes in particular that if $\Lambda$ is an open subset of a Lie group $A$ whose Lie algebra is Abelian then there is a unique flat affine connection on $\Lambda$ such that any restriction $\Lambda$ of a left invariant vector field on $A$ is covariant constant.

Now if $I \in S(g)^{G}$ and $\gamma \in \mathscr{I}(Z)$ then the vector field $\xi_{1}$ is tangent to the leaf $Z(\gamma)$ by Proposition 4.1. Thus $\xi_{I} \mid Z(\gamma)$ is a vector field on $Z(\gamma)$.

Proposition 4.2. Let $\gamma \in \mathscr{I}(Z)$ be arbitrary. Then there exists a unique flat affine connection in the isospectral leaf $Z(\gamma)$ such that $\xi_{I} \mid Z(\gamma)$ is covariant constant for any invariant $I \in S(g)^{G}$.

Proof. Let $\not \approx(\gamma)=\left\{\xi_{I}\right.$ restricted to $\left.Z(\gamma) \mid I \in S(g)^{G}\right\}$. By (4.1.3) it is clear that $z(\gamma)$ spans the tangent space at all points of $Z(\gamma)$. Thus it suffices to observe that $z(\gamma)$ is a commutative Lie algebra. But if $z$ is defined by (2.2.20) then $z$ is a commutative Lie algebra by Theorem 2.2. Thus $\approx(\gamma)$ is also a commutative Lie algebra since it is a homomorphic image of $z$.
Q.E.D.

Henceforth any isospectral leaf $Z(\gamma)$ will also be regarded as having the
structure of an affinely connected manifold where the affine connection is given by Proposition 4.2.
4.3. Now let $x \in g$ and $I \in S(g)^{G}$. Recalling the definition of $G$ if $F=\mathbb{R}$ (see (3.4.4)), put $G=G$ if $F=\mathbb{C}$. Now ( $\delta I)(x) \in g^{x}$ by (1.3.3) and hence $I$ defines a (left invariant) vector field $L_{I}^{x}$ on the centralizer $G^{x}$ of $x$ in $G$ such that if $g \in \tilde{G}^{x}$ and $\psi$ is a smooth function on $\mathcal{G}^{x}$ one has

$$
\begin{equation*}
\left(L_{I}^{x} \psi\right)(g)=\left.\frac{d}{d t} \psi(g \exp t(\delta I)(x))\right|_{t=0} . \tag{4.3.1}
\end{equation*}
$$

Proposition 4.3. Assume $x \in R$ (i.e., $x$ is regular) and $A$ is any open set in $\mathcal{G}^{x}$. Then there exists a unique flat affine connection on $A$ such that $L_{1}{ }^{x} \mid A$ is covariant constant for any $I \in S(g)^{G}$.

Proof. Recall that $I_{j} \in S(g)^{G}, j=1, \ldots, l$, are the fundamental invariants. Since $x \in R$ then, as already noted in the proof of Proposition 2.1, the elements $\left(\delta I_{j}\right)(x)$ are linearly independent. But since $(\delta I)(x) \in$ cent $g^{x}$ by (1.3.3) for any $I \in S(g)^{G}$ and $\operatorname{dim} g^{x}=l$ one recovers not only the well-known fact that $g^{x}$ is Abelian but also the fact that any element in $g^{x}$ is uniquely of the form $(\delta I)(x)$ for some $I \in S(g)^{G}$. Thus any left (or equivalently right) invariant vector field on $G^{x}$ is of the form $L_{I}^{x}$ for some $I \in S(g)^{G}$. The result then follows immediately as noted in Remark 4.2.
Q.E.D.

We now show that the isomorphisms of Theorems 2.4, 2.6, 3.2, and 3.6 are also isomorphisms of flat affinely connected manifolds.

Theorem 4.3. Let $g$ be a semi-simple Lie algebra over $F=\mathbb{R}$ or $\mathbb{C}$ which is split if $F=\mathbb{R}$. Let $l=\operatorname{rank}^{g}$ and let $\mathscr{I}: g \rightarrow \mathbb{F}^{l}$ be the map defined as (2.3.1) by the fundamental invariant polynomials on $g$. Let $h$ be as in Section 2.1 so that $h$ is a Cartan subalgebra of $g$. (If $F=\mathbb{R}$ then $h$ is split and $h_{+}$is an open Weyl chamber in $h$ defined as in (3.3.1).)

Now let $w_{o} \in h$, where $w_{o} \in h_{+}$if $F=\mathbb{R}$, and put $w=f+w_{o}$, where $f$ is defined by (1.5.4). Let $\gamma=\mathscr{I}\left(w_{o}\right)$ and let $Z(\gamma)$ be the corresponding isospectral leaf of normalized Jacobi elements defined as in Section 2.3. Let $Z(\gamma)$ have the structure of a flat affinely connected manifold as in Proposition 4.2. Put $G_{o}{ }^{w}=G_{(*)}^{w}$ if $F=\mathbb{C}$ and $G_{o}{ }^{w}=G_{0}{ }^{w}$ if $F=\mathbb{P}$ using the notation of (3.5.7) so that (since $w \in R$ and $G_{o}{ }^{w}$ is open in $\widehat{G}^{w}$ ) $G_{o}{ }^{w}$ has, as given by Remark 4.2, the structure of a flat affinely connected manifold. Let $\beta_{(w)}: G_{0}{ }^{w} \rightarrow Z(\gamma)$ be the map defined as in (2.6.15) and (3.6.2) so that if $a \in G_{o}{ }^{w}$ and $n(a)$ is defined by (2.6.2) and (3.4.14) then $\beta_{(w)}(a)=n(a) w$. Let $g \in G_{o}{ }^{w}$ and let $y=\beta_{(w)}(g)$. Put $G_{o}{ }^{y}=G_{*}^{y}$ if $F=\mathbb{C}$ and $G_{o}{ }^{y}=G_{e}{ }^{y}$ if $F=\mathbb{R}$ using the notation of (2.4.9) and (3.1.23) so that (since $y \in R$ and $G_{o}{ }^{\nu}$ is open in $\mathcal{G}^{v}$ ) $G_{o}{ }^{v}$ also has, as given by Remark 4.2, the structure of a flat affinely connected manifold. Now for any $a \in G_{o}{ }^{w}$ let $\psi_{s}(a)=$
$n(g) g^{-1} a n(g)^{-1}$. Then $\psi_{g}\left(G_{o}{ }^{u r}\right) \subseteq G_{o}{ }^{y}$. Finally let $\beta_{y}: G_{o}{ }^{y} \rightarrow Z(\gamma)$ be the map given by (2.4.11) and (3.2.15) so that if $d \in G_{o}{ }^{y}$ and $n_{d}$ is defined by (2.4.6) and (3.2.12) then $\beta_{y}(d)=n_{d} y$. Then one has a commutative diagram


Furthermore not only are the three maps in (4.3.2) diffeomorphisms but they are isomorphisms of flat affinely connected manifolds. In fact for any invariant $I \in S(g)^{G}$ one has for the corresponding mappings of vector fields, $\psi_{g}\left(L_{I}{ }^{w} \mid G_{o}{ }^{w}\right)=L_{I}{ }^{y} \mid G_{o}{ }^{y}$ and

$$
\begin{equation*}
\beta_{(w)}\left(L_{I}^{w} \mid G_{o}^{w}\right)=\beta_{y}\left(L_{l}{ }^{y} \mid G_{o}{ }^{y}\right)=\xi_{l} \mid Z(\gamma), \tag{4.3.3}
\end{equation*}
$$

where the vector fields $L_{I}{ }^{w}, L_{I}{ }^{y}$, and $\xi_{I}$ are defined by (4.3.1) and (2.2.19).
Proof. The fact that $w$ and $y$ are in $R$, as noted previously, is a consequence of (2.1.7) since $w, y \in b_{f}$. Now assume first that $F=\mathbb{C}$. Then $G_{(*)}^{w}$ and $G_{*}^{y}$ are open respectively in $\tilde{G}^{w}$ and $\tilde{G}^{y}$ since in fact they are Zariski open. One has $\psi_{g} G_{(*)}^{w} \subseteq G_{*}^{v}$ by (2.6.11). The fact that the diagram (4.3.2) is commutative and all three maps are diffeomorphisms (in fact they are algebraic isomorphisms) was established in the proof of Theorem 2.6. See (2.6.21).

Now let $I \in S(g)^{G}$. Then by (1.3.1) one has $n(g)(\delta I)(w)=(\delta I)(n(g) w)=$ $(\delta I)(y)$. Thus since $n(g) \tilde{G}^{w} n(g)^{-1}=\tilde{G}^{y}$, conjugation by $n(g)$ carries $L_{I}{ }^{w}$ to $L_{l}{ }^{y}$. On the other hand $L_{I}{ }^{w}$ is fixed by the left translation of $\tilde{G}^{w}$ by $g^{-1}$. Thus $\psi_{g}\left(L_{I}{ }^{w} \mid G_{o}{ }^{w}\right)=L_{I}{ }^{y} \mid G_{o}{ }^{y}$. Hence of course by the commutativity of (4.3.2) one has $\beta_{(w)}\left(L_{I} \mid G_{o}{ }^{w}\right)=\beta_{y}\left(L_{I}{ }^{y} \mid G_{o}{ }^{y}\right)$. It suffices only to show (for the case where $F=\mathbb{C}$ ) that the vector field is $\xi_{I} \mid Z(\gamma)$. But now since $g$ is an arbitrary element in $G_{o}{ }^{w}$ it thus suffices to show

$$
\begin{equation*}
\beta_{(w)}\left(\left(L_{l}^{w}\right)_{g}\right)=\left(\xi_{I}\right)_{y} \tag{4.3.4}
\end{equation*}
$$

Now let $g(t)=\exp t(\delta I)(w)$ for $t \in \mathbb{R}$. Put $a(t)=g g(t)$ so that $\left(L_{I}{ }^{w}\right)_{g}$ is the tangent vector to $a(t)$ at $t=0$. Let $d(t)=\psi_{g}(a(t))=n(g) g(t) n(g)^{-1}$. Thus

$$
\begin{equation*}
d(t)=\exp t((\delta I)(y)) \tag{4.3.5}
\end{equation*}
$$

since $n(g) w=y$. But if $v$ is the left side of (4.3.4) then from the commutativity of (4.3.2) it follows that $v$ is the tangent vector to $\beta_{(w)}(a(t))=\beta_{y}(d(t))$ at $t=0$. Now recalling (2.4.6) put $\bar{n}_{d(t)} h_{d(t)}=b_{d(t)}$ so that

$$
\begin{equation*}
d(t)=b_{d(t)} n_{d(t)} \tag{4.3.6}
\end{equation*}
$$

where we note that $\bar{b}_{\boldsymbol{a}(t)} \in \bar{B}$. One notes also that $\bar{b}_{a(0)}=n_{d(0)}=1$. Now let $x \in n$ be the tangent vector to $n_{d(t)}$ at $t=0$. Then since $\beta_{y}(d(t))=n_{d(t)} y$ it follows that

$$
\begin{equation*}
v=[x, y] . \tag{4.3.7}
\end{equation*}
$$

On the other hand if $z \in \bar{b}$ is the tangent vector to $\bar{b}_{d(t)}$ at $t=0$ then by (4.3.5) and (4.3.6) one has

$$
\begin{equation*}
(\delta I)(y)=z+x . \tag{4.3.8}
\end{equation*}
$$

But then recalling (1.2.9) one must have $x-\left(\delta_{n} I\right)(y)$ so that $v-\left[\left(\delta_{n} I\right)(y), y\right]$. But then $v=\left[y,\left(\delta_{\bar{\sigma}} I\right)(y)\right]$ by (1.3.3) so that $v=\left(\xi_{I}\right)_{y}$ by the definition of $\xi_{I}$. This proves (4.3.4) and hence the theorem is proved for $F=\mathbb{C}$.

Now assume $F=\mathbb{R}$. Then $G_{o}{ }^{w}$ and $G_{o}{ }^{y}$ are, by definition, the connected components $G_{0}{ }^{w}$ and $G_{e}{ }^{y}$ (recall that $G_{0}{ }^{w}$ is not the identity component), respectively, of $G^{w}$ and $\tilde{G}^{y}$. In particular they are, respectively, open subsets of $\vec{G}^{w}$ and $\bar{G}^{y}$.

Now the fact that the diagram (4.3.2) is commutative and that the three maps in question are diffeomorphisms has been established in the proof of Theorem 3.6. See Remark 3.6 and (3.6.7). In fact these statements were proved by first applying Theorem 2.6 (and its proof) to $g \mathbb{c}$ and then by restricting our considerations of

$$
\begin{equation*}
\left(G_{\mathbb{C}^{w}}\right)_{(*)}, \quad\left(G_{\mathbb{C}}\right)_{*}, \quad \text { and } \quad Z_{\mathbb{C}}(\gamma) \tag{4.3.9}
\end{equation*}
$$

to the respective real submanifolds

$$
\begin{equation*}
G_{0}{ }^{w}, \quad G_{e}{ }^{y}, \quad \text { and } \quad Z(\gamma) . \tag{4.3.10}
\end{equation*}
$$

But now again from what we have proved above in the complex case one has (4.3.3), where the vector fields in question are, in the obvious order, defined on the manifolds in (4.3.9). However, they are clearly tangent, in the same order, to the real submanifolds appearing in (4.3.10). Hence one has (4.3.3) for the case $F=\mathbb{R}$.
Q.E.D.

Remark 4.3. We point out here that at least in one way there is a significant difference between the real and the complex cases of Theorem 4.3. Namely, it is only in the real case that the manifolds in question, i.e., in (4.3.10), are complete with respect to their flat affine connections. This completeness will be more meaningful when we deal with the integration of Hamilton's equations in Section 7.

## 5. Representations and the Function $\Phi_{\lambda}\left(g_{0}, w_{0} ; t\right)$

5.1. Henceforth, that is, for the remainder of the paper, we assume $F=\mathbb{R}$ so that $g$ is a real split semi-simple Lie algebra and all of our previous results and notations for the case $F=\mathbb{R}$ apply.

Let $w_{o} \in h_{+}$and as usual put $w=f+w_{o}$. Let $a \in G_{0}{ }^{w}$ so that $h(a) \in H$ (see (3.4.14)). Now recalling (2.1.1) one has $h^{\prime}=h(H)$ so that $h(a)^{\lambda}$ is defined for any $\lambda \in h^{\prime}$. Our main results (see Theorem 7.5) depend upon a certain formula for $h(a)^{\lambda}$. Obtaining this formula will be the main objective of Section 5. The formula will arise from the finite-dimensional representation theory of $g$-which we now consider.

Now the restriction $Q \mid h_{c}$ induces a nonsingular bilinear form-also denoted by $Q$-on the dual space $\ell_{\mathbb{C}}^{\prime}$. We may of course regard $\ell^{\prime}$ as a real form of $\ell_{\mathbb{C}}^{\prime}$ and note that $Q$ is positive definite on $\ell^{\prime}$. Now let $\Lambda$ be the lattice in $\hbar^{\prime}$ defined by

$$
\begin{equation*}
\Lambda=\left\{\lambda \in \mathscr{h}^{\prime} \left\lvert\, \frac{2 Q\left(\lambda, \alpha_{i}\right)}{Q\left(\alpha_{i}, \alpha_{i}\right)} \in \mathbb{Z}\right., i=1, \ldots, l\right\} \tag{5.1.1}
\end{equation*}
$$

Now if $\nu_{j} \in \Lambda, j=1, \ldots, l$, are defined by the relation $2 Q\left(\nu_{j}, \alpha_{i}\right) / Q\left(\alpha_{i}, \alpha_{i}\right)=\delta_{i j}$ then of course one has the direct sum

$$
\begin{equation*}
\Lambda=\sum_{i=1}^{l} \mathbb{Z} \nu_{i} \tag{5.1.2}
\end{equation*}
$$

Now let $G_{\mathbb{C}}{ }^{8}$ be a fixed simply connected Lie group having gC as its Lie algebra. The adjoint representation defines a homomorphism Ad: $G_{\mathbb{C}}{ }^{8} \rightarrow G_{\mathbb{C}}$ and in fact it defines an exact sequence

$$
\begin{equation*}
(1) \rightarrow \operatorname{cent} G_{\mathbb{C}}^{s} \rightarrow G_{\mathbb{C}}^{s} \xrightarrow{\mathrm{Ad}} G_{\mathbb{C}} \rightarrow(1) \tag{5.1.3}
\end{equation*}
$$

If $x \in g_{C}$ then to avoid confusion we will write $\exp ^{s} x$ for its exponential image in $G_{\mathbb{C}^{8}}$ and retain the previous notation $\exp x$ for its exponential image in the adjoint group $G_{\mathbb{C}}$. Also let $H_{\mathbb{C}}{ }^{8}, N_{\mathbb{C}^{s}}$, and $\bar{N}_{\mathbb{C}}{ }^{8}$ be the subgroups of $G_{\mathbb{C}}{ }^{8}$ corresponding respectively to $h_{\mathbb{C}}, n_{\mathrm{C}}$, and $\bar{n}_{\mathbb{C}}$.

Now if $h_{\mathbb{C}}^{\prime}\left(H_{\mathbb{C}^{8}}\right) \subseteq \hbar_{\mathbb{C}}^{\prime}$ is defined by (2.1.1), where $\boldsymbol{h}_{\mathbb{C}}^{\prime}$ replaces $h^{\prime}, h_{\mathbb{C}}$ replaces $h_{\text {, }}$ and $\exp ^{8}$ replaces exp, one knows that

$$
\begin{equation*}
\Lambda=h_{\mathscr{C}}^{\prime}\left(H_{\mathbb{C}}^{\varepsilon}\right) \tag{5.1.4}
\end{equation*}
$$

so that $h^{\nu}$ is defined for any $\nu \in \Lambda, h \in H \mathbb{C}^{g}$. The elements of $\Lambda$ are often then referred to as the integral linear forms (on $H_{\mathbb{C}}{ }^{g}$ ). The cone $D \subseteq \Lambda$ of dominant integral linear forms is defined by

$$
\begin{equation*}
D=\left\{\lambda \in \Lambda \mid Q\left(\lambda, \alpha_{i}\right) \geqslant 0, i=1, \ldots, l\right\} \tag{5.1.5}
\end{equation*}
$$

If $\mathbb{Z}_{+}$is the set of nonnegative integers one clearly has

$$
\begin{equation*}
D=\sum_{i=1}^{l} \mathbb{Z}_{+} \nu_{i} \tag{5.1.6}
\end{equation*}
$$

Now a complex vector space $V_{\mathbb{C}}$ will be called a $G_{\mathbb{C}}{ }^{8}$ module if it is understood that there is a representation $\pi: G_{\mathbb{C}^{g}} \rightarrow$ Aut $V_{\mathbb{C}}$ of $G_{\mathbb{C}^{s}}$ on $V_{\mathbb{C}}$. In such a case we will write $g v \in V_{\mathbb{C}}$ for $\pi(g) v$, where $g \in G_{\mathbb{C}}{ }^{s}, v \in V_{\mathbb{C}}$. Assume $V_{\mathbb{C}}$ is a $G_{\mathbb{C}}{ }^{s}$ module. A vector $v \in V_{\mathbb{C}}$ is called a weight vector for $\mu \in \Lambda$ if $h v=h^{\mu} v$ for all $h \in H_{\mathbb{C}^{s}}$. For any $\mu \in \Lambda$ let $V_{\mathbb{C}}(\mu) \subseteq V_{\mathbb{C}}$ be the subspace of all weight vectors for $\mu$. The element $\mu \in \Lambda$ is called a weight of $V_{\mathbb{C}}($ or $\pi)$ if $V_{\mathbb{C}}^{-}(\mu) \neq 0$. Now if $\mu$ is a weight of $V$ then $0 \neq v \in V(\mu)$ is called a highest (resp. lowest) weight vector if $n v=v$ (resp. $\bar{n} v=v$ ) for all $n \in N_{\mathbb{C}}{ }^{s}$ (resp. $\bar{n} \in \bar{N}_{\mathbb{C}}{ }^{s}$ ). If $V_{\mathbb{C}}(\mu)$ contains a highest (resp. lowest) weight vector then $\mu$ is called a highest (resp. lowest) weight.

We now recall certain fundamental results of the Cartan-Weyl theory of representations. For each $\lambda \in D$ there exists an irreducible finite-dimensional (holomorphic) representation

$$
\begin{equation*}
\pi_{\lambda}: G_{\mathbb{C}}{ }^{s} \rightarrow \text { Aut } V_{\mathbb{C}}{ }^{\lambda} \tag{5.1.7}
\end{equation*}
$$

such that $\lambda$ is a highest weight. Furthermore (1) as such $V_{\mathbb{C}}{ }^{\lambda}$ is unique up to equivalence, (2) $\lambda$ is the only highest weight of $V_{\mathbb{C}^{\lambda}}$, and (3) $\operatorname{dim} V_{\mathbb{C}^{\lambda}}(\lambda)=1$.

Now for any $\lambda \in D$ let (5.1.7) be given and fixed. One knows then that if $\pi$ is any finite-dimensional irreducible holomorphic representation of ${\dot{G_{\mathbb{C}}}}^{s}$ then $\pi$ is equivalent to $\pi_{\lambda}$ for a unique $\lambda \in D$.

Now let $U=U(g)$ be the universal enveloping algebra of $g$ over $\mathbb{R}$ and let $U_{\mathbf{C}}$, its complexification, be the universal enveloping algebra of $g_{\mathbb{C}}$ over $\mathbb{C}$. If the adjoint representation is extended the group $G_{\mathbb{C}^{s}}$ operates as a group of automorphisms of $U_{\mathbb{C}}$ and one has the direct sum

$$
\begin{equation*}
U_{\mathbb{C}}=\sum_{\mu \in \Lambda} U_{\mathbb{C}}(\mu) \tag{5.1.8}
\end{equation*}
$$

where as above the $U_{\mathbb{C}}(\mu)$ are the weight spaces for this action. On the other hand given $\lambda \in D$ the representation $\pi_{\lambda}$ induces a representation of $g \mathrm{C}$ and hence of $U_{\mathbb{C}}$ on $V_{\mathbb{C}}^{\lambda}$, which we also denote by $\pi_{\lambda}$, so that $V_{\mathbb{C}^{\lambda}}$ becomes a $U_{\mathbb{C}}$ module and one has

$$
\begin{equation*}
U_{\mathbb{C}}(\mu) V_{\mathbb{C}^{\lambda}}(\nu) \subseteq V_{\mathbb{C}^{\lambda}}(\mu+\nu) \tag{5.1.9}
\end{equation*}
$$

for any $\mu, \nu \in \Lambda$.
Now let $G^{s}$ be the subgroup of $G_{\mathbb{C}^{s}}$ corresponding to $\mathscr{g}$ and let $N^{s}, H^{s}$ and
$\bar{N}^{s}$ be the subgroups of $G^{s}$ corresponding to $n, h$, and $\bar{n}$. Note that since $h^{u}$ is real (in fact positive) for any $h \in H^{s}$ (or $H$ ) and $\mu \in \Lambda$ one has the direct sum

$$
\begin{equation*}
U=\sum U(\mu) \tag{5.1.10}
\end{equation*}
$$

over $\Lambda$, where $U(\mu)=U \cap U_{\mathbb{C}}(\mu)$.
Now let $G_{u}{ }^{s}$ be the maximal compact subgroup of $G_{\mathbb{C}}{ }^{s}$ corresponding to the compact form $g_{u}$. See (3.1.6). One knows then (the "unitarian trick") that there exists a Hermitian (positive definite) inner product on $V_{\mathbb{C}}{ }^{\lambda}$ which is invariant under the action of $G_{u}{ }^{s}$. We will denote the inner product of $v$, $v^{\prime} \in V_{\mathbb{C}}{ }^{\lambda}$ by $\left\{v, v^{\prime}\right\}$. It is chosen so that it is linear in $v$ and conjugate linear in $v^{\prime}$.

Now since $G_{\mathbb{C}}{ }^{s}$ is simply connected the $*$-operation on $\mathscr{g}_{\mathbb{C}}$ (see Section 3.1) induces a unique $*$-operation $a \mapsto a^{*}$ on $G_{\mathbb{E}^{s}}$ such that one has (3.1.8) for $a, b \in G_{\mathbb{C}^{s}}$ and (3.1.9) for $x \in g \mathbb{C}$, where $\exp ^{s}$ replaces exp. It is clear that the *-operations on $G_{\mathbb{C}}{ }^{s}$ and $G_{\mathbb{C}}$ commute with Ad. It is also clear that the *-operation on $g \mathbb{C}$ extends uniquely to $U_{\mathrm{C}}$ as a conjugate linear map such that $\left(u_{1} u_{2}\right)^{*}=$ $u_{2}^{*} u_{1}^{*}$ for $u_{1}, u_{2} \in U_{\mathbb{C}}$. One easily has

$$
\begin{equation*}
\left\{u v, v^{\prime}\right\}=\left\{v, u^{*} v^{\prime}\right\} \tag{5.1.11}
\end{equation*}
$$

for any $v, v^{\prime} \in V_{\mathbb{C}^{\lambda}}$ and $u$ in $U_{\mathbb{C}}$ or $G_{\mathbb{C}}{ }^{s}$.
Now fix once and for all a highest weight vector $v^{\lambda} \in V_{\mathscr{C}^{\lambda}}(\lambda)$ such that $\left\{v^{\lambda}, v^{\lambda}\right\}=1$ and let $V^{\lambda}$ be the $\mathbb{R}$-subspace of $V \mathbb{C}^{\lambda}$ defined by

$$
\begin{equation*}
V^{\lambda}=U v^{\lambda} . \tag{5.1.12}
\end{equation*}
$$

Also for any $\mu \in \Lambda$ let $V^{\lambda}(\mu)=V^{\lambda} \cap V_{\mathbb{C}}^{-\lambda}(\mu)$.
Proposition 5.1. The $\mathbb{R}$-subspace $V^{\lambda}$ of $V_{\mathbb{C}}{ }^{\lambda}$ is a real form of $V_{\mathbb{C}}{ }^{\lambda}$. That is, $V_{\mathbb{C}^{\lambda}}=V^{\lambda}+i V^{\lambda}$ is a real direct sum. Furthermore if $Q_{*}^{\lambda}$ is the restriction of the inner product $\left\{v, v^{\prime}\right\}$ to $V^{\lambda}$ then $Q_{*}^{\lambda}$ is real valued (so that $V^{\lambda}$ is a real Hilbert space). Furthermore $V^{\lambda}$ is stable under the action of $G^{s}$ and

$$
\begin{equation*}
V^{\lambda}=\sum V^{\lambda}(\mu) \tag{5.1.13}
\end{equation*}
$$

summed over the weights of $\pi_{\lambda}$, is an orthogonal direct sum with respect to $Q_{*}^{\lambda}$.
Proof. Now $V_{\mathbb{C}^{\lambda}}=U_{\mathbb{C}} v^{\lambda}$ since of course $V_{\mathbb{C}}{ }^{\lambda}$ is $U_{\mathbb{C}}$ irreducible. Thus $V \mathbb{C}^{\lambda}=V^{\lambda}+i V^{\lambda}$. To show $V^{\lambda}$ is a real form it suffices to show $V^{\lambda} \cap i V^{\lambda}=0$. For this it clearly is enough to show that $Q_{*}^{\lambda}$ is real valued. Let $T^{s}$ be the subgroup (maximal torus) of $G_{u}{ }^{s}$ corresponding to $i \hbar$. Now if $\mu_{i} \in \Lambda, i=1,2$, are distinct then the characters they define on $T^{s}$ are distinct and hence one has

$$
\begin{equation*}
\left\{V_{\mathbb{C}^{\lambda}}\left(\mu_{1}\right), V_{\mathbb{C}^{\lambda}}\left(\mu_{2}\right)\right\}=0 \tag{5.1.14}
\end{equation*}
$$

since $T^{s}$ operates as a group of unitary operators on $V_{\mathbb{C}}$. Now let $v_{i} \in V^{\lambda}$, $i=1,2$. Write $v_{i}=u_{i} v^{\lambda}$ for $u_{i} \in U$. Put $u=u_{2}^{*} u_{1}$. Then $u \in U$ since $U$ is clearly stable under the $*$-operation. But now $\left\{v_{1}, v_{2}\right\}=\left\{u_{1} v^{\lambda}, u_{2} v^{\lambda}\right\}=$ $\left\{u_{2}^{*} u_{1} v^{\lambda}, v^{\lambda}\right\}=\left\{u v^{\lambda}, v^{\lambda}\right\}$. To show $Q_{*}^{\lambda}$ is real valued it suffices then to show that $\left\{u v^{\lambda}, v^{\lambda}\right\} \in \mathbb{R}$. Let $u_{0}$ be the component of $u$ in $U(0)(0$ here denotes the zero weight) relative to (5.1.10). By (5.1.9) and (5.1.14) one then has $\left\{u v^{\lambda}, v^{\lambda}\right\}=$ $\left\{u_{0} v^{\lambda}, v^{\lambda}\right\}$. However, if $U(h)$ is the enveloping algebra of $h$ over $\mathbb{R}$ then one knows

$$
\begin{equation*}
U(0)=U(f) \oplus((U n) \cap U(0)) \tag{5.1.15}
\end{equation*}
$$

is a direct sum, where $U_{n}$ is the left ideal in $U$ generated by $n$. In fact the validity of (5.1.15) follows from the validity of (5.1.15) for the complexification of the subspace involved. For the case over $\mathbb{C}$ see, e.g., Lemma 7.4.2 in [4, p. 230]. (In that case one depends on (5.1.15) to define the Harish-Chandra homomorphism.) But since $v^{\lambda}$ is a highest weight vector one has

$$
\begin{equation*}
x v^{\lambda}=0 \quad \text { for all } x \in n_{\mathbb{C}} \tag{5.1.16}
\end{equation*}
$$

and hence if $u_{1}$ is the component of $u_{0}$ in $U(\ell)$ relative to (5.1.15) one has $\left\{u_{0} v^{\lambda}, v^{\lambda}\right\}=\left\{u_{1} v^{\lambda}, v^{\lambda}\right\}$. However, if $y \in \mathscr{A}$ then $y v^{\lambda}=\langle\lambda, y\rangle v^{\lambda}$. But $\langle\lambda, y\rangle \in \mathbb{R}$ since $\lambda \in h^{\prime}$. Thus $u_{1} v^{\lambda}=r v^{\lambda}$ for some $r \in \mathbb{R}$ so that $\left\{u_{1} v^{\lambda}, v^{\lambda}\right\} \in \mathbb{R}$. This proves that $Q_{*}^{\lambda}$ is real valued and $V^{\lambda}$ is a real form of $V_{\mathbb{C}^{\lambda}}$. But now $U(\mu-\lambda) v^{\lambda} \subseteq V^{\lambda}(\mu)$ by (5.1.9) and hence one has the sum (5.1.13) by (5.1.10). The sum is orthogonal by (5.1.14). The subspace $V^{\lambda}$ is stable under $G^{s}$ since it is clearly stable under $g$.
Q.E.D.

Remark 5.1. Note that $V^{\lambda}(\mu)$, for any $\mu \in \Lambda$ is a real form of the weight space $V_{\mathbb{C}^{\lambda}}(\mu)$. This follows obviously from (5.1.13) and (5.1.14).

We will refer to $V^{\lambda}$ with respect to the action of $G^{s}$ by $\pi_{\lambda} \mid G^{s}$ as a $G^{s}$-module.
5.2. Now let $G_{*}^{s}=\bar{N}^{s} H^{s} N^{s}$ so that $G_{*}^{s}$, using the Bruhat decomposition of $G_{\mathbb{C}^{8}}$, is an open connected subset of $G^{s}$ and the map

$$
\begin{equation*}
\operatorname{Ad}: G_{*}^{s} \rightarrow G_{*} \tag{5.2.1}
\end{equation*}
$$

recalling (3.2.1) is a diffeomorphism.
Regarding Ad, as in (5.1.3), as a map from $G_{\mathbb{C}}{ }^{s}$ to $G_{\mathbb{C}}$
Lemma 5.2.1. One has

$$
\left(\operatorname{Ad}^{-1} G_{*}\right) \cap G^{s}=\left(\operatorname{cent} G^{s}\right) G_{*}^{s} .
$$

Furthermore if $c, c^{\prime} \in \operatorname{cent} G^{s}$ are distinct then $c G_{*}^{s}$ and $c^{\prime} G_{*}^{s}$ are disjoint so that
the connected components of $\left(\operatorname{cent} G^{s}\right) G_{*}^{s}$ are uniquely of the form $c G_{*}^{s}$ for $c \in \operatorname{cent} G^{s}$.

Proof. The first statement is an immediate consequence of (5.1.3) and the surjectivity of (5.2.1). Now if $a, a^{\prime} \in G_{s}^{*}$ and $a c=a^{\prime} c^{\prime}$, where $c, c^{\prime} \in$ cent $G^{s}$, then $\operatorname{Ad} a=\operatorname{Ad} a c=\operatorname{Ad} a^{\prime} c^{\prime}=\operatorname{Ad} a^{\prime}$. But then $a=a^{\prime}$ from the injectivity of (5.2.1) so that $c=c^{\prime}$. This proves the lemma.
Q.E.D.

Now let $w_{o} \in h_{+}$and as usual let $w=f+w_{o}$. We recall that $\bar{n}_{f}(w) \in \bar{N}$ satisfies $\bar{n}_{f}(w) w_{o}=w$. See (3.5.2). Now let $m \in \tilde{G}$ be as in (3.5.5) so that $m e_{\alpha_{i}}=-e_{\alpha_{i}}, i=1, \ldots, l$. Hence of course $m e_{-\alpha_{i}}=-e_{-\alpha_{i}}$ so that $m f=-f$. But clearly $m w_{o}=w_{o}$. Thus if we put

$$
\begin{equation*}
\bar{n}_{-f}(w)=m^{-1} \bar{n}_{f}(w) m \tag{5.2.2}
\end{equation*}
$$

then $\bar{n}_{-f}(w) \in \bar{N}$ and $\bar{n}_{-f}(w) w_{o}=-f+w_{o}$.
Identification 5.2. Henceforth to simplify notation we will identify $G_{*}^{s}$ with $G_{*}$ by the diffeomorphism (5.2.1). This means that $\bar{N}^{s}$ and $N, H^{s}$ and $H$ and also $N^{s}$ and $N$ are identified. The only possible confusion that can arise is with regard to multiplication of these elements. However, it should be clear from the context whether the multiplication is in $G_{\mathbb{C}}{ }^{s}$ or $\boldsymbol{G}_{\mathbb{C}}$.

Remark 5.2. One notes that both $G_{s}^{*}$ and $G_{*}$ are stable under the respective *-operations in $G_{\mathbb{C}^{8}}$ and $G_{\mathbb{C}}$ and that since Ad commutes with these operations no ambiguity with the $*$-operation is introduced by identifying $G_{*}^{s}$ with $G_{*}$.

Let $s(\kappa) \in \tilde{G}$ be as in (3.4.12).

Lemma 5.2.2. One has $m^{-1} s(\kappa) \in G$. Furthermore there exists a unique element $s_{0}(\kappa) \in G^{s}$ such that (1) $\operatorname{Ad} s_{0}(\kappa)=m^{-1} s_{0}(\kappa)$ and such that (2) for any $h \in H$ and $w_{o} \in h_{+}$one has the relation

$$
\begin{equation*}
s_{0}(\kappa)^{-1} \bar{n}_{-f}(w) h\left(\bar{n}_{f}(w)\right)^{-1} \in G_{*} \tag{5.2.3}
\end{equation*}
$$

in $G^{s}$.
Proof. Let $w_{o} \in h_{+}$and let $h \in H$. Then by (3.5.7) if $g \in G$ is defined by $g=\bar{n}_{f}(w) m h \bar{n}_{f}(w)^{-1}$ one has $g \in G_{0}{ }^{w}$. Thus $s(\kappa)^{-1} g \in G_{*}$ by Proposition 3.5 (recalling 3.4.13). But $s(\kappa)^{-1} g=s(\kappa)^{-1} \bar{n}_{f}(w) m h \bar{n}_{f}(w)^{-1}=s(\kappa)^{-1} m \bar{n}_{-f}(w) h \bar{n}_{f}(w)^{-1}$. Thus in $\tilde{G}$ one has

$$
\begin{equation*}
s(\kappa)^{-1} m \bar{n}_{-f}(w) h \bar{n}_{f}(w)^{-1} \in G_{*} \tag{5.2.4}
\end{equation*}
$$

However, since $G_{*} \subseteq G$ and of course $\bar{n}_{-f}(w) h \bar{n}_{f}(w)^{-1} \in G$ one necessarily
has $s(\kappa)^{-1} m \in G$. This proves the first statement. Now fix $s_{1}(\kappa) \in G^{s}$ so that $\operatorname{Ad} s_{1}(\kappa)=m^{-1} s(\kappa)$. Thus for any $c \in \operatorname{Cent} G^{s}$ one has

$$
\begin{equation*}
\operatorname{Ad}\left(\left(c s_{1}(\kappa)\right)^{-1} \bar{n}_{-f}(w) h \bar{n}_{f}(w)^{-1}\right) \in G_{*}, \tag{5.2.5}
\end{equation*}
$$

where the multiplication in (5.2.5) is in $G^{s}$. But then by Lemma 5.2.1 there exists a unique $c_{1} \in$ Cent $G^{s}$ such that if $s_{o}(\kappa)=c_{1} s_{1}(\kappa)$ then $s_{o}(\kappa)^{-1} \bar{n}_{-f}(w) h \bar{n}_{f}(w)^{-1} \in G_{*}^{s}=G_{*}$.

Now consider the map

$$
\begin{equation*}
\hbar_{+} \times H \rightarrow G^{s} \tag{5.2.6}
\end{equation*}
$$

where $\left(w_{0}^{\prime}, h^{\prime}\right) \mapsto s_{0}(\kappa)^{-1} \bar{n}_{-f}\left(w^{\prime}\right) h^{\prime} \bar{n}_{f}\left(w^{\prime}\right)^{-1}$ and $w^{\prime}=f+w_{0}^{\prime}$. Now by (5.2.4), where $w^{\prime}$ replaces $w$ and $h^{\prime}$ replaces $h$, and by Lemma 5.2.1 it follows that the image of (5.2.6) is in (Cent $G^{s}$ ) $G_{*}^{s}$. However, by the connectivity of $h_{+} \times H$ the image of (5.2.6) must lie in one component of (Cent $G^{s}$ ) $G_{*}^{s}$. Since ( $w_{o}, h$ ) maps into $G_{*}^{s}$ the entire image of (5.2.6) is in $G_{*}^{s}=G_{*}$. Except for the uniqueness of $s_{0}(\kappa)$ this proves the lemma. The ambiguity of $s_{0}(\kappa)$ satisfying (1) of the lemma is up to an element in Cent $G^{s}$. This uniqueness then follows immediately from Lemma 5.2.1.
Q.E.D.

Henceforth $s_{o}(\kappa)$ will denote the element of $G^{s}$ given by Lemma 5.2.2. Let $K^{s}$ be the (maximal compact) subgroup of $G^{s}$ corresponding to $k \subseteq g$.

Lemma 5.2.3. One has $s_{0}(\kappa) \in K^{s}$.
Proof. Now for $i=1, \ldots, l$ one has $s(\kappa) e_{\alpha_{i}}=e_{\kappa \alpha_{i}}$. However, since $\kappa^{2}=1$ in $W$ one then has $s(\kappa)^{2} e_{\alpha_{i}}=r_{i} e_{\alpha_{i}}$ for some $r_{i} \in \mathbb{R}$. Thus since $Q\left(e_{\alpha_{i}}, e_{-\alpha_{i}}\right)=1$ one has $r_{i}=Q\left(e_{\alpha_{i}}, s(\kappa)^{2} e_{\alpha_{i}}\right)=Q\left(s(\kappa)^{-1} e_{-\alpha_{i}}, s(\kappa) e_{\alpha_{i}}\right)=Q\left(s(\kappa)^{-1} e_{-\alpha_{i}}, e_{\kappa \alpha_{i}}\right)$ by the invariance of $Q$. But $-\kappa \alpha_{i}$ is also a simple positive root and hence $s(\kappa) e_{-\kappa \alpha_{i}}=e_{-\alpha_{i}}$. Thus $s(\kappa)^{-1} e_{-\alpha_{i}}=e_{-\kappa \alpha_{i}}$ so that $r_{i}=Q\left(e_{-\kappa \alpha_{i}}, e_{\kappa \alpha_{i}}\right)=1$. This proves

$$
\begin{equation*}
s(\kappa)^{2}=1 \tag{5.2.7}
\end{equation*}
$$

However, $m$ operates as -1 on $e_{\alpha_{i}}$ and $e_{-\alpha_{i}}$ so that clearly $s(\kappa)$ and $m$ commute and $m^{2}=1$. Hence $s_{1}(\kappa)^{2}=1$, where $s_{1}(\kappa)=m^{-1} s(\kappa)$. On the other hand for any $g \in G_{\mathbb{C}}$ and $x \in \mathscr{g} \mathbb{C}$ one easily has by (3.1.8) and (3.1.9),

$$
\begin{equation*}
\left(g^{*}\right)^{-1} x^{*}=(g x)^{*} \tag{5.2.8}
\end{equation*}
$$

Thus $\left(s(\kappa)^{-1}\right)^{*} e_{-\alpha_{i}}=e_{-\kappa \alpha_{i}}$ by (5.2.8) and (1.5.1). However, $s(\kappa)^{-1} e_{-\alpha_{i}}=e_{-\kappa \alpha_{i}}$. Similarly $s(\kappa)^{-1}(=s(\kappa))$ and $\left(s(\kappa)^{-1}\right)^{*}\left(=s(\kappa)^{*}\right)$ agree on $e_{\alpha_{i}}$ and $h$ so that $s(\kappa)=$ $s(\kappa)^{*}$. Relation (5.2.8) also implies that $m=m^{*}$ so that $s_{1}(\kappa)=s_{1}(\kappa)^{*}$. But
$s_{1}(\kappa) \in G$ by Lemma 5.2 .2 so that $s_{1}(\kappa) \in K$ by Lemma 3.1.2. However, as one knows, $K^{s}$ is the inverse image of $K$ in $G^{s}$ under Ad. Thus $s_{o}(\kappa) \in K^{s}$ since $\operatorname{Ad} s_{0}(\kappa)=s_{1}(\kappa)$ by Lemma 5.2.2.
Q.E.D.

Let $\lambda \in D$. Now recall we have fixed a highest weight vector $v^{\lambda} \in V^{\lambda}(\lambda)$. On the other hand one knows the set of weights of $\pi_{\lambda}$ is stable under the Weyl group and in fact since $s_{0}(\kappa) \in G^{s}$ corresponds to $\kappa$ one must have, for any $\mu \in \Lambda$

$$
\begin{equation*}
s_{0}(\kappa) V_{\mathbb{C}^{\lambda}}(\mu)=V_{\mathbb{C}^{\lambda}}(\kappa \mu) \tag{5.2.9}
\end{equation*}
$$

Furthermore $\kappa \lambda$ is the lowest weight of $\pi \lambda$ and hence if we fix $v^{\kappa \lambda}$ by putting

$$
\begin{equation*}
v^{\kappa \lambda}=s_{o}(\kappa) v^{\lambda} \tag{5.2.10}
\end{equation*}
$$

then $v^{\kappa \lambda}$ is a lowest weight vector. In fact one has
Proposition 5.2. Let $\lambda \in D$; then $V^{\lambda}(\kappa \lambda)$ is one dimensional over $\mathbb{R}$ and $\boldsymbol{v}^{\kappa \lambda} \in V^{\lambda}(\kappa \lambda)$, where

$$
\begin{equation*}
\left\{v^{\kappa \lambda}, v^{\kappa \lambda \lambda}\right\}=1 \tag{5.2.11}
\end{equation*}
$$

Proof. One has $\operatorname{dim} V^{\lambda}(\kappa \lambda)=1$ by (5.2.9) and Remark 5.1. But also one must have $s_{0}(\kappa) v^{\lambda} \in V^{\lambda}(\kappa \lambda)$ since, by Proposition 5.1, $V^{\lambda}$ is stable under $G^{s}$. But now one has (5.3.3) by Lemma 5.2.3 since, clearly, $K^{s} \subseteq G_{u}{ }^{s}$ so that $Q_{*}^{\lambda}$ is certainly invariant under $K^{s}$.
Q.E.D.
5.3. Now let $w_{o} \in h_{+}$and conforming to our standard notation we put $w=f+w_{o}$. Our main concern here is with the determination of $h(a)^{\lambda}$ for $a \in G_{0}{ }^{w}$ (see Section 3.5) and $\lambda \in \hbar^{\prime}$. We first note

Lemma 5.3. For any $a \in G_{0}{ }^{w}$ put $\rho_{w}(a)=m^{-1} \bar{n}_{f}(w)^{-1} a \bar{n}_{f}(w)$. Then $\rho_{w}(a) \in H$ and

$$
\begin{equation*}
\rho_{w}: G_{0}^{w} \rightarrow H, \quad a \mapsto \rho_{w}(a) \tag{5.3.1}
\end{equation*}
$$

is a diffeomorphism.
Proof. This is obvious from (3.5.7), which asserts that $\bar{n}_{f}(w) m H \bar{n}_{f}(w)^{-1}=$ $G_{0}{ }^{w}$.
Q.E.D.

We will generally conform to the notation of using the subscript $o$ to denote the image in $H$ of an element in $G_{0}{ }^{w}$ under $\rho_{w}$. Thus $\rho_{w}(a)=a_{0} \in H$ for any $a \in G_{0}{ }^{w}$.

Our determination of $h(a)^{\lambda}$ begins with

Proposition 5.3. Let $w_{o} \in h_{+}$and let $\lambda \in D$. Put $w=f+w_{o}$. Then for any $a \in G_{0}{ }^{2 w}$ one has (recalling (3.4.14) and (3.5.9))

$$
\begin{equation*}
h(a)^{\lambda}=\left\{a_{0} \bar{n}_{f}(w)^{-1} v^{\lambda}, \bar{n}_{-f}(w)^{*} v^{\kappa \lambda}\right\} \tag{5.3.2}
\end{equation*}
$$

where $\bar{n}_{f}(w)$ and $\bar{n}_{-f}(w)$ are defined by (3.5.2) and (5.2.2) and $a_{0}=\rho_{w}(a) \in H$, where $\rho_{w}$ is defined by Lemma 5.3. Also the right side of (5.3.2) is a $Q_{*}^{\lambda}$ inner product (see Proposition 5.1) of vectors in the $G^{s}$-module $V^{\lambda}$.

Proof. Let $a \in G_{0}{ }^{w}$. Then by (3.4.14) and (3.5.9) one has, in $\mathcal{G}^{*}$,

$$
\begin{equation*}
s(\kappa)^{-1} a=\bar{n}(a) h(a) n(a) \tag{5.3.3}
\end{equation*}
$$

On the other hand if $a_{o}=\rho_{w}(a) \in H$ then clearly $a=\bar{n}_{f}(w) m a_{o} \bar{n}_{f}(w)^{-1}=$ $m \bar{n}_{-f}(w) a_{0} \bar{n}_{f}(w)^{-1}$, recalling (5.2.2). Hence

$$
\begin{equation*}
s(\kappa)^{-1} a=\left(s(\kappa)^{-1} m\right) \bar{n}_{-f}(w) a_{0} \bar{n}_{f}(w)^{-1} \tag{5.3.4}
\end{equation*}
$$

in $\boldsymbol{G}$. But then by (5.3.3) one has

$$
\begin{equation*}
\bar{n}(a) h(a) n(a)=\left(s(\kappa)^{-1} m\right) \bar{n}_{-f}(w) a_{0} \bar{n}_{f}(w)^{-1} . \tag{5.3.5}
\end{equation*}
$$

On the other hand, in $G^{s}$, one has, recalling Lemma 5.2.2, $s_{o}(\kappa)^{-1} \bar{n}_{-f}(w) a_{o} \bar{n}_{f}(w) \in G_{*}=G_{*}^{*}$. Thus by Lemma 5.2.2 one has

$$
\begin{equation*}
\bar{n}(a) h(a) n(a)=s_{o}(\kappa)^{-1} \bar{n}_{-f}(w) a_{o} \bar{n}_{f}(w)^{-1} \tag{5.3.6}
\end{equation*}
$$

in $G^{s}$.
Now consider the $G^{s}$-module $V^{\lambda}$. Let $g \in G$ be the element given by (5.3.6). Now $n(a) v^{\lambda}=v^{\lambda}$ since $v^{\lambda}$ is a highest weight vector of $V^{\lambda}$. But $(\bar{n}(a))^{*} \in N$ by (1.5.1) and (3.1.9) so that one also has $(\bar{n}(a))^{*} v^{\lambda}=v^{\lambda}$. Thus from the left side of (5.3.6) one has $\left\{g v^{\lambda}, v^{\lambda}\right\}=\left\{h(a) v^{\lambda}, v^{\lambda}\right\}=\left\{h(a)^{\lambda} v^{\lambda}, v^{\lambda}\right\}=h(a)^{\lambda}$. But now $s_{o}(\kappa) \in K^{s}$ by Lemma 5.2.3. But of course $K^{s} \subseteq G^{s}$ so that $s_{o}(\kappa)^{-1}=s_{o}(\kappa)^{*}$ and hence $\left(s_{o}(\kappa)^{-1}\right)^{*} v^{\lambda}=s_{o}(\kappa) v^{\lambda}=v^{\kappa \lambda}$. Thus from the right side of (5.3.6) one has $\left\{g v^{\lambda}, v^{\lambda}\right\}=\left\{\bar{n}_{-f}(w) a_{0} \bar{n}_{f}(w)^{-1} v^{\lambda}, v^{\kappa \lambda}\right\}=\left\{a_{0} \bar{n}_{f}(w)^{-1} v^{\lambda}, \bar{n}_{-f}(w)^{*} v^{\kappa \lambda}\right\}$. This proves the lemma since $\left\{g v^{\lambda}, v^{\lambda}\right\}=h(a)^{\lambda}$.
Q.E.D.

Remark 5.3. Since $\operatorname{dim} V^{\lambda}(\kappa \lambda)=1$ there are exactly two vectors $v \in V^{\lambda}(\kappa \lambda)$ such that $\{v, v\}=1$. Thus if one were to choose a normalized lowest weight vector there would be an ambiguity up to sign. But now by (5.2.11), $v^{\kappa \lambda \lambda}$ is a choice of one of these two vectors. We now observe that since $h(a)^{\lambda}>0$ Proposition 5.3 implies that $v^{\kappa \lambda}$ is that choice such that $\left\{a_{o} \bar{n}_{f}(w)^{-1} v^{\lambda}, \bar{n}_{-f}(w)^{*} v^{\kappa \lambda}\right\}>0$ for all $a_{o} \in H$ (see (5.3.1)) and all $w_{o} \in h_{+}$.
5.4. Now let $w_{o} \in h_{+}$, and, as usual, put $w=f+w_{0}$. Let $g_{o} \in H$ and put $g=\rho_{w}^{-1}\left(g_{0}\right) \in G_{0}{ }^{w}$. See (5.3.1). Let $\bar{n}(g) \in \bar{N}, h(g) \in H$, and $\bar{n}(g) \in \bar{N}$ be defined by (3.4.14) recalling (3.5.9). Now since $H$ is a vector group we can consider the square root $h(g)^{1 / 2}$ of $h(g)$ in $H$. Let $P^{s}=\exp ^{s} p$. Using $K^{s}$ and $P^{s}$ we can, as one knows, take the polar decomposition of any element in $G^{s}$. Thus there exist unique elements $k\left(g_{o}, w_{o}\right) \in K^{s}$ and $p\left(g_{o}, w_{o}\right) \in P^{s}$ such that, in $G^{s}$,

$$
\begin{equation*}
h(g)^{1 / 2} n(g) \bar{n}_{f}(w)=k\left(g_{o}, w_{o}\right) p\left(g_{o}, w_{o}\right) \tag{5.4.1}
\end{equation*}
$$

where $\bar{n}_{f}(w)$ as usual is given by (3.5.2).
Now recalling the set Jac $g$ of Jacobi elements in $g$ (see (2.2.6)) let

$$
\begin{equation*}
\mathrm{Jac} p=\mathrm{Jac} g \cap \not p \tag{5.4.2}
\end{equation*}
$$

so that $\mathrm{Jac} \nsim$ is the set of all symmetric Jacobi elements. Using the notation (2.2.6) so that $y \in \mathrm{Jac} g$ if and only if $y$ is of the form

$$
\begin{equation*}
y=x+\sum_{i=1}^{l} a_{-i} e_{-\alpha_{i}}+\sum_{i=1}^{l} a_{i} e_{\alpha_{i}} \tag{5.4.3}
\end{equation*}
$$

where $a_{-i}>0$ and $a_{i}>0$ and $x \in h$, one notes by (1.5.1) that $y \in \mathrm{Jac} \nsim$ if and only if $a_{-i}=a_{i}, i=1, \ldots, l$. It is clear then that $\mathrm{Jac} \not \mu$ is a closed connected two-dimensional submanifold of $\mathrm{Jac} g$.

Now for any $\gamma \in \mathscr{I}\left(h_{+}\right)$(see (2.3.1)) let

$$
(\operatorname{Jac} \not p)(\gamma)=\mathscr{I}^{-1}(\gamma) \cap \mathrm{Jac} \not p
$$

so that the sets $(\mathrm{Jac} \not \mu)(\gamma), \gamma \in \mathscr{I}\left(h_{+}\right)$, are the "isospectral" equivalence classes of $\mathrm{Jac} \not p$. The following is a corollary of Theorem 3.6 and hence could have been proved earlier. It is proved now instead for notational convenience.

Theorem 5.4. For any $w_{o}$ in the open Weyl chamber $h_{+}$(see (3.3.1)) let $w=f+w_{o}$, where $f$ is given by (1.5.4). Let $G_{0}{ }^{w}$ be defined by (3.5.7) and for any $g_{o}$ in the split Cartan subgroup $H$ let $\rho_{w}^{-1}\left(g_{0}\right)=g \in G_{0}{ }^{w}$, where $\rho_{w}$ is defined as in (5.3.1). Then if $p\left(g_{o}, w_{o}\right)$ is defined by (5.4.1) one has

$$
\begin{equation*}
p\left(g_{o}, w_{0}\right) \in H \tag{5.4.4}
\end{equation*}
$$

On the other hand if $k\left(g_{o}, w_{o}\right) \in K^{8}$ is defined by (5.4.1) one has (regarding $g$ as a $G^{s}$-module using Ad )

$$
\begin{equation*}
k\left(g_{o}, w_{o}\right) w_{o} \in(\operatorname{Jac} \not \beta)(\gamma), \tag{5.4.5}
\end{equation*}
$$

where $\gamma=\mathscr{I}\left(w_{o}\right)$. Moreover the map

$$
\begin{equation*}
H \times h_{+} \rightarrow \mathrm{Jac} \not h, \quad\left(g_{o}, w_{o}\right) \mapsto k\left(g_{o}, w_{o}\right) w_{0} \tag{5.4.6}
\end{equation*}
$$

is a diffeomorphism. Furthermore if $w_{o} \in h_{+}$and $\gamma=\mathscr{I}\left(w_{0}\right)$ then $(\mathrm{Jac} \not p)(\gamma)$ is a closed connected submanifold of dimension 1 in Jac $\nsim$ and

$$
\begin{equation*}
H \rightarrow(\mathrm{Jac} \not p)(\gamma), \quad g_{o} \rightarrow k\left(g_{0}, w_{o}\right) w_{0} \tag{5.4.7}
\end{equation*}
$$

is a diffeomorphism so that as manifolds $(\mathrm{Jac} \nmid)(\gamma) \cong \mathbb{R}^{2}$. Finally

$$
\begin{equation*}
\operatorname{Jac} \not h=\bigcup_{w_{o} \in h_{+}}(\operatorname{Jac} \not n)\left(\mathscr{I}\left(w_{o}\right)\right) \tag{5.4.8}
\end{equation*}
$$

is a disjoint union and in fact (5.4.8) is the decomposition of Jac $\neq 4$ as the union of leaves of a foliation.

Proof. Let $Z \subseteq$ Jac $g$ be as in (2.2.4) (where of course $F=\mathbb{R}$ ). Let $y \in Z$ so that we may write $y=f+x+\sum_{i=1}^{l} r_{i} e_{\alpha_{i}}$. Let $a_{v} \in H$ be defined by the condition $\left(a_{y}\right)^{\alpha_{i}}=r_{i}^{-1 / 2}, i=1, \ldots, l$. See (2.2.9). Obviously the map

$$
\begin{equation*}
Z \rightarrow H, \quad y \mapsto a_{v} \tag{5.4.9}
\end{equation*}
$$

is smooth. But clearly $a_{v} y \in \mathrm{Jac} \not p$ and any element in Jac $\nsim$ is of this form. On the other hand it is easy to see that the bijection (2.2.8) is a diffeomorphism and hence by considering the graph of (5.4.9) the map

$$
\begin{equation*}
Z \rightarrow \mathrm{Jac} \not 2, \quad y \rightarrow a_{y} y \tag{5.4.10}
\end{equation*}
$$

is a diffeomorphism. Obviously then (5.4.10) induces a diffeomorphism

$$
\begin{equation*}
Z(\gamma) \rightarrow(\mathrm{Jac} \not p)(\gamma) \tag{5.4.11}
\end{equation*}
$$

for any $\gamma \in \mathscr{I}\left(h_{+}\right)$. But then, by Theorem 3.2, $(\mathrm{Jac} \not)(\gamma)$ is a connected, closed submanifold of dimension $l$ of Jac $p$. The final statement of Theorem 5.4 and the fact that

$$
\begin{equation*}
(\mathrm{Jac} \nsim)(\gamma) \cong \mathbb{R}^{l} \tag{5.4.12}
\end{equation*}
$$

follow from (3.6.3), (3.2.14), and Proposition 4.1.
Now let $g_{o} \in H$ and $w_{o} \in h_{+}$. Recall Theorem 3.7. Let $w=f+w_{o}$ so that $\bar{n}_{f}(w) w_{o}=w$ and let $g \in H$ be given by (3.7.1). Let $y=n(g) w$ so that $y \in Z(\gamma)$ where $\gamma=\mathscr{I}\left(w_{o}\right)$. But now $a_{y}=h(g)^{1 / 2}$ by (3.7.4). Thus if $z=h(g)^{1 / 2} y$ then
$z \in \mathrm{Jac} \not p$ by (5.4.9). But now $z=h(g)^{1 / 2} y=h(g)^{1 / 2} n(g) w=h(g)^{1 / 2} n(g) \bar{n}_{f}(w) w_{o}$. Thus

$$
\begin{equation*}
z=k p w_{o} \tag{5.4.13}
\end{equation*}
$$

by (5.4.1), where we have written $k=k\left(g_{0}, w_{o}\right)$ and $p=p\left(g_{o}, w_{o}\right)$. Thus if $v=k^{-1} z$ then $p w_{o}=v$. But $v \in \not \subset$ since $z \in \neq$ and $k \in K^{s}$. Thus if one applies the $*$-operation it follows from (5.2.8), which is clearly valid if $G_{\mathbb{C}}{ }^{s}$ is substituted for $G_{\mathbb{C}}$, that also $p^{-1} w_{o}=v$. Thus $p^{2} w_{o}=w_{o}$, which clearly implies $p w_{o}=w_{o}$ (since $\operatorname{Ad} p$ is diagonalizable with positive eigenvalues). Thus $\operatorname{Ad} p \in G^{w_{0}} \cap P$. But $G^{w_{o}} \cap P=H$ by (3.1.27). However, Ad $H=H$ by Identification 5.2 and as one knows the map $P^{s} \rightarrow P$ induced by Ad is bijective. Thus $p \in H$, proving (5.4.4). But also the relation $p w_{o}=w_{o}$ implies $k w_{o}=z$ by (5.4.13). This proves (5.4.5).

Now the map (5.4.7) is the composite of (5.3.1), (3.6.2), and (5.4.11). Since these three maps are diffeomorphisms it follows that (5.4.7) is a diffeomorphism. Furthermore the map (5.4.6) is a composite of (3.7.2) and the inverse to (5.4.10). Since these two maps are diffeomorphisms it follows also that (5.4.6) is a diffeomorphism.
Q.E.D.

Remark 5.4. A familiar question when dealing with symmetric Jacobi matrices is the problem of diagonalization, even if one knows the spectrum of the matrix. In a more general setting, if $x \in \mathrm{Jac} \nsim$ and $\lambda \in D$ find the eigenvectors of $\pi_{\lambda}(x)$ in terms of an orthonormal basis $v_{i}{ }^{\lambda}, i=1, \ldots, \operatorname{dim} V^{\lambda}$, of weight vectors in $V^{\lambda}$. Now let $\left(g_{0}, w_{o}\right) \in H \times h_{+}$correspond to $x$ under the bijection (5.4.6). Analogous to knowing the spectrum of a Jacobi matrix is the knowledge of the element $w_{o}$. Since the $v_{i}{ }^{\lambda}$ are eigenvectors of $\pi_{\lambda}\left(w_{o}\right)$ it is now clear from (5.4.5) that $\left\{\pi_{\lambda}\left(k\left(g_{o}, w_{o}\right)\right) v_{i}{ }^{\lambda}\right\}$ is indeed an orthonormal basis of eigenvectors of $\pi_{\lambda}(x)$. The problem then reduces to determining $k\left(g_{o}, w_{o}\right)$. By (5.4.1) it is then a question of determining $h(g), n(g)$, and $\bar{n}_{f}(w)$. The element $h(g)$ is given by (3.7.4). We will give a formula for $\bar{n}_{f}(w)^{-1}$ in Section 5.8. (The formula for $\bar{n}_{f}(w)$ is obtained from (5.8.7) by eliminating $(-1)^{(s)}$ and changing $\bar{s}$ to $s$.) Since $n(g)$ is unipotent one solves inductively for $n(g)$ by the relation $n(g) w=y$. See Section 7.8, where an example is worked out.
5.5. Let $\lambda \in D$. Onc notes that if $v$ is a weight vector of $\lambda$ then, by (5.1.9) and (5.1.14),

$$
\begin{equation*}
\{n v, v\}=\{\bar{n} v, v\}=\{v, v\} \tag{5.5.1}
\end{equation*}
$$

for any $n \in N$ or $\bar{n} \in \bar{N}$.
Proposition 5.5.1. Let $g_{o} \in H$ and $w_{o} \in h_{+}$. Write $p=p\left(g_{o}, w_{o}\right)$, where the latter is given by (5.4.1). We recall $p \in H$ by (5.4.4). Then $g_{o}^{-1} p^{2}$ is independent
of $g_{0}$, depending only on $w_{o}$ (or equivalently $w-f+w_{o}$ ) so that we can define $d(w) \in H$ by

$$
\begin{equation*}
g_{o}^{-1} p\left(g_{o}, w_{o}\right)^{2}=d(w) \tag{5.5.2}
\end{equation*}
$$

Furthermore if $\lambda \in D$ then

$$
\begin{equation*}
\left\{(d(w)) v,\left(\bar{n}_{f}(w)\right)^{-1} v^{\lambda}\right\}=\left\{\bar{n}_{-f}(w) v, v^{\kappa \lambda}\right\} \tag{5.5.3}
\end{equation*}
$$

for any $v \in V^{\lambda}$, where $\bar{n}_{f}(w)$ and $\bar{n}_{-f}(w)$ are defined by (3.5.2) and (5.2.2). Moreover one has

$$
\begin{equation*}
d(w)^{\lambda}=\left\{\bar{n}_{-f}(w) v^{\lambda}, v^{k \lambda \lambda}\right\} \tag{5.5.4}
\end{equation*}
$$

and

$$
\begin{equation*}
d(w)^{-\kappa \lambda}=\left\{\bar{n}_{f}(w)^{-1} v^{\lambda}, v^{\kappa \lambda}\right\} . \tag{5.5.5}
\end{equation*}
$$

Proof. Put $d=g_{o}^{-1} p^{2}$. Let $\lambda \in D$ and $v \in V^{\lambda}$ and put $C=\left\{d v, \bar{n}_{f}(w)^{-1} v^{\lambda}\right\}$. Then $C=\left\{p g_{o}^{-1} v, p\left(\bar{n}_{f}(w)\right)^{-1} v^{\lambda}\right\}$ since $g_{o}$ and $p$ commute and $p=p^{*}$. But $p=k^{-1} h^{1 / 2} n \bar{n}_{f}(w)$ in $G^{s}$ by (5.4.1), where we have written $k=k\left(g_{o}, w_{o}\right)$, $h=h(g)$, and $n=n(g)$. Then substituting for $p$ and canceling out $k^{-1}$ (since $\left(k^{-1}\right)^{*}=k$ ) one has $C=\left\{h^{1 / 2} n \bar{n}_{f}(w) g_{o}^{-1} v, h^{1 / 2} n v^{\lambda}\right\}$. But $n v^{\lambda}=v^{\lambda}$. Hence by transposing $h^{1 / 2}$ one has

$$
\begin{align*}
C & =\left\{h n \bar{n}_{f}(w) g_{o}{ }^{1} v, v^{\lambda}\right\}  \tag{5.5.6}\\
& =\left\{\bar{n} h n \bar{n}_{f}(w) g_{o}^{-1} v, v^{\lambda}\right\},
\end{align*}
$$

where $\bar{n}=\bar{n}(g)$, since $\bar{n}^{*} v^{\lambda}=v^{\lambda}$ (recalling that $\left.\bar{n}^{*} \in N\right)$. But $\bar{n} h n=$ $s_{o}(\kappa)^{-1} \bar{n}_{-f}(w) g_{o} \bar{n}_{f}(w)^{-1}$ by (5.3.6) in $G^{s}$. Thus $C=\left\{s_{o}(\kappa)^{-1} \bar{n}_{-f}(w) v, v^{\mu}\right\}=$ $\left\{\bar{n}_{-f}(w) v, v^{\kappa \lambda}\right\}$. That is,

$$
\begin{equation*}
\left\{d v, \bar{n}_{f}(w)^{-1} v^{\lambda}\right\}=\left\{\bar{n}_{-f}(w) v, v^{\kappa \lambda \lambda}\right\} . \tag{5.5.7}
\end{equation*}
$$

But now if we choose. $v=v^{\lambda}$ then $d v^{\lambda}=d^{\lambda} v^{\lambda}$. But $\left\{v^{\lambda}, \bar{n}_{f}(w)^{-1} v^{\lambda}\right\}=1$ by (5.5.1). Thus

$$
\begin{equation*}
d^{\lambda}=\left\{\bar{n}_{-f}(w) v^{\lambda}, v^{\kappa \lambda}\right\} . \tag{5.5.8}
\end{equation*}
$$

Obviously the right side of (5.5.8) depends only on $w_{o}$ and not on $g_{o}$. Furthermore this is true for all $\lambda \in D$. But $d$ is determined by $d^{v_{i}}, i=1, \ldots, l$, where $\nu_{i} \in D$ is as in (5.1.2), since the $\nu_{i}$ are a basis of $\ell^{\prime}$. Thus $d$ depends only on $w$ proving the first statement. Also, $d=d(w)$. But then (5.5.3) is just (5.5.7) and (5.5.4) is just (5.5.8). If we put $v=v^{\kappa \lambda}$ in (5.5.7) then since $\left\{\bar{n}_{-f}(w) v^{\kappa \lambda}, v^{\kappa \lambda \lambda}\right\}=1$ by (5.5.1) and $d(w) v^{\kappa \lambda}=d(w)^{\kappa \lambda} v^{\kappa \lambda}$ one has $d(w)^{\kappa \lambda}\left\{v^{\kappa \lambda}, \bar{n}_{f}(w)^{-1} v^{\lambda}\right\}=1$. But then (5.5.8) follows by inverting and reversing the order of the vectors. Q.E.D.

If $\mu, \nu \in \Lambda$ we will say that $\nu \geqslant \mu$ if $\nu-\mu$ is in the $\mathbb{Z}_{+}$-cone generated by the simple roots $\alpha_{1}, \ldots, \alpha_{l}$. If in addition $\nu \neq \mu$ we write $\nu>\mu$. One notes that

$$
\begin{equation*}
\nu>\mu \text { implies }\left\langle\nu, w_{o}\right\rangle>\left\langle\mu, w_{o}\right\rangle \text { for any } w_{o} \in \hbar_{+} . \tag{5.5.9}
\end{equation*}
$$

Now for any $\lambda \in D$ let $q(\lambda)=\operatorname{dim} V^{\lambda}$. Let $\lambda \in D$. By the orthogonal direct sum (5.1.13) there exists a $Q_{*}^{\lambda}$-orthonormal basis $v_{i}{ }^{\lambda}, i=1, \ldots, q(\lambda)$, of $V^{\lambda}$ such that the $v_{i}{ }^{\lambda}$ are weight vectors. We assume that such a basis is chosen once and for all where $v_{1}^{\lambda}=v^{\lambda}$ and $v_{q(\lambda)}^{\lambda}=v^{\kappa \lambda}$. See (5.2.10). Let $\lambda_{i} \in \Lambda$, $i=1, \ldots, q(\lambda)$, be the weight corresponding to $v_{i}{ }^{\lambda}$. Thus $\lambda_{1}=\lambda$ and $\lambda_{Q(\lambda)}=\kappa \lambda$. Furthermore since as one knows

$$
\begin{equation*}
V^{\lambda}=U(\bar{n}) v^{\lambda}=U(n) v^{\kappa \lambda} \tag{5.5.10}
\end{equation*}
$$

it follows from (5.1.9) that

$$
\begin{equation*}
\lambda>\lambda_{i}>\kappa \lambda \tag{5.5.11}
\end{equation*}
$$

for all $1<i<q(\lambda)$.
Now let $w_{o} \in h_{+}$and let $w=f+w_{o}$. Also let $\bar{n}_{f}(w)$ and $\bar{n}_{-f}(w)$ be as in Proposition 5.5. Now put

$$
\begin{equation*}
b_{i}(\lambda, w)=\left\{\bar{n}_{f}(w)^{-1} v^{\lambda}, v_{i}^{\lambda}\right\}\left\{\bar{n}_{-f}(w) v_{i}^{\lambda}, v^{\star \lambda}\right\} . \tag{5.5.12}
\end{equation*}
$$

Lemma 5.5.1. One has $b_{i}(\lambda, w) \geqslant 0$ for $1 \leqslant i \leqslant q(\lambda)$. Furthermore one has strict positivity at the extremes. That is,

$$
\begin{equation*}
b_{1}(\lambda, w)=\left\{\bar{n}_{-f}(w) v^{\lambda}, v^{\kappa \lambda}\right\}>0 \tag{5.5.13}
\end{equation*}
$$

and

$$
\begin{equation*}
b_{q(\lambda)}(\lambda, w)=\left\{\bar{n}_{f}(w)^{-1} v^{\lambda}, v^{\kappa \lambda \lambda}\right\}>0 . \tag{5.5.14}
\end{equation*}
$$

Proof. Let $d(w) \in H$ be defined by (5.5.2). Then $\left\{d(w) v_{i}{ }^{\lambda}, \bar{n}_{f}(w)^{-1} v^{\lambda}\right\}=$ $d(w)^{\lambda_{i}}\left\{\bar{n}_{f}(w)^{-1} v^{\lambda}, v_{i}^{\lambda}\right\}$. But this equals $\left\{\bar{n}_{-f}(w) v_{i}^{\lambda}, v^{k \lambda}\right\}$ by (5.5.3) and hence

$$
\begin{equation*}
d(w)^{\lambda_{i}} b_{i}(\lambda, w)=\left\{\bar{n}_{-f}(w) v_{i}^{\lambda}, v^{\kappa \lambda}\right\}^{2} . \tag{5.5.15}
\end{equation*}
$$

But then (5.5.15) is nonnegative. Since $d(w)^{\lambda_{i}}>0$ this proves $b_{i}(\lambda, w) \geqslant 0$. The equality in (5.5.13) and (5.5.14) is immediate from (5.5.1) and (5.5.12). The positivity then follows from the identities

$$
\begin{align*}
b_{1}(\lambda, w) & =d(w)^{\lambda}  \tag{5.5.16}\\
b_{q(\lambda)}(\lambda, w) & =d(w)^{-\kappa \lambda} \tag{5.5.17}
\end{align*}
$$

See (5.5.4) and (5.5.5).
Q.E.D.

Lemma 5.5.2. Let $a_{0} \in H, w_{o} \in \mathscr{h}_{+}$and let $\lambda \in D$. Put $a=\rho_{w}^{-1}\left(a_{o}\right) \in G_{0}{ }^{w}$ and $w=f+w_{o}$. Then

$$
\begin{equation*}
h(a)^{\lambda}=\sum_{i=1}^{q(\lambda)} b_{i}(\lambda, w) a_{o}^{\lambda_{i}} \tag{5.5.18}
\end{equation*}
$$

Proof. By (5.3.2) one has $h(a)^{\lambda}=\left\{a_{o} \bar{n}_{f}(w)^{-1} v^{\lambda}, \bar{n}_{-f}(w)^{*} v^{\kappa \lambda}\right\}$. Hence if $c_{i}=\left\{a_{0} \bar{n}_{f}(w)^{-1} v^{\lambda}, v_{i}^{\lambda}\right\}$ and $c_{i}^{\prime}=\left\{v_{i}^{\lambda}, \bar{n}_{-f}(w)^{*} v^{\kappa \lambda}\right\}=\left\{\bar{n}_{-f}(w) v_{i}^{\lambda}, v^{\kappa \lambda \lambda}\right\}$ one has

$$
\begin{equation*}
h(a)^{\lambda}=\sum_{i=1}^{q(\lambda)} c_{i} c_{i}^{\prime} . \tag{5.5.19}
\end{equation*}
$$

But $c_{i}=\left\{\bar{n}_{f}(w)^{-1} v^{\lambda}, a_{o} v_{i}^{\lambda}\right\}=a_{o}^{\lambda_{i}}\left\{\bar{n}_{f}(w)^{-1} v^{\lambda}, v_{i}^{\lambda}\right\}$ since $a_{o}=a_{o}^{*}$, recalling that $a_{0} \in H \subseteq P$. Thus $c_{i} c_{i}^{\prime}=b_{i}(\lambda, w) a_{o}^{\lambda_{i}}$. But then (5.5.19) implies (5.5.18). Q.E.D.

Proposition 5.5.2. Let $w_{o} \in \mathscr{h}_{+}$and let $w=f+w_{0}$. Let $G_{0} w$ be as in Theorem 3.6. Let $g \in G_{0}{ }^{w}$ and put $g_{o}=\rho_{w}(g) \in H$, where $\rho_{w}$ is defined by (5.3.1). Then for all $t \in \mathbb{R}$ one has $g \exp t w \in G_{0}{ }^{w}$ and

$$
\begin{equation*}
\rho_{w}(g \exp t w)=g_{a} \exp t w_{o} . \tag{5.5.20}
\end{equation*}
$$

Now let $\lambda \in D$ and let $b_{i}(\lambda, w) \in \mathbb{R}$ be defined by (5.5.12), $i=1, \ldots, q(\lambda)=$ : $\operatorname{dim} V^{\lambda}$. Then

$$
\begin{equation*}
h(g \exp t w)^{\lambda}=\sum_{i=1}^{q(\lambda)} b_{i}(\lambda, w) g_{0}^{\lambda_{i}} e^{t\left\langle\lambda_{i}, w_{o}\right\rangle} \tag{5.5.21}
\end{equation*}
$$

where $h(a)$ is defined by (3.4.14). See (3.5.9).
Proof. Now, recalling (3.5.7) and (3.5.8), $G_{1}{ }^{w}$ is the identity component of $G^{w}$. Thus $G_{0}{ }^{w} G_{1}{ }^{w} \subseteq G_{0}{ }^{w}$. But of course $\exp t w \in G_{1}{ }^{w}$. Thus $g \exp t w \in G_{0}{ }^{w}$. Furthermore $\rho_{w}(g \exp t w)=m^{-1} \bar{n}_{f}(w)^{-1} g \operatorname{cxp} t w \bar{n}_{f}(w)-\left(m^{-1} \bar{n}_{f}(w)^{-1} g \bar{n}_{f}(w)\right) \bar{n}_{f}(w)^{-1}$ $\exp t w \bar{n}_{f}(w)=\rho_{w}(g) \bar{n}_{f}(w)^{-1} \exp t w \bar{n}_{f}(w)$. But $\bar{n}_{f}(w)^{-1} \exp t w \bar{n}_{f}(w)=\exp t w_{o}$ since $\bar{n}_{f}(w) w_{o}=w$. Hence $\rho_{w}(g \exp t w)=\rho_{w}(g) \exp t w_{o}=g_{o} \exp t w_{o}$, proving (5.5.20). Thus if $a_{o}=g_{o} \exp t w_{o}$ then $a=g \exp t w$ using the notation of Lemma 5.5 .2 and (5.5.20). But then $a_{o}^{\lambda_{i}}=g_{o}^{\lambda_{i} e^{t\left(\lambda_{i}, w_{0}\right)} \text {. Proposition } 5.5 .2 \text { then }}$ follows from (5.5.18).
Q.E.D.

The following result will later play the key role in determining the "phase" for scattering in the generalized Toda lattice.

We recall that $\kappa$ is the element in the Weyl group $W$ which maps the set of positive roots into the set of negative roots. Now we may regard $W$ as operating on $H$ as well as $h$. If $d \in H$ we will write $d^{-\kappa}$ for $\kappa\left(d^{-1}\right)$. One notes that $\left(d^{-\kappa}\right)^{\lambda}=d^{-\kappa \lambda}$ for any $\lambda \in \Lambda$.

Theorem 5.5. Let $w_{o}$ be arbitrary in the open Weyl chamber $h_{+}$and let $g_{o}$ be arbitrary in the split Cartan subgroup H. Put $w=f+w_{0}$, where $f$ is defined by (1.5.4). Let $G_{0}{ }^{w}$ be defined by (3.5.5-3.5.7) and let $g=\rho_{w}^{-1}\left(g_{o}\right)$, where $\rho_{w}$ is defined by (5.3.1). Now for any $t \in \mathbb{R}$ one has $g \exp t w \in G_{0}{ }^{w}$ and $\rho_{w}(g \exp t w)=$ $g_{o} \exp t w_{o} \in H$. Let $h(g \exp t w) \in H$ be defined by (3.4.14). See (3.5.9), (3.4.13), and (3.2.1). Then the curve $t \rightarrow h(g \exp t w)\left(g_{o} \exp t w_{o}\right)^{-1}$ in $H$ converges as $t \rightarrow \infty$ and the curve $t \rightarrow h(\exp t w)\left(g_{o} \exp w_{o}\right)^{-\kappa}$ in $H$ converges as $t \rightarrow-\infty$. Furthermore the limits are independent of $g_{o}$. In fact if $d(w) \in H$ is defined by (5.5.2) then, in $H$,

$$
\begin{equation*}
\lim _{t \rightarrow \infty} h(g \exp t w)\left(g_{o} \exp t w_{o}\right)^{-1}=d(w) \tag{5.5.22}
\end{equation*}
$$

and

$$
\begin{equation*}
\lim _{t \rightarrow-\infty} h(g \exp t w)\left(g_{o} \exp t w_{a}\right)^{-\kappa}=(d(w))^{-\kappa} \tag{5.5.23}
\end{equation*}
$$

Proof. Let $c(t)=h(g \exp t w)\left(g_{o} \exp t w_{o}\right)^{-1} \quad$ and $\quad$ let $\quad \tilde{c}(t)=h(g \exp t w)$ $\left(g_{v} \exp t w_{v}\right)^{-\kappa}$. Now for any $\lambda \in D$ one has by (5.5.21)

$$
c(t)^{\lambda}=b_{1}(\lambda, w)+\sum_{i=2}^{q(\lambda)} b_{i}(\lambda, w) g_{o}^{\left.\lambda_{i}-\lambda^{t} t \lambda_{i}-\lambda, w_{o}\right\rangle}
$$

and

$$
\tilde{c}(t)^{\lambda}=b_{q(\lambda)}(\lambda, w)+\sum_{i=1}^{q(\lambda)-1} b_{i}(\lambda, w) g_{o}^{\lambda_{i}-\kappa \lambda} e^{t\left\langle\lambda_{i}-\kappa \lambda, w_{o}\right\rangle}
$$

But now $\left\langle\lambda_{i}-\lambda, w_{o}\right\rangle<0$ and $\left\langle\lambda_{i}-\kappa \lambda, w_{o}\right\rangle>0$ by (5.5.11). Thus $\lim _{t \rightarrow \infty} c(t)^{\lambda}=$ $b_{1}(\lambda, w)$ and $\lim _{t \rightarrow-\infty} \tilde{c}(t)^{\lambda}=b_{q(\lambda)}(\lambda, w)$. But then $\lim _{t \rightarrow \infty} c(t)^{\lambda}=d(w)^{\lambda}$ and $\lim _{t \rightarrow-\infty} \tilde{c}(t)^{\lambda}=(d(w))^{-\kappa \lambda}$ by (5.5.16) and (5.5.17).

On the other hand, recalling (5.1.2), the map

$$
\begin{equation*}
H \rightarrow\left(\mathbb{R}^{*}\right)^{l}, \quad h \rightarrow\left(h^{\nu_{1}}, \ldots, h^{\nu_{l}}\right) \tag{5.5.24}
\end{equation*}
$$

is a diffeomorphism since the elements $\nu_{i}$ are a basis of $\hbar^{\prime}$. Thus since we can put $\lambda=\nu_{i}, i=1, \ldots, l$, this proves $\lim _{t \rightarrow \infty} c(t)=d(w)$ and $\lim _{t \rightarrow \infty} \tilde{c}(t)=$ $d(w)^{-\kappa}$.
Q.E.D.

Related to the two maps $H \rightarrow H, g_{o} \mapsto h(g) g_{o}^{-1}$ and $g_{o} \mapsto h(g) g_{o}^{-\kappa}$ mentioned in Theorem 5.5 one has

Proposition 5.5.3. Let $\lambda \in D$ and let the notation be as in Theorem 5.5. Then

$$
\begin{equation*}
\left(h(g) g_{0}{ }^{\kappa}\right)^{\lambda}=\left\{v^{\lambda}, n(g) v^{\kappa \lambda}\right\} \tag{5.5.25}
\end{equation*}
$$

Furthermore there exists a scalar $b \in \mathbb{R}$ such that $s(\kappa)^{2} v_{\lambda}=b v_{\lambda}$ and one has

$$
\begin{equation*}
\left(h(g) g_{o}^{-1}\right)^{-\lambda}=b\left\{\bar{n}(g) v^{\lambda}, v^{\kappa \lambda}\right\} \tag{5.5.26}
\end{equation*}
$$

Proof. By (5.3.6) one has $\bar{n} h n=s_{o}(\kappa)^{-1} \bar{n}_{-f} g_{o} \bar{n}_{f}^{-1}$, where we have written $\bar{n}, h$, and $n$ for $\bar{n}(g), h(g)$, and $n(g)$. Also, we have written $\bar{n}_{-f}$ and $\bar{n}_{f}^{-1}$ for $\bar{n}_{-f}(w)$ and $\bar{n}_{f}^{-1}(w)$. But then since $\tilde{n}^{*} \in N$ and $h=h^{*}$ one has

$$
\begin{align*}
h^{\lambda}\left\{v^{\lambda}, n v^{\kappa \lambda}\right\} & =\left\{v^{\lambda}, \bar{n} h n v^{\kappa \lambda}\right\}  \tag{5.5.27}\\
& =\left\{v^{\lambda}, s_{o}(\kappa)^{-1} \bar{n}_{-f} g_{0} \bar{n}_{f}^{-1} v^{k \lambda}\right\} .
\end{align*}
$$

But clearly $\bar{n}_{-f} g_{o} \bar{n}_{f}^{-1} v^{\kappa \lambda}=g_{o}^{\kappa \lambda} v^{\kappa \lambda}$. Thus since $\left(s_{o}(\kappa)^{-1}\right)^{*}=s_{o}(\kappa)$ by Lemma 5.2.3 one has $h^{\lambda}\left\{v^{\lambda}, n v^{\kappa \lambda \lambda}\right\}=g_{o}^{\kappa \lambda}\left\{v^{\kappa \lambda}, v^{\kappa \lambda}\right\}$ or $\left\{v^{\lambda}, n v^{\kappa \lambda}\right\}=h^{-\lambda} g_{o}^{-\kappa \lambda}$, proving (5.5.25). But now on the other hand $h^{\lambda}\left\{\tilde{n} v^{\lambda}, v^{\kappa \lambda}\right\}=\left\{\bar{n} h n v^{\lambda}, v^{\kappa \lambda}\right\}=\left\{s_{o}(\kappa)^{-1} \bar{n}_{f} g_{o} \bar{n}_{f}^{-1} v^{\lambda}, v^{\kappa \lambda}\right\}$. But one has $b \in \mathbb{R}$ such that $s_{0}(\kappa) v^{\kappa \lambda}=s_{0}(\kappa)^{2} v^{\lambda}=b v^{\lambda}$ since of course $\kappa^{2}=1$. But $\bar{n}_{-f}^{*}$ and $\left(\bar{n}_{f}^{-1}\right)^{*}$ are in $N$ and hence they fix $v^{\lambda}$. Also $g_{o}^{*}=g_{o}$. Thus $h^{\lambda}\left\{\tilde{n} v^{\lambda}, v^{\kappa \lambda}\right\}=b g_{o}{ }^{\lambda}$. This proves (5.5.26), since $b= \pm 1$ by Lemma 5.2.3. Q.E.D.

For any $\mu \in \Lambda$ let

$$
\begin{equation*}
o(\mu)=\left\langle\mu, x_{\theta}\right\rangle \tag{5.5.28}
\end{equation*}
$$

where $\boldsymbol{x}_{o} \in \mathscr{K}$ is defined as in Section 2.1.
Remark 5.5. Since $\lambda-\kappa \lambda$ is spanned by roots (using, e.g., (5.1.9)) one has $o(\lambda-\kappa \lambda) \in \mathbb{Z}$. We shall not use the following, but one can show that the scalar $b$ in Proposition 5.5.3 is given by $b=(-1)^{o(\lambda-\kappa \lambda)}$.
5.6. In this section we wish to show the connection between the development here and the results and calculations of Moser in [19]. Let $x \in \mathrm{Jac} / p$ so that $x$ is a "symmetric" Jacobi element. By Theorem 5.4 there exist unique $w_{o} \in h_{+}$and $g_{o} \in H$ such that $x=k\left(g_{o}, w_{o}\right) w_{o}$. Let $g$ and $w$ be defined as usual so that $w=f+w_{o}$ and $\rho_{w}^{-1}\left(g_{o}\right) \in G_{0}{ }^{w}$. Now for any $\lambda \in D$ let, for $1 \leqslant i \leqslant$ $q(\lambda)=\operatorname{dim} V^{\lambda}$,

$$
\begin{equation*}
r_{i}{ }^{\lambda}\left(g_{o}, w_{o}\right)=\left\{k v_{i}^{\lambda}, v^{k \lambda}\right\}, \tag{5.6.1}
\end{equation*}
$$

where $k=k\left(g_{o}, w_{o}\right)$. We use the letter $r$ here because of the connection, as will soon be seen, between $r_{i}{ }^{\lambda}\left(g_{o}, w_{o}\right)$ and the coordinates $r_{i}$ in [19]. Now as noted in Remark 5.4 the eigenvectors for $\pi_{\lambda}(x)$ are just $k v_{i}{ }^{\lambda}$. The corresponding eigenvalues are clearly $\zeta_{i}=\left\langle\lambda_{i}, w_{o}\right\rangle$. See Section 5.5. Thus if $\zeta$ is indeterminate one has (as an entry in the resolvent)

$$
\begin{equation*}
\left\{(\zeta 1-x)^{-1} v^{\kappa \lambda}, v^{\kappa \lambda}\right\}=\sum_{i=1}^{a(\lambda)} \frac{\left(r_{i}^{\lambda}\right)^{2}}{\zeta-\zeta_{i}}, \tag{5.6.2}
\end{equation*}
$$

where 1 is the identity operator on $V^{\lambda}$ and $r_{i}{ }^{\lambda}=r_{i}{ }^{\lambda}\left(g_{o}, w_{o}\right)$. Equation (5.6.2) follows immediately from the relation

$$
\begin{equation*}
\left(\zeta 1-\pi_{\lambda}(x)\right)^{-1}=\pi_{\lambda}(k)\left(\zeta 1-\pi_{\lambda}\left(w_{o}\right)\right)^{-1} \pi_{\lambda}\left(k^{-1}\right) \tag{5.6.3}
\end{equation*}
$$

The $r_{i}{ }^{\lambda}$ may be given by
Proposition 5.6. One has

$$
\begin{equation*}
\left(r_{i}^{\lambda}\left(g_{o}, w_{o}\right)\right)^{2}=h(g)^{\kappa \lambda} g_{o}^{-\lambda_{i}} b_{i}(\lambda, w) \tag{5.6.4}
\end{equation*}
$$

where we recall $b_{i}(\lambda, w)=\left\{\bar{n}_{f}(w)^{-1} v^{\lambda}, v_{i}^{\lambda}\right\}\left\{\bar{n}_{-f}(w) v^{\lambda_{i}}, v^{k \lambda \lambda}\right\}$. See (5.5.12).
Proof. If we put $h=h(g), n=n(g), \bar{n}_{f}=\bar{n}_{f}(w)$, and $p=p\left(g_{o}, w_{o}\right)$ then recalling (5.4.1) one has $k=h^{1 / 2} n \bar{n}_{f} p^{-1}$. But $r_{i}^{\lambda}=\left\{k v_{i}^{\lambda}, v^{\kappa \lambda}\right\}$. Thus since $\left(h^{1 / 2}\right)^{*}=h^{1 / 2}$ and $n^{*} \in \bar{N}$ so that $n^{*} v^{\kappa \lambda}=v^{\kappa \lambda}$ one has

$$
\begin{equation*}
r_{i}^{\lambda}=h^{\kappa \lambda / 2} p^{-\lambda_{i}}\left\{\bar{n}_{f}(w) v_{i}^{\lambda}, v^{\kappa \lambda}\right\} \tag{5.6.5}
\end{equation*}
$$

Now let $m \in H_{\mathbb{C}}$ be as in (3.5.5) and let $m_{0} \in H_{\mathbb{C}^{8}}$ be such that Ad $m_{o}=m$. But then $m_{o} \bar{n}_{f}(w)\left(m_{o}\right)^{-1}=\bar{n}_{-f}(w)$ in $G_{\mathbb{C}^{s}}$ using (5.2.2), recalling also that $m^{2}=1$. But then if $c^{\prime}, c^{\prime \prime} \in \mathbb{C}^{*}$ are defined by $m_{o}^{-1} v_{i}^{\lambda}=c^{\prime} v_{i}^{\lambda}$ and $m_{o}^{*} v^{\kappa \lambda}=c^{\prime \prime} v^{\kappa \lambda}$ one has $\left\{\bar{n}_{-f}(w) v_{i}^{\lambda}, v^{\kappa \lambda}\right\} c=\left\{\bar{n}_{f}(w) v_{i}^{\lambda}, v^{\kappa \lambda}\right\}$, where $c=c^{\prime} c^{\bar{n}}$. But $|c|=1$ since $m_{0}$ and $m_{o}^{*}$ clearly have finite order. But then one must have

$$
\begin{equation*}
\left(\left\{\bar{n}_{-f}(w) v_{i}^{\lambda}, v^{\kappa \lambda}\right\}\right)^{2}=\left\{\bar{n}_{f}(w) v_{i}^{\lambda}, v^{\kappa \lambda}\right\}^{2} \tag{5.6.6}
\end{equation*}
$$

since $c^{2}$ times the left side of (5.6.6) is the right side of (5.6.6), both sides are nonnegative, and $\left|c^{2}\right|=1$. (That is, if (5.6.6) is nonzero, one has $c^{2}=1$.)

But now if $d=d(w)$, recalling (5.5.2) one has $d v_{i}{ }^{\lambda}=d^{\lambda} v_{i}{ }^{\lambda}$ and hence $d^{\lambda_{i}}\left\{v_{i}{ }^{\lambda},\left(\bar{n}_{f}(w)\right)^{-1} v^{\lambda}\right\}=\left\{\bar{n}_{-f}(w) v_{i}{ }^{\lambda}, v^{\kappa \lambda}\right\}$. But then upon substituting for just one of the factors in the left side of $(5.6 .6)$ one has

$$
\begin{equation*}
d^{\lambda_{i}} b_{i}(\lambda, w)=\left\{\bar{n}_{f}(w) v_{i}^{\lambda}, v^{\kappa \lambda}\right\}^{2} \tag{5.6.7}
\end{equation*}
$$

But then by squaring (5.6.5) and substituting (5.6.7) one has

$$
\begin{equation*}
\left(r_{i}^{\lambda}\right)^{2}=h^{\kappa \lambda} p^{-2 \lambda_{i}} d^{\lambda_{i}} b_{i}(\lambda, w) . \tag{5.6.8}
\end{equation*}
$$

But $p^{-2} d=g_{o}^{-1}$. This proves (5.6.4).
Q.E.D.

Now clearly one has

$$
\begin{equation*}
\sum_{i=1}^{\alpha(\lambda)}\left(r_{i}^{\lambda}\right)^{2}=1 \tag{5.6.9}
\end{equation*}
$$

by (5.2.11) and Lemma 5.2.3.

In [19] Moser considers the case where for us $G^{s}=S l(n, \mathbb{R})$ and $\pi_{\lambda}$ is thc standard representation. That is, $V^{\lambda}=\mathbb{R}^{n}$. In that case if we write $\boldsymbol{r}_{\boldsymbol{i}}=\boldsymbol{r}_{\boldsymbol{i}}{ }^{\boldsymbol{\lambda}}$ then $r_{i}{ }^{2}$ has the same meaning as in [19]. He then changes the definition of $r_{i}{ }^{2}$ so that they become "homogeneous." That is, in our notation define $s_{i}^{\lambda}\left(g_{o}, w_{o}\right)$ by putting

$$
\begin{equation*}
s_{i}^{\lambda}\left(g_{o}, w_{o}\right)=h(g)^{-\kappa \lambda / 2} r_{i}^{\lambda}\left(g_{o}, w_{o}\right) \tag{5.6.10}
\end{equation*}
$$

so that by Proposition 5.6 one has, if $s_{i}{ }^{\lambda}=s_{i}{ }^{\lambda}\left(g_{o}, w_{o}\right)$,

$$
\begin{equation*}
\left(s_{i}^{\lambda}\right)^{2}=g_{o}^{-\lambda_{i}} b_{i}(\lambda, w) . \tag{5.6.11}
\end{equation*}
$$

But then by (5.6.9) Eq. (5.6.2) becomes (see 3.9 in [19])

$$
\begin{equation*}
\left\{(\zeta 1-x)^{-1} v^{\kappa \lambda}, v^{\kappa \lambda \lambda}\right\}=\left(\sum_{i} \frac{\left(s_{i}^{\lambda}\right)^{2}}{\zeta-\zeta_{i}}\right) \frac{1}{\sum_{i}\left(s_{i}^{\lambda}\right)^{2}} . \tag{5.6.12}
\end{equation*}
$$

Now the problem of solving the generalized Toda lattice, as we shall see in Section 7, rests with determining $h(g \exp t w)^{\nu}$ for suitable $\nu \in h$. For the standard Toda lattice one needs only the case $\nu=\alpha_{i}, i=1,2, \ldots, l$. However, if $x(t)=k\left(g_{o} \exp t w_{o}\right) w_{o}$ then the functions $h(g \exp t w)^{-\alpha_{i} / 2}$ appear as matrix entries of $x(t)$. See (3.7.4) and (5.4.10). Furthermore as Moser notes, going back to Stieltjes, the left side of (5.6.12) with $x(t)$ replacing $x$ is computed using a continued fraction expansion in the matrix entries. One notes that $v^{\kappa \lambda}=e_{n}$, using the notation of [19]. On the other hand (5.6.12) clearly becomes an equation for all $t \in \mathbb{R}$ if we substitute $\left(s_{i}{ }^{\lambda}\right)^{2} e^{-t\left\langle\lambda_{i}, w_{o}\right\rangle}$ for $\left(s_{i}^{\lambda}\right)^{2}$. These being simple exponential functions the method of [19] is then to use (5.6.12) to inductively solve for $h(g \exp t w)^{x_{i}}$.
Now it is clear that a more direct and systematic approach to determine $h(g \exp t w)^{\lambda}$, for $\lambda \in D$ arises from (5.5.21). However, (5.5.21) is a formula for $h(g \exp t w)^{\lambda}$ only insofar as one can determine $\bar{n}_{f}(w)^{-1}$ and $\bar{n}_{-f} f(w)$ (and their images under $\pi_{\lambda}$ ) explicitly. Furthermore the scattering phase change also will depend upon such a determination. We will solve this problem in Section 5.8 using the machinery of representation theory.
5.7. Let $\mathbb{R}\left[G^{s}\right]$ be the group algebra of $G^{s}$ when the latter is regarded as an abstract group. Let $D\left(G^{s}\right)=\mathbb{R}\left[G^{s}\right] \otimes_{\mathbb{R}} U$. In this section only we will write $a \cdot u$ for $a u$, where $a \in G^{s}, u \in U$, to denote the adjoint action of $a$ on $u$. This is done to avoid confusion with the multiplication in $D\left(G^{s}\right)$, which will now be defined. One makes $D\left(G^{s}\right)$ into an algebra by retaining the given algebra structures in $\mathbb{R}\left[G^{s}\right]$ and $U$ and putting $a n a^{-1}=a \cdot u$ for $a \in G^{s}, u \in U$. One knows that we may identify $D\left(G^{s}\right)$ with the algebra (under convolution) of all distributions (in the sense of Schwartz) of finite support on $G^{s}$. Thus if $a \in G^{s}$ then $a U=U a$ is the set of all distributions on $G^{s}$ with support at $a$.

Now if $V$ is a finite-dimensional real vector space then we will say that $V$ is a smooth $D\left(G^{s}\right)$ module if it is a $D\left(G^{s}\right)$ module with respect to a representation $\pi: D\left(G^{s}\right) \rightarrow$ End $V$ such that $\pi \mid G^{s}$ is a representation of Lie groups and $\pi \mid g$ is the differential of $\pi \mid G^{s}$. Now a linear functional $\psi$ on $D\left(G^{s}\right)$ will be called a representative functional if it is in the span of linear functionals of the form $\psi_{v . v^{\prime}}$, where $V$ is a smooth $D\left(G^{s}\right)$ module, $v \in V, v^{\prime} \in V^{\prime}$ (the dual to $V$ ), and $\psi_{v, v^{\prime}}(a)=\left\langle a v, v^{\prime}\right\rangle$ for all $a \in D\left(G^{s}\right)$. Let $D\left(G^{s}\right)^{\prime}$ be the space of all representative functionals on $D\left(G^{s}\right)$. Then as one knows, $D\left(G^{s}\right)$ and $D\left(G^{s}\right)^{\prime}$ are non-singularly paired. (This is clear since (1) $\subseteq G^{s} \subseteq G_{\mathbb{C}}{ }^{s}$ and $G_{\mathbb{C}^{s}}$ has a faithful finite-dimensional representation and (2) $D\left(G^{s}\right)^{\prime}$ has the structure of an algebra-using the coalgebra structure on $D\left(G^{s}\right)$. Given distinct elements $g_{i} \in G^{s}, i=1, \ldots, k$, one easily then constructs $\psi \in D\left(G^{s}\right)^{\prime}$, which vanishes to any given preassigned order at $g_{2}, \ldots, g_{l}$ and is such that $\psi(u) \neq 0$, where $u$ is a given distribution with support at $g_{1}$.)

We now topologize $D\left(G^{s}\right)$ (the weak * topology) so that $D\left(G^{s}\right)^{\prime}$ is its continuous dual. One notes then, for example, that the series $\sum_{j=0}^{\infty}\left(x^{j} / j\right.$ !) converges to $\exp ^{s} x$ for any $x \in g$. Let $\hat{D}\left(G^{s}\right)$ be the completion of $D\left(G^{s}\right)$ with respect to this topology. If $a \in D\left(G^{s}\right)$ then left and right multiplication by $a$ in $D\left(G^{s}\right)$ is clearly continuous and hence such an operator extends to $\hat{D}\left(G^{s}\right)$. This defines on $\hat{D}\left(G^{s}\right)$ the structure of a two-sided $D\left(G^{s}\right)$-module.

Now let $\mathscr{S}$ be the set of all finite sequences

$$
\begin{equation*}
s=\left(i_{1}, \ldots, i_{k}\right), \tag{5.7.1}
\end{equation*}
$$

where $k \geqslant 0$ is arbitrary and $i_{j}$ is an integer such that $1 \leqslant i_{j} \leqslant l$ for $j \leqslant k$. We will write $|s|$ for the length $k$ of the element $s \in \mathscr{S}$.

Now let $U(\bar{x})$ and $U(n)$ be respectively the enveloping algebras of $\bar{n}$ and $n$. For convenience write $e_{i}=e_{\alpha_{i}}$ and $e_{-i}=e_{-\alpha_{i}}$ for $i=1, \ldots, l$. Furthermore if $s \in \mathscr{S}$ and $s$ is given by (5.7.1) put

$$
\begin{equation*}
e_{-s}=e_{-i_{k}} \cdots e_{-i_{2}} e_{-i_{1}} \tag{5.7.2}
\end{equation*}
$$

and

$$
\begin{equation*}
e_{s}=e_{i_{1}} e_{i_{2}} \cdots e_{i_{k}} \tag{5.7.3}
\end{equation*}
$$

If $|s|=0$ then one puts $e_{s}=e_{-s}=1$. In any case

$$
\begin{equation*}
e_{s}^{*}=e_{-s} \tag{5.7.4}
\end{equation*}
$$

and since the $e_{i}$ (resp. $e_{-i}$ ) generate $n$ (resp. $\bar{n}$ ) it is clear that the $e_{s}$ (resp. $e_{-s}$ ) for $s \in \mathscr{S}$ span $U(x)$ (resp. $U(\bar{x})$ ). It should be noted, however, that the $e_{-s}$ (for example) are not linearly independent.

Now let $x_{0} \in \ell$ be as in Section 2.1 so that $\left[e_{-i}, x_{0}\right]=e_{-i}$. One thus has

$$
\begin{equation*}
\left[e_{-s}, x_{o}\right]=|s| e_{-s} \tag{5.7.5}
\end{equation*}
$$

for any $s \in \mathscr{S}$. Now let $U(\bar{n})_{k}-\left\{u \in U(\bar{n}) \mid\left[u, x_{o}\right] \cdots k u\right\}$ so that one has the direct sum

$$
\begin{equation*}
U(\bar{x})=\sum_{k=0}^{\infty} U(\bar{x})_{k} \tag{5.7.6}
\end{equation*}
$$

and $U(\bar{n})_{k}$ is spanned by all $e_{-s}, s \in \mathscr{S}$, where $|s|=k$.
Now let $D(\bar{N})$ be the subalgebra of $D\left(G^{s}\right)$ generated by $\bar{N}$ and $\bar{n}$ and let $\hat{D}(\bar{N})$ be the closure of $D(\bar{N})$ in $D\left(G^{s}\right)$.

Proposition 5.7.1. For every $k \in \mathbb{Z}_{+}$choose an arbitrary element $u_{k} \in(U(\bar{n}))_{k}$. Then the infinite sum $\sum_{k \in \mathbb{Z}_{+}} u_{k}$ converges in $\hat{D}(\bar{N})$. Furthermore any $u \in \hat{D}(\bar{N})$ can be uniquely written as an infinite sum

$$
\begin{equation*}
u=\sum_{k \in \mathbb{Z}_{+}} u_{k} \tag{5.7.7}
\end{equation*}
$$

where $u_{k} \in U(\bar{x})_{k}$.
Proof. Let $U_{+}(\bar{n})=\sum_{k=1}^{\infty} U(\bar{n})_{k}$ so that $U_{+}(\bar{n})$ is the augmentation ideal. Furthermore it is clear that the elements $e_{-i}, i=1, \ldots, l$, are a basis of $U(\bar{n})_{1}$ so that if $j \in \mathbb{Z}_{+}$and $U_{+}(\bar{n})^{j}$ is the $j$ th power of the augmentation ideal one has

$$
\begin{equation*}
U_{+}(\bar{n})^{j}=\sum_{k=j}^{\infty} U(\bar{n})_{k} \tag{5.7.8}
\end{equation*}
$$

Now let $D(\bar{N})^{\prime}$ be the set of all restrictions $\psi \mid D(\bar{N})$, where $\psi \in D\left(G^{s}\right)^{\prime}$. But if $\pi$ is any faithful finite-dimensional holomorphic representation of $G_{\mathbb{C}^{s}}$ then $\pi(\bar{N})$ is a group of unipotent operators and $\pi \mid \bar{N}_{\mathbb{C}}$ is faithful. It follows easily then that if $\psi \in D(\bar{N})^{\prime}$ then the restriction $\psi \mid U(\bar{n})$ vanishes on a power of $U_{+}(\bar{x})$ and that every such linear functional on $U(\bar{x})$ is uniquely of this form. Thus by (5.7.6) and (5.7.8) if $U(\bar{x})_{k}^{\prime}$ is the dual space to $U(\bar{n})_{k}$ we may regard $U(\bar{n})_{k}^{\prime} \subseteq D(\bar{N})^{\prime}$, where $U(\bar{x})_{k}^{\prime}$ is identified with the set of all elements in $D(\bar{N})^{\prime}$ which vanish on $D(\bar{r})_{j}$ for $j \neq k$. One then has the direct sum

$$
\begin{equation*}
D(\bar{N})^{\prime}=\sum_{k \in \mathbb{Z}_{+}} U(\bar{n})_{k}^{\prime} \tag{5.7.9}
\end{equation*}
$$

Now let $y_{k} \in U(\bar{n})_{k}$ for all $k \in \mathbb{Z}_{+}$. Then for any $\psi \in D(\bar{N})^{\prime}$ one has $\psi\left(y_{k}\right)=0$ for $k$ sufficiently large. This proves that the sum $\sum y_{k}$ converges. On the other hand if $\bar{n} \in \bar{N}$ then we may write $\bar{n}=\exp ^{s} x$ for $x \in \bar{n}$. Thus $\bar{n}=\sum\left(x^{j} / j!\right)$. This (and in fact also (5.7.9)) implies that $U(\bar{x})$ is dense in $D(\bar{N})$. Now let $u \in \hat{D}(\bar{N})$ and let $u(j) \in U(\bar{x}), j=1,2, \ldots$, be a sequence such that $u(j)$ converges to $u$. By (5.7.6) we may write $u(j)=\sum u(j)_{k}$, where $u(j)_{k} \in U(\bar{n})_{k}$. But now $u(j)_{k}$, as $j \rightarrow \infty$, must, by (5.7.9), have a limit $u_{k} \in U(\bar{n})_{k}$ and also $u=\sum u_{k}$. This is clear since $\psi(u(j))=\psi\left(u(j)_{k}\right)$ converges for any $\psi \in U(\bar{n})_{k}^{\prime}$. This proves
the existence in Proposition 5.7. The uniqueness follows easily since $\boldsymbol{u}_{\boldsymbol{k}}$ is characterized by the relation $\psi\left(u_{k}\right)=\psi(u)$ (which is of course defined) for any $\psi \in U(\bar{n})_{k}^{\prime}$.
Q.E.D.

Proposition 5.7.2. For any $s \in \mathscr{S}$ choose $a$ scalar $b_{s} \in \mathbb{R}$. Then the infinite sum $\sum b_{s} e_{-8}$ relative to any simple ordering in $\mathscr{S}$ converges to an element in $\hat{D}(\bar{N})$ and this element is independent of the ordering. Furthermore any element in $\bar{D}(\bar{N})$ may be written as such a sum.

Proof. Given any $\psi \in D(\bar{N})^{\prime}$ it suffices in order to prove the first statement only to note that $\psi\left(e_{-s}\right)=0$ for all $s \in \mathscr{S}$ such that $|s|$ is sufficiently large. But this is clear from (5.7.9) since $e_{-s} \in U(\bar{n})_{|s|}$. See (5.7.5). The second statement follows from Proposition 5.7 .1 since $U(\bar{x})_{k}$ is spanned by all $e_{-s}$ such that $|s|=k$.
Q.E.D.
5.8. Now let $s \in \mathscr{S}$. Assume $s$ is given by (5.7.1) so that in particular $|s|=k$. Now if $k \geqslant 1$ let $\varphi(s) \in \Lambda$ be the linear form on the Cartan subalgebra $h$ given by

$$
\begin{equation*}
\varphi(s)=\sum_{j=1}^{|s|} \alpha_{i_{j}} \tag{5.8.1}
\end{equation*}
$$

One puts $\varphi(s)$ equal to the constant function 1 on $h$ if $|s|=0$.
On the other hand for any $0 \leqslant j \leqslant|s|$ let $s_{j} \in \mathscr{S}$ be the sequence obtained from $s$ by "cutting off" the last $j$ terms. That is,

$$
\begin{equation*}
s_{j}=\left(i_{1}, i_{2}, \ldots, i_{|s|-j}\right) \tag{5.8.2}
\end{equation*}
$$

One notes of course that $s_{0}=s$ and

$$
\begin{equation*}
\left|s_{j}\right|=|s|-j \tag{5.8.3}
\end{equation*}
$$

We now observe that $s$ defines a polynomial of degree $|s|$ on $\hbar$ by putting for any $w_{o} \in h$

$$
\begin{equation*}
p\left(s, w_{o}\right)=\prod_{j=0}^{|s|}\left\langle\varphi\left(s_{j}\right), w_{o}\right\rangle \tag{5.8.4}
\end{equation*}
$$

Remark 5.8.1. It is useful to think of $p\left(s, w_{o}\right)$ as some sort of generalized "factorial" expression. Indeed putting $w_{0}=x_{0}$ (see Section 2.1) note that $p\left(s, x_{o}\right)=|s|!$

Proposition 5.8.1. If $w_{o} \in h_{+}$then $p\left(s, w_{o}\right)>0$ for any $s \in \mathscr{S}$.

Proof. This is immediate from (5.8.1) and (5.8.2) since $\left\langle\alpha_{i}, w_{o}\right\rangle>0$, $i=1, \ldots, l$, for $w_{o} \in h_{+}$by definition of $h_{+}$. Q.E.D.

Now for any $s \in \mathscr{S}$ let $\bar{s} \in \mathscr{S}$ be the sequence obtained from $s$ by reversing the order. Thus if $s$ is given by (5.7.1) then

$$
\begin{equation*}
\bar{s}=\left(i_{k}, i_{k-1}, \ldots, i_{1}\right) \tag{5.8.5}
\end{equation*}
$$

Now recalling Lemma 3.5.2 we can determine the element $\bar{n}_{f}(w) \in \bar{N}$ and observe the dependence on $w_{o} \in h_{+}$.

Proposition 5.8.2. Let $w_{o} \in h_{+}$and let $w=f+w_{o}$, where $f=\sum_{i=1}^{l} e_{-\alpha_{i}}$. Let $\bar{n}_{f}(w) \in \bar{N}$ be the element defined by (3.5.2) so that $\bar{n}_{f}(w) \cdot w_{o}=w$. Now let $\mathscr{S}$ be the set of all finite sequences $s=\left(i_{1}, \ldots, i_{k}\right)$ of integers where $1 \leqslant i_{j} \leqslant l$. For any $s \in \mathscr{S}$ let $p\left(s, w_{o}\right)$ be defined by (5.8.4) so that $p\left(s, w_{o}\right)>0$ (recalling Proposition 5.8.1). Let $e_{-s} \in U(\bar{n})$ be defined by (5.7.2). Then recalling (5.2.2) (and also Proposition 5.7.2) one has

$$
\begin{equation*}
\bar{n}_{-f}(w)=\sum_{s \in \mathscr{S}}(-1)^{|s|} \frac{e_{-s}}{p\left(s, w_{o}\right)} . \tag{5.8.6}
\end{equation*}
$$

Furthermore recalling (5.8.5) one has

$$
\begin{equation*}
\bar{n}_{f}(w)^{-1}=\sum_{s \in \mathscr{P}}(-1)^{|s|} \frac{e_{-\bar{s}}}{p\left(s, w_{o}\right)} . \tag{5.8.7}
\end{equation*}
$$

Proof. Now let $s \in S$. Then if $\left\{t^{1}, \ldots, t^{l}\right\}$ is the set of all $t \in S$ such that $t_{1}=s$ (using the notation of (5.8.2)) one clearly has

$$
\begin{equation*}
f e_{-s}=\sum_{i=1}^{l} e_{-t^{i}} \tag{5.8.8}
\end{equation*}
$$

On the other hand one also has

$$
\begin{equation*}
e_{-\bar{s}} f=\sum_{i=1}^{l} e_{-\bar{t}^{i}} \tag{5.8.9}
\end{equation*}
$$

Now let $v \in D(\bar{N})$ be the element (see Proposition 5.7.2) defined by the right side of (5.8.6). Then by (5.8.8) one notes that

$$
\begin{equation*}
-f v=\sum_{\substack{s \in \mathscr{S} \\|s| \geqslant 1}}(-1)^{|s|} \frac{e_{-s}}{p\left(s_{1}, w_{o}\right)} \tag{5.8.10}
\end{equation*}
$$

On the other hand since $\left[e_{-i}, w_{0}\right]=\left\langle\alpha_{i}, w_{a}\right\rangle e_{-i}$ one certainly has

$$
\begin{equation*}
\left[e_{-s}, w_{o}\right]=\left\langle\varphi(s), w_{o}\right\rangle e_{-s} \tag{5.8.11}
\end{equation*}
$$

However, note that

$$
\begin{equation*}
\frac{\left\langle\varphi(s), w_{o}\right\rangle}{p\left(s, w_{o}\right)}=\frac{1}{p\left(s_{1}, w_{o}\right)} . \tag{5.8.12}
\end{equation*}
$$

But then by (5.8.10) one has in $\hat{D}\left(G^{s}\right)$,

$$
\begin{equation*}
\left[v, w_{o}\right]=\sum_{\substack{s \in \mathscr{C} \\|s| \geqslant 1}}(-1)^{|s|} \frac{e_{-s}}{p\left(s_{1}, w_{o}\right)} . \tag{5.8.13}
\end{equation*}
$$

The term for $|s|=0$ does not occur in (5.8.13) since $\left[1, w_{o}\right]=0$. Thus by (5.8.10) one has (in $\hat{D}\left(G^{s}\right)$ )

$$
\begin{equation*}
\left[v, w_{o}\right]=-f v . \tag{5.8.14}
\end{equation*}
$$

But in $\hat{D}\left(G^{s}\right)$ one has the relation (see the statement following (5.2.2)) $\bar{n}_{-f}(w) w_{0}\left(\bar{n}_{-f}(w)\right)^{-1}-w_{0}=-f$. Substituting then for $-f$ in (5.8.14) one has $\left[v, w_{o}\right]=\bar{n}_{-f}(w) w_{o}\left(\bar{n}_{-f}(w)\right)^{-1} v-w_{o} v$. Thus $v w_{o}=\bar{n}_{-f}(w) w_{o}\left(\bar{n}_{-f}(w)\right)^{-1} v$ and hence if we put $u=\vec{n}_{-f}(w)^{-1} v$ then $u w_{o}=w_{0} u$. That is, $u \in \hat{D}(\bar{N})$ and $u$ commutes with $w_{o}$. Now let $u_{k} \in U(\bar{n})_{k}$ be defined by (5.7.7) for any $k \in \mathbb{Z}_{+}$. Since $\left[w_{o}, u\right]=0$ and since $U(\bar{n})_{k}$ is stable under ad $w_{o}$ it follows from the uniqueness of (5.7.7) that $\left[w_{o}, u_{k}\right]=0$ for all $k$. However, since ad $w_{o}$ is diagonalizable with a strictly negative spectrum in $\bar{n}$ it follows that the same is true in $U(\bar{n})_{k}$ for any $k>0$. Thus $u_{k}=0$ for all $k>0$. On the other hand the leading term of $v$ is 1 , using (5.7.7), since $p\left(s, w_{o}\right)=1$ if $|s|=0$. By considering the exponential series the same is then true for $\left(\bar{n}_{-f}(w)\right)^{-1}$. Thus $u_{a}=1$. Indeed if $s, s^{\prime} \in \mathscr{S}$ then $e_{-s^{\prime}} e_{-s^{\prime}}=e_{-s^{\prime \prime}}$ for some $s^{\prime \prime} \in \mathscr{S}$ where $\left|s^{\prime \prime}\right|=$ $|s|+\left|s^{\prime}\right|$ so that using the expansions given by (5.7.7) one easily see that $u_{o}$ must be the product of the constant terms of $\vec{n}_{f}(w)^{-1}$ and $v$. Thus $u=1$ so that $\bar{n}_{-f}(w)=v$. This proves (5.8.6).

Now let $y$ be the right side of (5.8.7). Then by (5.8.9) one has

$$
\begin{equation*}
y f=-\sum_{\substack{s \in \mathscr{P} \\|s| \geqslant 1}}(-1)^{|s|} \frac{e_{-\bar{s}}}{p\left(s_{1}, w_{o}\right)} . \tag{5.8.15}
\end{equation*}
$$

On the other hand $\left[w_{o}, y\right]$ is also given by the right side of (5.8.15) using (5.8.11) and (5.8.12). Thus

$$
\begin{equation*}
\left[w_{o}, y\right]=y f \tag{5.8.16}
\end{equation*}
$$

But $\bar{n}_{f}(w) w_{o}\left(\bar{n}_{f}(w)\right)^{-1}-w_{o}=f$. Hence if we substitute for $f$ in (5.8.16) we obtain $y \bar{n}_{f}(w)=1$ in a manner similar to the proof of (5.8.6). That is, $y=$ $\left(\bar{n}_{f}(w)\right)^{-1}$.
Q.E.D.

Now let $D(N)$ be the subalgebra of $D\left(G^{s}\right)$ generated by $N$ and $n$ and let $\hat{D}(N)$ be the closure of $D(N)$ in $\hat{D}\left(G^{s}\right)$. It is obvious that Proposition 5.7.2 is valid if $N$ replaces $\bar{N}$ and $e_{s}$ replaces $e_{-s}$ for all $s \in \mathscr{S}$.

But now $\left(\bar{n}_{-f}(w)\right)^{*} \in N$ by (1.5.1) and (3.1.9).
Proposition 5.8.3. Let the notation be as in Proposition 5.8.2. See also (5.7.3). Then

$$
\begin{equation*}
\left(\bar{n}_{-f}(w)\right)^{*}=\sum_{s \in \mathscr{S}}(-1)^{|s|} \frac{e_{s}}{p\left(s, w_{n}\right)} \tag{5.8.17}
\end{equation*}
$$

Proof. Let $a \in G^{s}$ and $x_{j} \in U, j=1,2, \ldots$. Assume that $x_{j}$ converges to $a$ in $D\left(G^{s}\right)$. We then assert that $x_{j}^{*}$ converges to $a^{*}$. Indeed recalling the definition of representative functional in Section 5.7 and the particular representative functional $\psi_{v, v^{\prime}}$ one notes that $\psi_{v, v^{\prime}}\left(x_{j}^{*}\right)$ converges to $\psi_{v, v^{\prime}}\left(a^{*}\right)$ in case $V=V^{\lambda}$ by (5.1.11). But then the result follows easily for all $V$ using the complete reducibility of smooth $D\left(G^{s}\right)$ modules. Thus $x_{j}^{*}$ converges to $a^{*}$. But then (5.8.17) follows from (5.8.6) and (5.7.4).
Q.E.D.

Remark 5.8.2. One major advantage of formulas (5.8.6), (5.8.7), and (5.8.17) is that for any given smooth $D\left(G^{s}\right)$ module we need consider only a finite number of summands.
5.9. Now let $\lambda \in D$ and let

$$
\begin{equation*}
\mathscr{S}^{\lambda}=\{s \in \mathscr{S} \mid \varphi(s)=\lambda-\kappa \lambda\} . \tag{5.9.1}
\end{equation*}
$$

Now $\lambda-\kappa \lambda$ is the difference between the highest and lowest weights of $V^{\lambda}$ and hence we can write $\lambda-\kappa \lambda=\sum_{i=1}^{i} m_{i} \alpha_{i}$, where $m_{i} \in \mathbb{Z}_{+}$. If $o(\lambda-\kappa \lambda)$ is defined as in (5.5.28) note that $o(\lambda-\kappa \lambda)=\sum m_{i}$ and hence the cardinality of $\mathscr{S}^{\lambda}$ is given by

$$
\begin{equation*}
\left|\mathscr{S}^{\lambda}\right|=\frac{o(\lambda-\kappa \lambda)!}{m_{1}!\cdots m_{l}!} \tag{5.9.2}
\end{equation*}
$$

One notes also that $\mathscr{S}^{\lambda}$ is stable under the map $\mathscr{S} \rightarrow \mathscr{S}, s \rightarrow \bar{s}$.
Now for any $s \in \mathscr{S}^{\lambda}$ clearly $e_{-s}$ carries $v^{\lambda}$ into a multiple of $v^{k \lambda}$. Thus there exists $c_{s, \lambda} \in \mathbb{R}$ such that

$$
\begin{equation*}
e_{-s} v^{\lambda}=c_{s, \lambda} v^{\kappa \lambda} \tag{5.9.3}
\end{equation*}
$$

Proposition 5.9.1. Let $\lambda \in D$ and $w_{o} \in h_{+}$. Then if $p\left(s, w_{o}\right)$ is defined by (5.8.4) for any $s \in \mathscr{S}$ one has

$$
\begin{align*}
(-1)^{o(\lambda-\kappa \lambda)} \sum_{s \in \mathscr{S}_{\lambda}} \frac{c_{s, \lambda}}{p\left(s, w_{o}\right)} & =\left\{\bar{n}_{-f}(w) v^{\lambda}, v^{\kappa \lambda}\right\} \\
& =d(w)^{\lambda}, \tag{5.9.4}
\end{align*}
$$

where we recall that $d(w) \in H$ is defined by (5.5.2). Furthermore

$$
\begin{align*}
(-1)^{o(\lambda-\kappa \lambda)} \sum_{s \in \mathscr{P}_{\lambda}} \frac{c_{s, \lambda}}{p\left(\bar{s}, w_{o}\right)} & =\left\{\bar{n}_{f}(w)^{-1} v^{\lambda}, v^{\kappa \lambda}\right\}  \tag{5.9.5}\\
& =d(w)^{-\kappa \lambda}
\end{align*}
$$

In particular the numbers given by the sums in (5.9.4) and (5.9.5) are positive.
Proof. If $s \in \mathscr{P}$ then by (5.1.14) and (5.2.11) clearly $\left\{e_{-s} v^{\lambda}, v^{\kappa \lambda}\right\}=0$ if $s \notin \mathscr{S}^{\lambda}$ and $\left\{e_{-s} v^{\lambda}, v^{\kappa \lambda}\right\}=c_{s, \lambda}$ if $s \in \mathscr{S}$. But then (5.9.4) and (5.9.5) follow from (5.8.6), (5.8.7), (5.5.4), and (5.5.5). Note that we have interchanged $s$ and $\bar{s}$ in (5.8.7).
Q.E.D.

Since the numbers $c_{s, \lambda}, s \in \mathscr{S}^{\lambda}$ play such an important role in our solution of the generalized Toda lattice we will give another expression for them purely in terms of $U$. For any $s, t \in \mathscr{S}^{\lambda}$ the element $e_{t} e_{-s}$ clearly commutes with $h$. That is, in the notation of (5.1.10) one has $e_{t} e_{-s} \in U(0)$. Now by (5.1.15) there exists a unique element $u_{t, s} \in U(h)$ such that $e_{t} e_{-s}-u_{t, s}$ is in the left ideal $U_{n}$. On the other hand we can regard $U(h)$ as the algebra of real-valued polynomial functions on $\ell^{\prime}$ so that in particular if $u \in U(h)$ then $u v^{\lambda}=u(\lambda) v^{\lambda}$.

Proposition 5.9.2. Let $\lambda \in D$. Then the $o(\lambda-\kappa \lambda) \times o(\lambda-\kappa \lambda)$ matrix $\left\{u_{t, s}(\lambda)\right\}$, indexed by $t, s \in \mathscr{S}^{\lambda}$ has rank 1 and hence there exists a unique, up to sign, vector $\left\{c_{s, \lambda}\right\}$ in $\mathbb{R}^{o(\lambda-\kappa \lambda)}$ such that

$$
\begin{equation*}
c_{s, \lambda} c_{t, \lambda}=u_{t, s}(\lambda) \tag{5.9.6}
\end{equation*}
$$

for all $s, t$. The sign is, however, determined by the relation

$$
\begin{equation*}
(-1)^{o(\lambda-\kappa \lambda)} \sum_{s \in \mathscr{S}_{\lambda}} c_{s, \lambda}>0 \tag{5.9.7}
\end{equation*}
$$

Proof. Now clearly $\left\{e_{-s} v^{\lambda}, e_{-t} v^{\lambda}\right\}=\left\{e_{t} e_{-s} v^{\lambda}, v^{\lambda}\right\}=\left\{u_{s, t} v^{\lambda}, v^{\lambda}\right\}=u_{s, t}(\lambda)$ since $U_{n}$ annihilates $v^{\lambda}$. But $c_{s, \lambda} c_{t, \lambda}=\left\{e_{-s} v^{\lambda}, e_{-t} v\right\}$. Thus one has (5.9.6) and hence $\left\{u_{t, s}(\lambda)\right\}$ has rank 1 or 0 . Now by Remark 5.8 .1 choosing $w_{o}=x_{o}$ one has $p\left(s, x_{o}\right)=|s|!$ Thus, by (5.9.4), the left side of (5.9.7) is just $o(\lambda-\kappa \lambda)!d(x)^{\lambda}$, where $x=f+x_{o}$. This proves (5.9.7) and also that the rank of $\left\{u_{t, s}(\lambda)\right\}$ cannot be zero.
Q.E.D.
5.10. Now given $\lambda \in D$ let $\Delta^{\lambda} \subseteq \Lambda$ be the set of weights of $V^{\lambda}$. Now for any $\nu \in \Delta^{\lambda}$ let

$$
\begin{equation*}
\mathscr{S}^{\lambda}(\nu)=\left\{s \in \mathscr{S}^{\lambda} \mid \varphi\left(s_{o(\nu-\kappa \lambda)}\right)=\lambda-\nu\right\} \tag{5.10.1}
\end{equation*}
$$

where we use the notation of (5.8.1) and (5.8.2). That is, if $s \in \mathscr{P}^{\lambda}$ is given by (5.7.1) then $s \in \mathscr{S}^{\lambda}(\nu)$ in case $\sum_{j=1}^{\circ(\lambda-\nu)} \alpha_{i_{j}}=\lambda-\nu$.

Now for any $s \in \mathscr{S}$ and integer $j$, where $0 \leqslant j \leqslant|s|$, let $p_{j}\left(s, w_{o}\right)$ be the polynomial of degree $|s|$ in $w_{o}$ defined by putting

$$
\begin{equation*}
\left.p_{j}\left(s, w_{o}\right)=p\left(\overline{s_{|s|-j}}\right), w_{o}\right) p\left(\left(\left(\overline{\bar{s})_{j}}\right), w_{o}\right),\right. \tag{5.10.2}
\end{equation*}
$$

using the notation of (5.8.2) and (5.8.5). Explicitly if $s=\left(i_{1}, \ldots, i_{|s|}\right)$ then

$$
\begin{equation*}
p_{j}\left(s, w_{o}\right)=\left(\prod_{k=0}^{j}\left\langle\sum_{p=k+1}^{j} \alpha_{i_{p}}, w_{o}\right\rangle\right)\left(\prod_{m=j}^{|s|}\left\langle\sum_{q=j+1}^{m} \alpha_{i_{q}}, w_{o}\right\rangle\right) . \tag{5.10.3}
\end{equation*}
$$

It is clear of course that $p_{j}\left(s, w_{o}\right)>0$.
Remark 5.10. Recall Remark 5.8.1, where it was suggested that $p\left(s, w_{0}\right)$ should be regarded as a sort of factorial "function." In the same sense one should perhaps think of $p_{j}\left(s, w_{0}\right)$ in terms of binomial coefficients. Indeed if $w=x_{o}$, recalling Section 2.1, one easily has

$$
p_{j}\left(s, x_{o}\right)=j!(|s|-j)!
$$

For the two extreme cases one clearly has

$$
\begin{equation*}
p_{0}\left(s, w_{o}\right)=p\left(s, w_{o}\right) \tag{5.10.4}
\end{equation*}
$$

and

$$
\begin{equation*}
p_{|s|}\left(s, w_{o}\right)=p\left(\bar{s}, w_{o}\right) . \tag{5.10.5}
\end{equation*}
$$

Now for any $g_{o} \in H, w_{o} \in h_{+}, \lambda \in D$, and $t \in \mathbb{R}$ let

$$
\begin{equation*}
\Phi_{\lambda}\left(g_{0}, w_{o} ; t\right)=(-1)^{o(\lambda-\kappa \lambda)} \sum_{\nu \in \Delta^{\lambda}}\left(\sum_{s \in \mathscr{S}^{\circ} \lambda(\nu)} \frac{c_{s, \lambda}}{p_{o}(\lambda-\nu)\left(s, w_{o}\right)}\right) g_{o}^{\nu} e^{-t\left\langle\nu, w_{o}\right\rangle} \tag{5.10.6}
\end{equation*}
$$

where $c_{s, \lambda}$ is defined by (5.9.3), $\Delta^{\lambda}$ is the set of weights of $V^{\lambda}$ and $o(\mu)$ is defined by (5.5.28).

Theorem 5.10. Let $g$ be a real split semi-simple Lie algebra. Let $\hbar \subseteq g$ be a split Cartan subalgebra and let $h_{+} \subseteq \hbar$ be an open Weyl chamber. Let $G$ be the adjoint group and let $H \subseteq G$ be the subgroup corresponding to h. Let $g_{0} \in H, w_{o} \in h_{+}$ and let $\lambda$ be the highest weight of an irreducible representation of $g$ and for any $t \in \mathbb{R}$ let $\Phi_{\lambda}\left(g_{o}, w_{o} ; t\right)$ be defined by (5.10.6) where $c_{s, \lambda}$ is defined by (5.9.3) and $p_{j}\left(s, w_{o}\right)$ is defined by (5.10.3), recalling that $\alpha_{1}, \ldots, \alpha_{l}$ are the simple positive roots. Now let $w=\sum_{i=1}^{l} e_{-\alpha_{i}}+w_{o}$ for any fixed choice of negative simple root vectors $e_{-x_{d}}$. Let $G_{0}{ }^{w}$ be the connected component of the centralizer of $w$ in $G$
defined by (3.5.5)-(3.5.7) so that $g \exp t w \in G_{0}{ }^{w}$, where $g=\rho_{w}^{-1}\left(g_{o}\right)$ and $\rho_{w}$ is defined by (5.3.1). Now let $h(g \exp t w) \in H$ be defined by (2.6.2) using (2.4.5), (2.6.1), and (3.5.9). Then one has

$$
\begin{equation*}
h(g \exp (-t) w)^{\lambda}=\Phi_{\lambda}\left(g_{o}, w_{o} ; t\right) \tag{5.10.7}
\end{equation*}
$$

Proof. For the statement about the arbitrariness of the $e_{-\alpha_{i}}$ see Remark 1.5.1. Now for any $\nu \in \Delta^{\lambda}$ let $P_{\nu}: V^{\lambda} \rightarrow V^{\lambda}(\mu)$ be the orthogonal projection relative to $Q_{*}^{\lambda}$. Now if $b_{i}\left(\lambda, w_{o}\right)$ is defined by (5.5.12) then clearly

$$
\begin{equation*}
\left\{P_{\nu} \bar{n}_{f}(w)^{-1} v^{\lambda}, P_{\nu} \bar{n}_{-f}(w)^{*} v^{\kappa \lambda}\right\}=\sum_{\lambda_{i}=\nu} b_{i}\left(\lambda, w_{o}\right) \tag{5.10.8}
\end{equation*}
$$

and hence by $(5.5 .21)$ if $b_{\nu}\left(\lambda, w_{o}\right)$ is the left side of $(5.10 .8)$ then

$$
\begin{equation*}
h(g \exp t w)^{\lambda}=\sum_{\nu \in \Delta^{\lambda}} b_{v}\left(\lambda, w_{o}\right) g_{o}^{\nu} e^{t\left\langle\nu, w_{o}\right\rangle} \tag{5.10.9}
\end{equation*}
$$

Now for any $\mu \in \mathscr{R}^{\prime}$ let $\mathscr{S}^{(\mu)}=\{s \in \mathscr{S} \mid \varphi(s)=\mu\}$. Then by (5.8.7) and (5.8.17) for any $\nu \in \Delta^{\lambda}$ one has upon interchanging $s^{\prime}$ and $\bar{s}^{\prime}$

$$
\begin{equation*}
P_{\nu} \bar{n}_{f}(w)^{-1} v^{\lambda}=\sum_{s^{\prime} \in \mathcal{S}_{(\lambda-v)}}(-1)^{\left|s^{\prime}\right|} \frac{e_{-s^{\prime}, v^{\lambda}}}{p\left(\bar{s}^{\prime}, w_{o}\right)} \tag{5.10.10}
\end{equation*}
$$

and

$$
\begin{equation*}
P_{\nu} \bar{n}_{-f}(w)^{*} v^{\kappa \lambda}=\sum_{\left.s^{*} \in \mathscr{\mathscr { P } ^ { \prime } \nu - \kappa \lambda}\right)}(-1)^{\left|s^{*}\right|} \frac{e_{s^{\prime \prime} v^{\kappa \lambda}}}{p\left(s^{\prime \prime}, w_{o}\right)} \tag{5.10.11}
\end{equation*}
$$

Now composition of sequences clearly induces a bijection

$$
\begin{equation*}
\mathscr{S}^{(\lambda-\nu)} \times \mathscr{S}^{(\nu-\mu \lambda)} \rightarrow \mathscr{S}^{\lambda}(\nu), \quad\left(s^{\prime}, s^{\prime \prime}\right) \mapsto s \tag{5.10.12}
\end{equation*}
$$

where we note that $e_{-s}=e_{-s^{n}} e_{-s^{\prime}}$. One then has $\left\{e_{-s^{\prime} v^{\lambda}}, e_{s^{n}} v^{\kappa \lambda}\right\}=\left\{e_{-s} v^{\lambda}, v^{\kappa \lambda}\right\}=$ $c_{s, \lambda}$. Furthermore $\left.p\left(\overline{s^{\prime}}, w_{o}\right)=p\left(\overline{\left(\overline{s_{0}(\nu-\kappa \lambda)}\right)}\right), w_{o}\right)$ and $p\left(s^{\prime \prime}, w_{o}\right)=p\left(\overline{\left.(\bar{s})_{o(\lambda-\nu)}\right)}, w_{o}\right)$ so that

$$
\begin{equation*}
p\left(\bar{s}^{\prime}, w_{o}\right) p\left(s^{\prime \prime}, w_{o}\right)=p_{o(\lambda-p)}\left(s, w_{o}\right) . \tag{5.10.13}
\end{equation*}
$$

But also $\left|s^{\prime}\right|+\left|s^{\prime \prime}\right|=o(\lambda-\kappa \lambda)$ so that

$$
\begin{equation*}
b_{\nu}(\lambda, w)=(-1)^{o(\lambda-\nu \lambda)} \sum_{s \in \mathscr{S}^{\lambda_{(v)}}} \frac{c_{s, \lambda}}{p_{o(\lambda-\nu)}\left(s, w_{o}\right)} . \tag{5.10.14}
\end{equation*}
$$

But then (5.10.7) follows from (5.10.9) and (5.10.14).
Q.E.D.
5.11. Formula (5.10.7) expresses $h(g \exp (-t) w)^{\lambda}$ as a finite sum of
exponentials with explicit but clearly complicated coefficients. On the other hand the asymptotic values of $h(g \exp (-t) w)^{\lambda}$, or rather $\log h(g \exp (-t) w)^{\lambda}$, as $t \rightarrow \pm \infty$ may be given by simple linear formulas. We may also replace $\lambda \in D$ by any $\mu \in \ell^{\prime}$.

Theorem 5.11. Let the notation be as in Theorem 5.10 and let $\kappa$ be the Weyl group element which takes positive roots to negative roots. See Section 2.4. Let $\mu \in \ell^{\prime}$ and let $d(w) \in H$ be defined by (5.5.2). (See also (5.5.4) and (5.5.5).) Then if $R_{ \pm}\left(g_{o}, w_{0} ; t\right)$ are the real-valued remainder functions defined by

$$
\begin{aligned}
\log h(g \exp (-t) w)^{\mu} & =\left\langle-\kappa \mu, w_{o}\right\rangle t+\log \left(g_{o} d(w)^{-1}\right)^{\kappa \mu}+R_{+}\left(g_{o}, w_{o} ; t\right) \\
& -\left\langle-\mu, w_{o}\right\rangle t+\log \left(g_{o} d(w)\right)^{\mu}+R_{-}\left(g_{o}, w_{o} ; t\right)
\end{aligned}
$$

one has $\lim _{t \rightarrow+\infty} R_{+}\left(g_{o}, w_{o} ; t\right)=\lim _{t \rightarrow-\infty} R_{-}\left(g_{o}, w_{o} ; t\right)=0$.
Proof. By definition, one has with exponentiation

$$
\begin{equation*}
\frac{h(g \exp (-t) w)^{\mu}}{e^{t\left\langle-\kappa \mu, w_{o}\right)}\left(g_{0} d(w)^{-1}\right)^{\kappa \mu}}=e^{R_{+}\left(g_{0}, w_{0} ; t\right)} \tag{5.11.1}
\end{equation*}
$$

and

$$
\begin{equation*}
\frac{h(g \exp (-t) w)^{\mu}}{e^{t\left\langle-\mu, w_{0}\right\rangle}\left(g_{o} d(w)\right)^{\mu}}=e^{R_{-}\left(g_{o}, w_{0} ; t\right)} . \tag{5.11.2}
\end{equation*}
$$

However, the denominators in (5.11.1) and (5.11.2), respectively, are $\left(g_{o} \exp (-t) w_{o} d(w)^{-1}\right)^{\kappa \mu}$ and $\left(g_{o} \exp (-t) w_{o} d(w)\right)^{\mu}$. But then the left sides of (5.11.1) and (5.11.2), respectively, approach 1 as $t \rightarrow+\infty$ and $t \rightarrow-\infty$ by (5.5.23) and (5.5.22). This proves the theorem.
Q.E.D.

## 6. The Symplectic Structure of $\left(Z, \omega_{z}\right)$ and the Integration <br> $$
\text { of } \xi_{I}, I \in S(g)^{G}
$$

6.1. We recall some aspects of the theory of symplectic manifolds. See Chapter 4 in [15] for more details. Let ( $X, \omega_{X}$ ), or more simply $X$, if $\omega_{X}$ is understood, be a symplectic manifold of dimension $2 n$. That is, $X$ is a smooth (i.e., $C^{\infty}$ ) manifold of dimension $2 n$ and $\omega_{X}$ is a closed nonsingular smooth differential 2-form on $X$. Thus if $C^{\infty}(X)$ is the space of all smooth functions on $X$ and $\operatorname{Der} C^{\infty}(X)$ is the Lie algebra of all smooth vector fields on $X$ then $\omega_{X}(\xi, \eta) \in C^{\infty}(X)$ for any $\xi, \eta \in \operatorname{Der} C^{\infty}(X)$ and $\omega_{X}(\xi, \eta)$ is alternating in $\xi$ and $\eta$.

Now for any $\varphi \in C^{\infty}(X)$ one defines a (Hamiltonian) vector field $\xi_{\infty}$ on $X$ by the relation

$$
\begin{equation*}
\omega_{x}\left(\xi_{\varphi}, \eta\right)=\eta \varphi \tag{6.1.1}
\end{equation*}
$$

for all $\eta \in \operatorname{Der} C^{\infty}(X)$. Furthermore $C^{\infty}(X)$ inherits a Poisson structure (see Section 1.1) where for any $\varphi, \psi \in C^{\infty}(X)$ one defines $[\varphi, \psi]=\xi_{\varphi} \psi$. Moreover if $\operatorname{Ham} X=\left\{\xi_{\varphi} \mid \varphi \in C^{\infty}(X)\right\}$ then Ham $X$ is a Lie subalgebra of $\operatorname{Der} C^{\infty}(X)$ and (assuming $X$ is connected)

$$
\begin{equation*}
0 \rightarrow \mathbb{R} \xrightarrow{\text { inj }} C^{\infty}(X) \xrightarrow{\varphi \rightarrow \xi_{\varphi}} \text { Ham } X \rightarrow 0 \tag{6.1.2}
\end{equation*}
$$

is an exact sequence of Lie algebras defining $C^{\infty}(X)$ as a central extension of $\operatorname{Ham} X$. We have identified $\mathbb{R}$ here with the constant functions on $X$.

As one knows, one may express $\omega_{X}$ in terms of the Poisson structure.

Proposition 6.1. Assume that $\varphi_{1}, \ldots, \varphi_{2 n}$ is a global coordinate system on $X$. Let $\Psi_{i j} \in C^{\infty}(X)$ be defined by $\Psi_{i j}=\left[\varphi_{i}, \varphi_{j}\right]$. Then the $2 n \times 2 n$ matrix $\left\{\Psi_{i j}\right\}$ is invertible at all points and if $\left\{\Psi^{i j}\right\}$ is the inverse matrix of functions one has

$$
\begin{equation*}
\omega_{X}=\frac{1}{2} \sum_{i, j}^{2 n} \Psi^{i j} d \varphi_{i} \wedge d \varphi_{j} . \tag{6.1.3}
\end{equation*}
$$

Proof. That $\left\{\Psi_{i j}\right\}=\left\{\xi_{\Phi_{i}} \varphi_{j}\right\}$ is invertible is obvious from the tangent spacecotangent space isomorphism induced by $\omega_{X}$. Now for any $p, \psi \in C^{\infty}(X)$ one easily has

$$
\begin{equation*}
\omega_{X}\left(\xi_{\psi}, \xi_{\varphi}\right)=[\varphi, \psi] \tag{6.1.4}
\end{equation*}
$$

from (6.1.1). In particular $\omega_{X}\left(\xi_{\varphi_{D}}, \xi_{\varphi_{q}}\right)=\left[\varphi_{q}, \varphi_{p}\right]$. On the other hand it is immediate that if $\omega$ is the right side of (6.1.3) one also has $\omega\left(\xi_{\varphi_{p}}, \xi_{q_{q}}\right)=$ [ $\varphi_{q}, \varphi_{p}$ ]. This proves $\omega=\omega_{X}$.
Q.E.D.
6.2. Now symplectic manifolds arise in a number of ways. One of these is from the coadjoint orbits of Lie groups. We recall some of the details. See Chapter 5 in [15] for a more complete account. Let $A$ be a connected Lie group and let $a$ be its Lie algebra so that $a$ is an $A$-module with respect to the adjoint representation. Also $a^{\prime}$, the dual to $a$, is an $A$-module with respect to the coadjoint representation. The latter is defined by contragradience so that if $a \in A, x \in a$, and $g \in g^{\prime}$

$$
\begin{equation*}
\langle g, x\rangle=\langle a g, a x\rangle . \tag{6.2.1}
\end{equation*}
$$

Now let $O \subseteq a^{\prime}$ be an orbit of $A$ in $a^{\prime}$. That is, $O$ is a homogeneous space for $A$ of the form $O=A h$ for some $h \in a^{\prime}$. Now for any $y \in a$ let $\varphi^{y} \in C^{\infty}(O)$ be the function defined by $\varphi^{y}(g)=\langle g, y\rangle$ for any $g \in O$. Also for any $x \in a$ let $\xi_{x} \in \operatorname{Der} C^{\infty}(O)$ be the vector field on $O$ corresponding to the action of the
one-parameter group $\exp (-t) x$ on $O$. The vector field $\xi_{x}$ is characterized by the relation

$$
\begin{equation*}
\xi_{x} \varphi^{y}=\varphi^{[x, y]} \tag{6.2.2}
\end{equation*}
$$

for any $x, y \in a$.

Proposition 6.2.1. There exists a unique symplectic structure $\omega_{0}$ on $O$ (defining a symplectic manifold $\left(O, \omega_{O}\right)$ ) such that

$$
\begin{equation*}
\left[\varphi^{x}, \varphi^{y}\right]=\varphi^{[x, y]} \tag{6.2.3}
\end{equation*}
$$

for any $x, y \in a$. Furthermore if $\xi_{\Phi^{x}}$ is the Hamiltonian vector field corresponding to $\varphi^{x}$ then

$$
\begin{equation*}
\xi_{\Phi^{x}}=\xi_{x} \tag{6.2.4}
\end{equation*}
$$

Proof. The existence of a symplectic structure $\omega_{O}$ on $O$ where (6.2.3) is satisfied is established by Theorem 5.3.1 in [15]. The uniqueness follows from Proposition 6.1 since locally the set of functions $\left\{\varphi^{x}\right\}, x \in a$, contains a coordinate system. The relation (6.2.4) follows from the equality of the left sides of (6.2.2) and (6.2.3) for all $y \in a$.
Q.E.D.

Henceforth $\left(O, \omega_{o}\right)$, for any coadjoint orbit $O$, will be the symplectic manifold given by Proposition 6.2.1.

Now let the notation be as in Section 1 so the symmetric algebra $S(a)$ (here $F=\mathbb{R}$ ) has a Poisson structure. Also, $g_{i}, x_{j}$ are respectively a basis and a dual basis of $a^{\prime}$ and $a$. Furthermore $S(a)$ is regarded as the algebra of polynomial functions on $a^{\prime}$. Now for any $u \in S(a)$ let $\varphi^{u}$ be the restriction $u \mid O$.

Proposition 6.2.2. For any $u \in S(a)$, coadjoint orbit $O \subseteq a^{\prime}$ and $g \in O$, one has

$$
\begin{equation*}
\left(\xi_{w^{u}}\right)_{g}=\sum_{i}\left(\left(\partial\left(g_{i}\right) u\right)(g)\right)\left(\xi_{x_{i}}\right)_{g} \tag{6.2.5}
\end{equation*}
$$

Furthermore for $u, v \in S(a)$ one has

$$
\begin{equation*}
\left[\varphi^{u}, \varphi^{v}\right]=\varphi^{[u, v]} \tag{6.2.6}
\end{equation*}
$$

with respect to the Poisson structure on $S(a)$ defined by Proposition 6.2.1.
Proof. It is immediate from (1.1.2) that

$$
\begin{equation*}
\left(d \varphi^{u}\right)_{g}=\sum_{i}\left(\left(\partial\left(g_{i}\right) u\right)(g)\right)\left(d \varphi^{x_{i}}\right)_{g} \tag{6.2.7}
\end{equation*}
$$

But now (6.2.5) follows from (6.2.4) and (6.2.7). But (6.2.5) and (6.2.7) (for $v$ instead of $u$ ) implies

$$
\begin{equation*}
\left[\varphi^{u}, \varphi^{v}\right]=\sum_{i, j}\left(\partial\left(g_{i}\right) u\right)\left(\partial\left(g_{j}\right) v\right)\left[\varphi^{x_{i}}, \varphi^{x_{j}}\right] \tag{6.2.8}
\end{equation*}
$$

But $\left[\varphi^{x_{i}}, \varphi^{x_{j}}\right]=\varphi^{\left[x_{i}, x_{j}\right]}$ by (6.2.3) and $\varphi^{\left[x_{i}, x_{j}\right]}(g)=\left\langle g,\left[x_{i}, x_{i}\right]\right\rangle$. But now substituting (1.1.2) in (1.1.3) and then comparing (1.1.3) with (6.2.8) one has $\varphi^{[u, v]}(g)=[u, v](g)=\left[\varphi^{u}, \varphi^{v}\right](g)$. This proves (6.2.6).
Q.E.D.
6.3. Now let the notation be as in Section 1.2 where $F=\mathbb{R}$ so that $g$ is a real semi-simple Lie algebra with a fixed Cartan decomposition, and $Q$ is an invariant bilinear form which on any simple component is a positive multiple of the Killing form. Furthermore $a$ is an arbitrary subalgebra of $g$. Using $Q$, the inner product $Q_{*}$, and the Cartan decomposition we have associated to $a$ three subspaces $a^{0}, a^{\perp}$ and $a^{*}$ of $g$. We recall that $a^{o}$ is the $Q$-orthogonal subspace to $a, a^{\perp}$ is $Q_{*}$-orthocomplement to $a$ and $a^{*}$ is nonsingularly paired to $a$ by $Q$. Now let $f \in g$ be arbitrary and let $\left(a^{*}\right)_{f}=f+a^{*}$ using the notation of Section 1.6. Even though $\left(a^{*}\right)_{f}$ is only an affine subspace and not necessarily a linear subspace we note that

$$
\begin{equation*}
\mathfrak{g}=\left(a^{*}\right)_{f} \oplus a^{o} \tag{6.3.1}
\end{equation*}
$$

is still a direct sum. (That is, any element in $g$ has a unique sum decomposition relative to the summands in (6.3.1).) This is clear since $a^{0}=\left(a^{*}\right)^{\perp}$ by (1.2.5) and of course $f+g=g$. Let

$$
\begin{equation*}
P=g \rightarrow\left(a^{*}\right)_{f} \tag{6.3.2}
\end{equation*}
$$

be the projection on $\left(a^{*}\right)_{f}$ according to the decomposition (6.3.1). Now it is immediate from (6.3.1) that if $g \in a^{\prime}$ there exists a unique element $y_{g} \in\left(a^{*}\right)_{f}$ such that

$$
\begin{equation*}
Q\left(y_{g}, x\right)=\langle g, x\rangle \tag{6.3.3}
\end{equation*}
$$

for all $x \in a$ and the map

$$
\begin{equation*}
T: a^{\prime} \rightarrow\left(a^{*}\right)_{f}, \quad g \mapsto y_{g} \tag{6.3.4}
\end{equation*}
$$

is bijective and is in fact, clearly, a diffeomorphism.
Now let $A \subseteq G$ be the subgroup corresponding to $a$. For any $a \in A$ and $y \in\left(a^{*}\right)_{f}$ let $a \cdot y \in\left(a^{*}\right)_{f}$ be defined by

$$
\begin{equation*}
a \cdot y=P a y \tag{6.3.5}
\end{equation*}
$$

Proposition 6.3. Equation (6.3.5) defines an action of $A$ on $\left(a^{*}\right)_{f}$. Furthermore regarding $a^{\prime}$ as an $A$-manifold (i.e., a manifold on which $A$ operates) with *respect to the coadjoint representation the map (6.3.4) is an isomorphism of A-manifolds.

Proof. Let $a \in A, g \in a^{\prime}$, and $x \in a$. Then substituting $a^{-1} x$ for $x$ one has from the invariance of $Q$ that $Q\left(a y_{g}, x\right)=\langle a g, x\rangle=Q\left(y_{a g}, x\right)$. But now $a y_{g}-a \cdot y_{g} \in a^{o}$ so that $Q\left(a \cdot y_{g}, x\right)=Q\left(y_{a g}, x\right)$. But then $a \cdot y_{g}=y_{a g}$ from the injectivity of (6.3.4). This proves both statements of the proposition.
Q.E.D.

Proposition 6.3 enables one to carry over the coadjoint-symplectic theory of $A$ from $a^{\prime}$ to $\left(a^{*}\right)_{f}$. Thus any orbit $Y$ of $A$ in $\left(a^{*}\right)_{f}$ is a symplectic where the symplectic structure $\omega_{Y}$ is defined by

$$
\begin{equation*}
\omega_{Y}=\Gamma\left(\omega_{\Gamma^{-1}}\right), \tag{6.3.6}
\end{equation*}
$$

where of course the action of the diffeomorphism $\Gamma$ on differential forms is defined in the usual differential-geometric way.
6.4. We apply Proposition 6.3 to the case of the example of Sections 1.5 and $2-5$, where $g$ is a real split semi-simple Lie algebra. The notation is as in, say, Section 1.5, so that $a=\bar{b}, A=\bar{B}, a^{*}=b$ and $f$ is given by (1.5.4). Then $\left(a^{*}\right)_{f}=b_{f}=f+\mathfrak{b}$. However, we are interested in only one orbit of $\bar{B}$ in $b_{f}$.

Let $h_{i}, i=1, \ldots, l$, be the basis of $h$ defined by putting

$$
\begin{equation*}
h_{i}=\left[e_{\alpha_{i}}, e_{-\alpha_{i}}\right] . \tag{6.4.1}
\end{equation*}
$$

Now recalling the definition of the $2 l$-dimensional submanifold $Z \subseteq \mathscr{C}_{f}$ (see (2.2.3) and (2.2.4)) one has a global coordinate system $\rho_{i}, \gamma_{j} \in C^{\infty}(Z), i, j=$ $1, \ldots, l$, such that $y \in Z$ if and only if $y$ is of the form

$$
\begin{equation*}
y=f+\sum_{i=1}^{l} \rho_{i}(y) h_{i}+\sum_{j=1}^{l} \gamma_{j}(y) e_{\alpha_{j}} \tag{6.4.2}
\end{equation*}
$$

where $\rho_{i}(y) \in \mathbb{R}$ and $\gamma_{j}(y) \in \mathbb{R}^{*}$ are arbitrary. Of course since the $\gamma_{j}$ are positive valued they are of course invertible. We recall also that for any invariant $I \in S(g)^{G}$ we defined a vector field $\xi_{I}$ on $Z$ by (2.2.18) and (2.2.19) so that for any $y \in Z$ one has $\left(\xi_{I}\right)_{y}=\left[y,\left(\delta_{\bar{B}} I\right)(y)\right]$ where $\delta_{\beta} I$ is defined by (1.2.7). See (1.2.10).

Proposition 6.4. $Z$ is an orbit of $\bar{B}$ in $b_{f}$. Furthermore the corresponding symplectic structure on $Z$ is given by

$$
\begin{equation*}
\omega_{Z}=\sum_{i=1}^{l} d \rho_{i} \wedge \gamma_{i}^{-1} d \gamma_{i} \tag{6.4.3}
\end{equation*}
$$

Moreover if $\xi_{I \mid Z}$ is the Hamiltonian vector field corresponding to $I \mid Z \subseteq C^{\infty}(Z)$ for any $I \in S(g)^{G}$ then one has

$$
\begin{equation*}
\xi_{I \mid Z}=\xi_{I} \tag{6.4.4}
\end{equation*}
$$

Proof. Let $\bar{b} \in \mathscr{b}$. We may write $\bar{b}=h \exp x$, where $x \in \bar{n}$ and $h \in H$. Let $c_{i} \in \mathbb{R}$ be defined by $x+\sum c_{i} e_{-\alpha_{i}} \in[\bar{n}, \bar{n}]$. Now let $y=f+e$, where $e=$ $\sum_{i=1}^{l} e_{\alpha_{i}}$. Then $y \in Z \subseteq b_{f}$. Furthermore if $z=f+\sum_{i=1}^{l} c_{i} h_{i}+\sum a_{i} e_{\alpha_{i}}$, where $h_{i}$ is given by (6.4.1) and $a_{i}=h^{\alpha_{i}}$, then $z \in Z$. But one notes that

$$
\begin{equation*}
\bar{b} y-z \in \bar{n} . \tag{6.4.5}
\end{equation*}
$$

On the other hand $\bar{n}=(\bar{k})^{o}$ so that $b \cdot y=z$ using the notation of (6.3.5). However, it is clear that $z$ is an arbitrary element of $Z$ and hence $Z=\bar{B} \cdot y$. This proves that $Z$ is a $\bar{B}$ orbit in $b_{f}$.

Now let $\Gamma: \mathscr{C}^{\prime} \rightarrow \mathscr{b}_{f}$ be the map (6.3.4) where of course $\bar{C}=a$. Thus if $O=\Gamma^{-1}(Z)$ then $O$ is a coadjoint orbit of $\bar{B}$ in $\vec{b}^{\prime}$. For any $u \in S\left(Z^{\prime}\right)$ let $\varphi^{u} \in C^{\infty}(O)$ be defined as in Section 6.1 and let $\psi^{u} \in C^{\infty}(Z)$ be defined by putting $\psi^{u}=u \mid Z$. It is clear from (6.3.3) that $\psi^{u} \circ \Gamma=\varphi^{u}$ and hence, using (6.2.3), with regard to the Poisson structure in $S\left({ }^{6}\right)$ and $C^{\infty}(Z)$ one has

$$
\begin{equation*}
\left[\psi^{u}, \psi^{v}\right]=\psi^{[u, v]} \tag{6.4.6}
\end{equation*}
$$

for $u, v \in S(\bar{b})$.
Now let $z_{j} \in \mathscr{\ell}, j=1, \ldots, l$, be the basis of $\ell$ such that

$$
\begin{equation*}
Q\left(z_{j}, h_{i}\right)=\delta_{i j} . \tag{6.4.7}
\end{equation*}
$$

But from (1.5.2) it follows that

$$
\begin{equation*}
Q\left(x, h_{i}\right)=\left\langle\alpha_{i}, x\right\rangle \tag{6.4.8}
\end{equation*}
$$

for any $x \in h$. Thus one must have

$$
\begin{equation*}
\left[e_{-\alpha_{i}}, z_{j}\right]=\delta_{i j} e_{-\alpha_{i}} . \tag{6.4.9}
\end{equation*}
$$

On the other hand it follows from (6.4.7) and (6.4.2) that

$$
\begin{equation*}
\psi^{z_{j}}=\rho_{j} . \tag{6.4.10}
\end{equation*}
$$

Furthermore using (1.5.2) one notes that

$$
\begin{equation*}
\psi^{e-\alpha_{i}}=\gamma_{i} \tag{6.4.11}
\end{equation*}
$$

Thus by (6.4.6) and (6.4.9) one has, in $C^{\infty}(Z)$,

$$
\begin{equation*}
\left[\gamma_{i}, \rho_{j}\right]=\delta_{i j} \gamma_{i} \tag{6.4.12}
\end{equation*}
$$

On the other hand since clearly $\psi^{z}=0$ for any $z \in[\bar{n}, \bar{n}]$ and since $\ell$ is commutative it follows from (6.4.10) and (6.4.11) that for all $i, j$

$$
\begin{equation*}
\left[\gamma_{i}, \gamma_{j}\right]=\left[\rho_{i}, \rho_{j}\right]=0 \tag{6.4.13}
\end{equation*}
$$

On the other hand the $\gamma_{i}$ together with the $\rho_{j}$ define a coordinate system in $Z$. We may then apply Proposition 6.1 to compute $\omega_{Z}$ in terms of $d \gamma_{i}$ and $d p_{j}$, using of course (6.4.12) and (6.4.13). In fact the computation is immediate and yields (6.4.3).

Now for any $x \in \mathscr{C}$ let $\zeta_{x}$ be the vector field in $Z$ which corresponds to $\xi_{x}$ by the isomorphism $\Gamma$. See (6.3.4). Thus by (6.2.2) one has

$$
\begin{equation*}
\zeta_{x} \psi^{z}=\psi^{[x, z]} \tag{6.4.14}
\end{equation*}
$$

for $x, z \in \mathscr{\ell}$. Now let $x_{i}$ be a basis of $\mathscr{\ell}$ and $y_{j}$ be a basis of $\mathscr{\ell}$ such that $Q\left(x_{i}, y_{j}\right)=$ $\delta_{i j}$. Now for any $\varphi \in C^{\infty}(Z)$ let $\xi_{\varphi}$ here be the Hamiltonian vector field on $Z$ corresponding to $\varphi$. Now for any $u \in S(g)$ and $z \in g$ let $i(z) u \in S(g)$ be defined as in Section 1.2. For any $y \in Z$ we now assert that

$$
\begin{equation*}
\left(\xi_{u \mid z}\right)_{y}=\sum_{j}\left(i\left(y_{j}\right) u\right)(y)\left(\zeta_{x_{j}}\right)_{y} . \tag{6.4.15}
\end{equation*}
$$

Indeed for any $u \in S(g)$ it is clear from (6.3.1) that there exists a unique element $\bar{u} \in S(\bar{b})$ such that $\bar{u}\left|b_{f}=u\right| \ell_{f}$. Furthermore since $b_{f}$ is stable under translation by elements in $\ell$ one has upon differentiation, $\overline{i(z) u}=i(z) \bar{u}$ for $z \in \ell$. But then both sides of (6.4.15) do not change if $\bar{u}$ is substituted for $u$. However, for $\bar{u}$ one has (6.4.15) by (6.2.5). One recalls (6.3.3) and (6.3.6). This proves (6.4.15).

But now we may identify the tangent space to $Z$ at $y$ with $\ell_{i}+d_{1} \subseteq b$ as in (2.2.10). But by (6.2.2) $\left(\zeta_{x} \psi^{z}\right)(y)=Q(y,[x, z])=Q([y, x], z)$. Thus if $P_{\ell}: g \rightarrow b$ is the $Q_{*}$-orthogonal projection (i.e., Ker $P_{\ell}=\bar{n}=b^{\circ}$ ) then clearly

$$
\begin{equation*}
\left(\zeta_{x}\right)_{y}=P_{\epsilon}[y, x] \in \hbar+d_{1} \tag{6.4.16}
\end{equation*}
$$

Now recalling the definition of $\delta_{\beta}^{-} u$ in Section 1.2 one then has

$$
\begin{equation*}
\left(\xi_{u \mid z}\right)_{y}=P_{\ell}\left[y,\left(\delta \delta_{\beta}^{-} u\right)(y)\right] \tag{6.4.17}
\end{equation*}
$$

by (1.2.10). However, if $u=I \in S(g)^{G}$ then $\left(\xi_{1}\right)_{y}=\left[y,\left(\delta_{\bar{d}} I\right)(y)\right] \in \ell$ by (2.2.22). Thus

$$
\xi_{I}=\xi_{I \mid Z}
$$

6.5. Now using the isomorphism (3.5.27) which maps $\ell$ to $\ell^{\prime}$ the bilinear form $Q \mid h$, as one knows, induces a positive definite bilinear form $Q \mid h^{\prime}$ on $h^{\prime}$. For $i, j=1, \ldots, l$ let

$$
\begin{equation*}
b_{i j}=Q\left(\alpha_{i}, \alpha_{j}\right) . \tag{6.5.1}
\end{equation*}
$$

Remark 6.5. Clearly $b_{i j}$ depends upon the choice of $Q$. However, if $C$ is the $l \times l$ matrix defined by $C_{i j}=b_{i j} / b_{i i}$ then $C$ is of course the Cartan matrix and is independent of $Q$.

Now among the fundamental invariants $I_{j} \in S(g)^{G}, j=1, \ldots, l$, we may clearly fix $I_{1}$ so that

$$
\begin{equation*}
I_{1}(x)=\frac{1}{2} Q(x, x) \tag{6.5.2}
\end{equation*}
$$

for any $\boldsymbol{x} \in \mathcal{g}$.
Proposition 6.5. Let $\rho_{i}, \gamma_{j} \in C^{\infty}(Z)$ be defined by (6.4.2). Since $\gamma_{j}(Z) \subseteq \mathbb{R}^{*}$ we may define $\varphi_{j} \in C^{\infty}(Z)$ by putting $\varphi_{j}=\log \gamma_{j}$. Then $\rho_{i}, \varphi_{j}, i, j=1, \ldots, l$, define a coordinate system on $Z$. In fact the map

$$
\begin{equation*}
Z \rightarrow \mathbb{R}^{2 l} \tag{6.5.3}
\end{equation*}
$$

given by $y \mapsto\left(\rho_{1}(y), \ldots, \rho_{l}(y), \varphi_{1}(y), \ldots, \varphi_{l}(y)\right)$ is a diffeomorphism. Moreover

$$
\begin{equation*}
\omega_{Z}=\sum_{i=1}^{l} d \rho_{i} \wedge d \varphi_{i} \tag{6.5.4}
\end{equation*}
$$

Furthermore if $b_{i j}$ is defined by (6.5.1) then

$$
\begin{equation*}
I_{1} \left\lvert\, Z=\frac{1}{2} \sum_{i, j=1}^{l} b_{i j} \rho_{i} \rho_{j}+\sum_{i=1}^{l} e^{\Phi_{i}}\right. \tag{6.5.5}
\end{equation*}
$$

Proof. The first statement is immediate from the definition of $Z$. See (2.2.3) and (2.2.4). But now $d \varphi_{i}=\gamma_{j}^{-1} d \gamma_{i}$ and hence (6.5.4) follows from (6.4.3). On the other hand if $y \in Z$ is given by (6.4.2) then $I_{1}(y)=\frac{1}{2} Q(y, y)=$ $Q\left(f, \sum_{i} \gamma_{i}(y) e_{\alpha_{i}}\right)+\frac{1}{2} \sum_{i, j} \rho_{i}(y) \rho_{j}(y) Q\left(h_{i}, h_{j}\right)$. But $Q\left(f, e_{\alpha_{i}}\right)=1$ and hence $Q\left(f, \sum_{i} \gamma_{i}(y) e_{\alpha_{i}}\right)=\sum_{i} \gamma_{i}(y)=\sum_{i} e^{\phi_{i}(y)}$. On the other hand $h_{i} \rightarrow \alpha_{i}$ by (3.5.27), recalling (3.5.28), (3.1.14), and (6.4.1) so that $b_{i j}=Q\left(h_{i}, h_{j}\right)$. One thus obtains (6.5.5).
Q.E.D.
6.6. Now let $\left(X, \omega_{X}\right)$ be a symplectic manifold. Let $2 n=\operatorname{dim} X$. Let $T_{p}(X)$, for any $p \in X$, be the tangent space to $X$ at $p$. Assume that for any $p \in X$ one has an $n$-dimensional subspace $F_{p} \subseteq T_{p}(Z)$ and that the map given by $p \mapsto F_{p}$ is a smooth involutory distribution $F$ (in the sense of Cartan) on $X$.

By a leaf of $F$ we mean a maximal connected integral submanifold $M \subseteq X$. Thus $F$ is called a real polarization of $\left(X, \omega_{X}\right)$ in case $\omega_{X} \mid M=0$ for any leaf $M$ of $F$. That is, the leaves of $F$ are Lagrangian submanifolds of $X$.

Now assume that $F$ is a real polarization of $\left(X, \omega_{X}\right)$. For any open subset $U \subseteq X$ let $C_{F}^{\infty}(U)=\left\{\varphi \in C^{\infty}(U) \mid v \varphi=0\right.$ for all $p \in U$ and $\left.v \in F_{p}\right\}$. It follows easily from (6.1.1) and the maximal isotropic property of $F_{p}$ that, using the notation of Section 6.1, one has

$$
\begin{equation*}
\left(\xi_{\varphi}\right)_{p} \in F_{p} \text { for any } p \in U \text { and } \varphi \in C_{F}^{\infty}(U) \tag{6.6.1}
\end{equation*}
$$

Thus if $M$ is any leaf of $F$ and $U \subseteq X$ is any open subset then the restriction

$$
\begin{equation*}
\xi_{\varnothing} \mid M \cap U \text { is tangent to } M \cap U \text { for any } \varphi \in C_{F}^{\infty}(U) \tag{6.6.2}
\end{equation*}
$$

But then (6.6.1) implies $[\varphi, \psi]=0$, and hence

$$
\begin{equation*}
\left[\xi_{\varphi}, \xi_{\psi}\right]=0 \text { for any } \varphi, \psi \in C_{F}^{\infty}(U) \tag{6.6.3}
\end{equation*}
$$

But now if $M$ is a leaf of $F$ then it is an easy consequence of (6.6.2) and (6.6.3) that there exists a unique flat affine connection on $M$ such that $\xi_{\infty} \mid M \cap U$ is covariant constant for any open set $U \subseteq X$ and $\varphi \in C_{F}{ }^{\infty}(U)$. See Section 4.2. We will refer to this as the affine connection on $M$ induced by $F$. The real polarization $F$ will be said to be complete in case the affine connection induced by $F$ on $M$ is complete for every leaf of $F$. That is, all parametrized geodesics on $M$ are defined for all values of the parameter.

Remark 6.6.1. In case $\left(X, \omega_{X}\right)$ is a coadjoint orbit of a simply connected exponential solvable Lie group $A$ and $F$ is an $A$-invariant real polarization the condition of completeness has been called the Pukansky condition. A theorem of Pukansky (see [21]) asserts that the unitary representation of $A$ associated to $F$ is irreducible if and only if $F$ is complete.

An important property of completeness is that the vector field $\xi_{\varphi}$ for any $\varphi \in C_{F}{ }^{\infty}(X)$, can be globally integrated.

Proposition 6.6. Assume that $F$ is a complete, real polarization of a symplectic manifold $\left(X, \omega_{X}\right)$. Then for any $\varphi \in C_{F}{ }^{\infty}(X)$ there exists a one-parameter group, $\exp t \xi_{\infty}, t \in \mathbb{R}$, of symplectic diffeomorphisms of $X$ such that for any $\psi \in C^{\infty}(X)$ and $p \in X$

$$
\begin{equation*}
\left.\frac{d}{d t} \psi\left(\exp (-t) \xi_{\infty} \cdot p\right)\right|_{t=0}=\left(\xi_{\Phi} \psi\right)(p) \tag{6.6.4}
\end{equation*}
$$

where the dot denotes the diffeomorphism action on $X$.

Proof. Let $\varphi \in C_{F}^{\infty}(X)$ and let $p \in X$. Let $M$ be the leaf of $F$ containing $p$. Let $c(t) \in M$, for all $t \in \mathbb{R}$, be the geodesic on $M$ whose tangent vector at $t=0$ is $\left(-\xi_{\varphi}\right)_{\mathcal{D}}$. Then $\left(\exp t \xi_{\varphi}\right) \cdot p$ is defined so that $c(t)=\left(\exp t \xi_{\varphi}\right) \cdot p$. It is immediate then that $\exp t \xi_{\Phi}$ is a one-parameter group of symplectic diffeomorphisms satisfying (6.6.4).
Q.E.D.

Remark 6.6.2. In general if $\xi$ is a globally integrable vector field on a manifold $M$ then the corresponding one parameter group of diffeomorphisms will be written as $\exp t \xi$. The direction of the trajectories $\exp t \xi \cdot p$ of $\xi$ will also be taken in the reverse direction of $\xi$. That is, if $\psi \subset C^{\infty}(M)$ one has $\left.(d / d t) \psi(\exp (-t) \xi \cdot p)\right|_{t=0}=(\xi \psi)(p)$ for any $p \in M$. This of course guarantees the correct functorial properties for a Lie group of diffeomorphisms. Furthermore it produces the correct signs for Hamilton's equations. See (7.1.5).
6.7. Now recalling Section 4.1, $\mathscr{Z}$ is an involutory distribution of $\operatorname{dim} l$ on $Z$. To conform to the notation of Section 6.6 we write $F=\mathscr{Z}$ and recall then that for any $y \in Z$ (see (2.2.21))

$$
\begin{equation*}
F_{y}=\left\{\left(\xi_{I}\right)_{y} \mid I \in S(g)^{G}\right\} \tag{6.7.1}
\end{equation*}
$$

where we recall that (see (2.2.22))

$$
\begin{equation*}
\left(\xi_{I}\right)_{y}=\left[y,\left(\delta_{\bar{\delta}} I\right)(y)\right], \tag{6.7.2}
\end{equation*}
$$

where $\delta_{G} I$ is defined by (1.2.7).
Furthermore by Proposition 4.1 and (3.5.26) the leaves of $F$ are all the submanifolds of the form $Z(\gamma), \gamma \in \mathscr{I}\left(h_{+}\right)$, and any such leaf has been given (see end of Section 4.2) the structure of a flat affinely connected manifold.

Theorem 6.7.1. Let $g$ be a split semi-simple Lie algebra. Let $l=\operatorname{rank}_{g}$ and let $Z \subseteq g$, defined as in Section 2.2, be the 2l-dimensional submanifold of normalized Jacobi elements. Let $\omega_{Z}$ be the symplectic structure on $Z$ defined as in Section 6.3 and given explicitly by (6.4.3) and let $F$ be the l-dimensional involutory distribution on $Z$ defined by (6.7.1). Let $\mathscr{I}: g \rightarrow \mathbb{R}^{l}$ be the map defined as in (2.3.1) and let $h_{+}$be the open Weyl chamber defined by (3.3.1). For each $\gamma \in \mathscr{I}\left(h_{+}\right)$let $Z(\gamma)$ be defined by (2.3.2) so that by Proposition 4.2 and (3.5.26) the $Z(\gamma)$ are flat affinely connected manifolds and are the leaves of $F$.

Then $F$ is a complete real polarization of $\left(Z, \omega_{Z}\right)$. Furthermore the affine connecon $Z(\gamma)$ is the same as the one induced by $F$. Moreover for any invariant $I \in S(g)^{G}$ one has $I \mid Z \in C_{F}{ }^{\infty}(Z)$ and $\xi_{I}$ is the corresponding Hamiltonian vector field. In particular $\xi_{1}$ is globally integrable. In fact let $y \in Z$ and let $I \in S(g)^{G}$. Let $\gamma=\mathscr{F}(y)$ so that by (3.5.26) and Proposition 3.3.1 there exists a unique $w_{o} \in \hbar_{+}$such that $\mathscr{I}\left(w_{o}\right)=\gamma$. Let $w=f+w_{0}$, where $f$ is given by (1.5.4), and let $z=(\delta I)(w)$
so that $\boldsymbol{z} \in \mathfrak{g}^{w \prime}$ using the notation of Section 1.2 and recalling (1.3.3). Let $G_{0}{ }^{w}$ be as (3.5.5)-(3.5.7) and let $\beta_{(w)}: G_{0}{ }^{w} \rightarrow Z(\gamma)$ be the isomorphism (3.6.2). Let $g=\beta_{(w)}^{-1}(y) \in G_{0}{ }^{w}$. Then for all $t \in \mathbb{R}$ one has $g \exp (-t) z \in G_{0}{ }^{w}$ and

$$
\begin{equation*}
\beta_{(w)}(g \exp (-t) z)=\exp t \xi_{I} \cdot y, \tag{6.7.3}
\end{equation*}
$$

where the dot is the same as in Proposition 6.6.
Proof. For any $I \in S(g)^{g}$ we recall that $I^{f} \in S(f)$ has been defined by the relation $I^{f}(x)=I(f+x)$ for $x \in 6$. See Section 1.4. For any $u \in S(\bar{b})$ let, as in the proof of Proposition 6.4, $\psi^{u} \in C^{\infty}(Z)$ be defined by $\psi^{u}=u \mid Z$. We assert that

$$
\begin{equation*}
\psi^{y^{\prime}}=I \mid Z . \tag{6.7.4}
\end{equation*}
$$

Indeed if $y \in Z$ we may write $y=f+x$, where $x \in \mathscr{\ell}$. Then $I(y)=I(f+x)=$ $I^{f}(x)=I^{f}(f+x)=I^{f}(y)$ using (1.2.6) since $f \in(\bar{E})^{0}=\bar{n}$. This proves (6.7.4). But now if $I, J \in S(g)^{G}$ and we put $u=I^{f}, v=J^{f}$ then $\xi_{d^{u}}=\xi_{I}$ and $\xi_{u^{v}}=\xi_{J}$ by (6.4.4). But then by (6.1.4) and (6.4.6)

$$
\begin{equation*}
\omega_{Z}\left(\xi_{I}, \xi_{J}\right)-\left[\psi^{n}, \psi^{n}\right]=\psi^{[p, u]} . \tag{6.7.5}
\end{equation*}
$$

But $[v, u]=\left[J^{f}, I^{f}\right]=0$ by Theorem 1.4. Thus (6.7.5) vanishes and hence $Z(\gamma)$, recalling Proposition 4.1, is Lagrangian for all $\gamma \in \mathscr{I}\left(f_{+}\right)$. Thus $F$ is a real polarization. Furthermore by definition $I \mid Z$ is constant on $Z(\gamma)$ so that $I \mid Z \in C_{F}{ }^{\infty}(Z)$ and hence $\xi_{I}$ is covariant constant with respect to the affine connection induced by $F$. Recalling the definition (see end of Section 4.2) of the given affine connection on $Z(\gamma)$ this proves that both affine connections are the same.

Now let $y \in Z$. Put $\gamma=\mathscr{I}(y)$. By (3.5.26) and Proposition 3.3.1 there exists a unique $w_{o} \in h_{+}$such that $\mathscr{I}\left(w_{o}\right)=\gamma$. Then $Z(\gamma)$ is the leaf of $F$ containing $y$ and by Theorem 4.3

$$
\begin{equation*}
\beta_{(w)}: G_{0}{ }^{w} \rightarrow Z(\gamma) \tag{6.7.6}
\end{equation*}
$$

is an isomorphism of flat affinely connected manifolds where the affine connection on $G_{0}{ }^{w} \subseteq \tilde{G}^{w}$ is defined by the Abelian group structure on $\tilde{G}^{w}$. In particular the geodesics in $G_{0}{ }^{w}$ are translates of one-parameter groups $\exp t x$, where $x \in \mathcal{g}^{w}$. Since $G_{0}{ }^{w}$ is a connected component of $\vec{G}^{w}$ it is obviously complete. Thus $Z(\gamma)$ is complete and hence $F$ is complete.

Now let $I \in S(g)^{G}$ and let $g \in G_{0}{ }^{w}$ be such that $\beta_{(w)}(g)=y$. Let $L_{I}{ }^{w}$ be the vector field on $\mathcal{G}^{w}$, defined by (4.3.1), so that if $\exp t L_{I^{w}}$ is the one-parameter group of diffeomorphisms of $\boldsymbol{G}^{w}$ defined by $L_{I}^{w}$, one has $\exp t L_{I}^{w} \cdot \boldsymbol{a}=$
$a \exp (-t)(\delta I)(w)$ for any $a \in G^{w}$. But $\beta_{(w)}\left(L_{I}^{w} \mid G_{0}{ }^{w}\right)=\xi_{I} \mid Z(\gamma)$ by (4.3.3). This proves (6.7.3) since $z=(\delta I)(w)$.
Q.E.D.

One also would know the flow defined by $\xi_{I}$ if one knew the coordinate values of $\exp t \xi_{I} \cdot y$. For "half" the coordinates this is given explicitly by

Theorem 6.7.2. Let $y \in Z$ and let $w_{0} \in h_{+}$, as in Theorem 6.7.1, be such that $\mathscr{I}\left(w_{o}\right)=\mathscr{I}(y)$. Also as in Theorem 6.7.1 let $w=f+w_{o}$ and let $g \in G_{0}{ }^{w}$ be such that $\beta_{(w)}(g)=y$. Now let $H$ be the split Cartan subgroup defined as in Section 1.5 and also as in Section 1.5 let $\alpha_{1}, \ldots, \alpha_{l}$ be a corresponding set of simple positive roots. For any $a \in G_{0}{ }^{w}$ let $h(a) \in H$ be defined by (3.4.14). Let $I \in S(g)^{G}$ and let $(\delta I)(w)=z \in g^{w}$. Finally let $\varphi_{1}, \ldots, \varphi_{l} \in C^{\infty}(Z)$ be "half" the coordinates on $Z$ defined as in Proposition 6.5. Then (see Theorem 6.7.1)

$$
\begin{equation*}
\varphi_{i}\left(\exp t \xi_{I} \cdot y\right)=\log h(g \exp (-t) z)^{-\alpha_{i}} \tag{6.7.7}
\end{equation*}
$$

Proof. Now, recalling Theorem 3.6, if $a \in G_{0}{ }^{w}$ then we may write $a=$ $s(\kappa) \bar{n}(a) h(a) n(a)$, where the four factors are defined as in the statement of Theorem 3.6. Furthermore $\beta_{(w)}(a)=n(a) w \in Z$ by Theorem 3.6 and hence we can write $n(a) w=f+x+\sum_{i=1}^{l} c_{i} e_{\alpha_{i}}$, where $x \in h$. But by (3.7.4) one has in fact

$$
\begin{equation*}
c_{i}=h(a)^{-\alpha_{i}} . \tag{6.7.8}
\end{equation*}
$$

But now if $\gamma_{i} \in C^{\infty}(Z)$ is defined by (6.4.2) then (6.7.8) implies $\gamma_{i}\left(\beta_{(w)}(a)\right)=$ $h(a)^{-\alpha_{i}}$. Thus one has

$$
\begin{equation*}
\varphi_{i}\left(\beta_{(w)}(a)\right)=\log h(a)^{-\alpha_{i}} \tag{6.7.9}
\end{equation*}
$$

But then (6.7.7) follows from (6.7.3) by putting $a=g \exp (-t) z . \quad$ Q.E.D.
6.8. We are mainly interested in applying Theorems 6.7 .1 and 6.7 .2 for the case where $I$ is the quadratic invariant $I_{1}$. See (6.5.2). In that case the element $z$ becomes $w$ itself.

Our solution to the generalized Toda lattice will depend on the following theorem together with the explicit formula for $h(g \exp (-t) w)^{\lambda}$ given in Theorem 5.10. See (5.10.6) and (5.10.7).

Theorem 6.8.1. Let $\left\{\rho_{i}, \varphi_{j}\right\}, i, j=1, \ldots, l$, be the coordinate system on $Z$ given in Proposition 6.5. Let $\omega_{Z}$ be the symplectic structure on $Z$ defined in Section 6.3 and given explicitly (see (6.5.4)) by

$$
\begin{equation*}
\omega_{Z}=\sum_{i=1}^{l} d \rho_{i} \wedge d \varphi_{i} \tag{6.8.1}
\end{equation*}
$$

Let $I_{1} \in S(g)^{G}$ be defined by (6.5.2) so that $I_{1} \mid Z$ has been given explicitly (see (6.5.5)) by

$$
\begin{equation*}
I_{1} \left\lvert\, Z=\frac{1}{2} \sum_{i, j=1}^{l} b_{i j} \rho_{i} \rho_{j}+\sum_{i=1}^{l} e^{\varphi_{i}}\right. \tag{6.8.2}
\end{equation*}
$$

where $b_{i j}=Q\left(\alpha_{i}, \alpha_{j}\right)$ Let $\xi$ be the Hamiltonian vector field on $Z$ corresponding to $I_{1} \mid Z$. Then $\xi$ is globally integrable. Let $\exp t \xi$ be the corresponding flow. (See Theorem 6.7.1.) Let $y \in Z$ and let $w_{o} \in h_{+}$be such that $\mathscr{I}(y)=\mathscr{I}\left(w_{o}\right)$. Let $w=f+w_{0}$ and let $g \in G_{0}{ }^{w}$ be such that $\beta_{(w)}(g)=y$. (We are using the notation of Theorem 6.7.1.) Then

$$
\begin{equation*}
\exp t \xi \cdot y=\beta_{(w)}(g \exp (-t) w) \tag{6.8.3}
\end{equation*}
$$

Furthermore

$$
\begin{equation*}
\varphi_{i}(\exp t \xi \cdot y)=\log h(g \exp (-t) w)^{-\alpha_{i}} \tag{6.8.4}
\end{equation*}
$$

$i=1, \ldots, l$.
Proof. Theorem 6.8.1 follows immediately from Theorems 6.7.1 and 6.7.2 as soon as one shows that $\left(\delta I_{i}\right)(w)=w$. Let $\left\{x_{i}\right\}$ and $\left\{y_{j}\right\}$ both be bases of $g$ such that $Q\left(x_{i}, y_{j}\right)=\delta_{i j}$. Then clearly

$$
\begin{equation*}
I_{1}=\frac{1}{2} \sum_{i=1}^{\mathrm{dim} \mathcal{g}} x_{i} y_{i} \tag{6.8.5}
\end{equation*}
$$

But then for any $z \in \mathscr{g}$ one has

$$
\begin{equation*}
\left(\delta I_{1}\right)(z)=\frac{1}{2}\left(\sum_{i} Q\left(x_{i}, z\right) y_{i}+\sum_{i} Q\left(y_{i}, z\right) x_{i}\right) \tag{6.8.6}
\end{equation*}
$$

by (1.2.1) and (1.2.2). But clearly both sums in (6.8.6) are equal to $z$. Thus $\left(\delta I_{1}\right)(z)=z$.
Q.E.D.

Using the diagram (3.6.7) one may give a conceptually simpler expression (in that it only involves $y$ itself) for the trajectory $\exp t \xi \cdot y$. However, it does not seem to lead to an explicit formula.

Theorem 6.8.2. Let the notation be as in Theorem 6.8.1. Then there exist for all $t \in \mathbb{R}$ unique elements $\bar{n}_{t} \in \bar{N}, h_{t} \in H$, and $n_{t} \in N$, where $\bar{N}, H$, and $N$ are as in Section 3.2, such that

$$
\begin{equation*}
\exp t y=\bar{n}_{t} h_{t} n_{t} \tag{6.8.7}
\end{equation*}
$$

Moreover with respect to the adjoint action one has $n_{-t} y \in Z$ and in fact

$$
\begin{equation*}
\exp t \xi \cdot y=n_{-t} y \tag{6.8.8}
\end{equation*}
$$

Proof. Recalling the notation and statement of Theorem 3.6 one has $\beta_{(w)}(g \exp t w)=\beta_{y}\left(\psi_{g}(g \exp t w)\right)$. See (3.6.7). But $\psi_{g}(g \exp t w)=$ $n(g) \exp \operatorname{twn}(g)^{-1}$. However, $n(g) w=\beta_{(w)}(g)=y$. Thus $\psi_{g}(g \exp t w)=\exp t y$ and hence

$$
\begin{equation*}
\beta_{(w)}(g \exp t w)=\beta_{y}(\exp t y) . \tag{6.8.9}
\end{equation*}
$$

But now if we write $y(t)=\exp t y$ then $y(t) \in G_{e}{ }^{y}$ and hence by Theorem 3.2 one has (6.8.7), where $\bar{n}_{t}=\bar{n}_{y(t)}, h_{t}=h_{y(t)}$, and $n_{t}=n_{y(t)}$. But for any $d \in G_{e}{ }^{y}$ one has $\beta_{y}(d)=n_{d} y$. Thus $\beta_{y}(\exp t y)=n_{t} y$. But then (6.8.8) follows from (6.8.9) and (6.8.3).
Q.E.D.

Remark 6.8. The solution to the generalized Toda lattice will be based on (6.8.4) and formula (5.10.7) for $h(g \exp (-t) w)^{\lambda}$. Given the "initial condition" $y$ what is needed is only the "spectrum" of $y$, namely, $w_{o}$. The "input" is $w_{o}$ together with the constants (the exponential action angle coordinates of the initial condition $y$ ) $g_{o}{ }^{\nu}$ for $v \in \Lambda$. This, however, is determined from (5.5.25). That is, one easily determines $n(g)$ inductively by the relation $n(g) w=y$. See, e.g., (7.8.25). But then (5.5.25) yields $g_{o}^{\nu_{i}}$ and hence $g_{o}^{\nu}$ for any $\nu \in \Lambda$. One notes (see (3.7.4)) that $y$ itself has $h(g)^{-\alpha_{i}}$ are coordinates so that $h(g)^{-\lambda}$ is known in (5.5.25). Also we remark that $n(g)$ is naturally determined by $y$. That is, one has

$$
\begin{equation*}
n(g) \bar{n}_{f}(w) w_{o}=y \tag{6.8.10}
\end{equation*}
$$

and one sees easily from the injectivity of (2.4.5) that $n(g)$ and $\bar{n}_{f}(w)$ are the unique elements $n \in N$ and $\bar{n} \in \bar{N}$ such that $n \bar{n} w_{o}=y$. Finally it is then to be noted that $\bar{n}$ depends only on the isospectral leaf containing $y$ whereas $n$ depends on the "action angle" coordinates of $y$.

## 7. Dénouement; the Formula for $q_{i}(x(t))$

7.1. Let $\left(\mathbb{R}^{2 n}, \omega\right)$ be the classical $2 n$-dimensional phase space. That is, one has linear (canonical) coordinates $p_{i}, q_{j} \in C^{\infty}\left(\mathbb{R}^{2 n}\right), i, j=1, \ldots, n$, and the symplectic structure is given by

$$
\begin{equation*}
\omega=\sum_{i=1}^{n} d p_{i} \wedge d q_{i} \tag{7.1.1}
\end{equation*}
$$

Onc easily sees that if $H \in C^{\infty}\left(\mathbb{R}^{2 n}\right)$ then the corresponding Hamiltonian vector field $\xi_{H}$ (see (6.1.1)) is given explicitly by

$$
\begin{equation*}
\xi_{H}=\sum_{i=1}^{l} \frac{\partial H}{\partial q_{i}} \frac{\partial}{\partial p_{i}}-\frac{\partial H}{\partial p_{i}} \frac{\partial}{\partial q_{i}} . \tag{7.1.2}
\end{equation*}
$$

Now recall certain aspects of the Hamilton-Jacobi theory. We envision a mechanical system consisting of $n$ particles moving on a line. The space $\mathbb{R}^{2 n}$ is the set of all classical states, $p_{i}$ and $q_{i}$ are respectively the momentum and position of the $i$ th particle and $H=H\left(p_{1}, \ldots, p_{n}, q_{1}, \ldots, q_{n}\right) \in C^{\infty}\left(\mathbb{R}^{2 n}\right)$ is the total energy of the system. Assume $\xi_{H}$ is globally integrable so that one has the action on $\mathbb{R}^{2 n}$ of the one-parameter group $\exp t \xi_{H}$ of symplectic diffeomorphisms where if $\psi \in C^{\infty}\left(\mathbb{R}^{2 n}\right)$ and $z \in \mathbb{R}^{2 n}$ then

$$
\begin{align*}
{[\psi, H](z) } & =-\left(\xi_{H} \psi\right)(z)  \tag{7.1.3}\\
& =\left.\frac{d}{d t} \psi\left(\exp t \xi_{H} \cdot z\right)\right|_{t=0} .
\end{align*}
$$

That is, the map $t \rightarrow z(t)=\exp t \xi_{H} \cdot z$ is a trajectory of $-\xi_{H}$. In physical terms $z(t)$ is the state the mechanical system would occupy at time $t$ if it occupied the state $z$ at time $t=0$. Now it is immediate from (6.1.1) and (7.1.1) that

$$
\begin{equation*}
\xi_{q_{i}}=\frac{\partial}{\partial p_{i}} \quad \text { and } \quad \xi_{p_{i}}=-\frac{\partial}{\partial q_{i}} \tag{7.1.4}
\end{equation*}
$$

Now also regard $p_{i}$ and $q_{i}$, respectively, as the functions $t \rightarrow p_{i}(z(t)), t \rightarrow q_{i}(z(t)$, of $t$. It then follows from (7.1.3), recognizing that $[\psi, H]$ is also $\xi_{\psi} H$, that

$$
\begin{equation*}
\frac{\partial H}{\partial p_{i}}(z(t))=\frac{\partial q_{i}}{d t}, \quad \frac{\partial H}{\partial q_{i}}(z(t))=-\frac{\partial p_{i}}{d t} . \tag{7.1.5}
\end{equation*}
$$

These relations are of course just Hamilton's equations.
Now let $m_{i}$ be the mass of the $i$ th particle. We assume that $H$ takes the usual form as the sum of the kinetic and potential energies. That is, where $V \in C^{\infty}\left(\mathbb{R}^{2 n}\right)$ depends only on the $q_{i}$ (i.e., $\partial V / \partial p_{i}=0$ ), we assume that

$$
\begin{equation*}
H=\sum_{i=1}^{n} \frac{p_{i}{ }^{2}}{2 m_{i}}+V \tag{7.1.6}
\end{equation*}
$$

It follows then from (7.1.6) that

$$
\begin{equation*}
p_{i}(z(t))=m_{i} \frac{d q_{i}}{d t} \tag{7.1.7}
\end{equation*}
$$

so that the trajectory $z(t)$ is very simply determined as soon as one knows the velocities $d q_{i} / d t$.

Now let $\mathscr{R} \subseteq C^{\infty}\left(\mathbb{R}^{2 n}\right)$ be the $(2 n+1)$-dimensional space of functions spanned by the $p_{i}, q_{j}$, and the constant function 1 . One defines an alternating bilinear form $B_{\omega}$ on $\mathscr{R}$ by the relation

$$
\begin{equation*}
[r, s]=B_{\omega}(r, s) 1 \tag{7.1.8}
\end{equation*}
$$

for $r, s \in \mathscr{R}$. Let $\mathscr{P}, \mathscr{Q} \subseteq \mathscr{R}$ be respectively the $n$-dimensional subspaces spanned by the $p_{i}$ and $q_{j}$. It is clear that $\mathscr{P}$ and $\mathscr{Q}$ are totally singular subspaces which, however, are non-singularly paired with respect to $B_{\omega}$. In fact one readily has

$$
\begin{equation*}
B_{\omega}\left(q_{i}, p_{j}\right)=\delta_{i j} . \tag{7.1.9}
\end{equation*}
$$

We now assume that the potential function $V$ takes a certain special form. Let $l \leqslant n$ and assume that $\psi_{i} \in \mathscr{Q}, i=1, \ldots, l$, are linearly independent. That is, there exists an $l \times n$ matrix $A=\left(a_{i j}\right)$ of rank $l$ such that

$$
\begin{equation*}
\psi_{i}=\sum_{j=1}^{n} a_{i j} g_{j}, \quad i=1, \ldots, l . \tag{7.1.10}
\end{equation*}
$$

Let $r_{i} \in \mathbb{R}^{*}, i=1, \ldots, l$, be some positive constants. We assume that $V=$ $\sum_{i} r_{i} e^{\psi_{i}}$ so that the Hamiltonian $H$ has the form

$$
\begin{equation*}
H=\sum_{i=1}^{n} \frac{p_{i}^{2}}{2 m_{i}}+\sum_{i=1}^{l} r_{i} e^{\psi_{i}} . \tag{7.1.11}
\end{equation*}
$$

Let $z \in \mathbb{R}^{2 n}$ be arbitrary. We will consider the question of determining the trajectory $z(t)=\exp t \xi_{H} \cdot z$ when the $\psi_{i}$ satisfy a certain property-to be stated below.

Now $[q, H]$ is contained in $\mathscr{P}$ for any $q \in \mathscr{Q}$. In fact $\left[q_{i}, H\right]=p_{i} / m_{i}$ by (7.1.9) and hence

$$
\beta_{H}: \mathscr{Q} \rightarrow \mathscr{P}
$$

is a linear isomorphism where $\beta_{H}(q)=[q, H]$. One then defines a bilinear form $B_{H}$ on $\mathscr{Q}$ by putting

$$
\begin{equation*}
B_{H}\left(q, q^{\prime}\right)=B_{\omega}\left(q, \beta_{H} q^{\prime}\right) \tag{7.1.12}
\end{equation*}
$$

That is,

$$
\begin{equation*}
B_{H}\left(q, q^{\prime}\right) 1=\left[q,\left[q^{\prime}, H\right]\right] . \tag{7.1.13}
\end{equation*}
$$

It is clear then that $B_{H}$ is symmetric and positive definite. In fact one easily has

$$
\begin{equation*}
B_{H}\left(q_{i}, q_{j}\right)=\delta_{i j} / m_{i} \tag{7.1.14}
\end{equation*}
$$

Now let $\mathscr{Q}_{1} \subseteq \mathscr{Q}$ be the $l$-dimensional subspace spanned by the $\psi_{i}$ and put $\mathscr{P}_{1}=\beta_{H} \mathscr{Q}_{1}$. Since $\beta_{H}$ is an isomorphism one has $\operatorname{dim} \mathscr{P}_{1}=l$. Furthermore since $B_{H}$ is positive definite it follows from (7.1.12) that $\mathscr{Q}_{1}$ and $\mathscr{P}_{1}$ are nonsingularly paired by $B_{\omega}$. It follows that there exist uniquely $p_{i}^{\prime} \in \mathscr{P}_{1}, i=1, \ldots, l$, such that $B_{\omega}\left(\psi_{i}, p_{j}^{\prime}\right)=\delta_{i j}$. Thus if we define $q_{i}^{\prime} \in \mathscr{R}, i=1, \ldots, l$, by putting

$$
\begin{equation*}
q_{i}^{\prime}=\psi_{i}+\log r_{i} 1 \tag{7.1.15}
\end{equation*}
$$

then for $i, j=1, \ldots, l$

$$
\begin{equation*}
\left[q_{i}^{\prime}, p_{j}^{\prime}\right]=\delta_{i j} 1 \tag{7.1.16}
\end{equation*}
$$

Now let $\mathscr{Q}_{2}$ be the orthocomplement of $\mathscr{Q}_{1}$ in $\mathscr{Q}$ with respect to $B_{H}$ and put $\mathscr{P}_{2}=\beta_{H} \mathscr{Q}_{2}$. Obviously one has the direct sum

$$
\begin{equation*}
\mathscr{P}=\mathscr{P}_{1}+\mathscr{P}_{2} . \tag{7.1.17}
\end{equation*}
$$

Now let $q_{i}^{\prime}, i=l+1, \ldots, n$, be a $B_{H}$ orthonormal basis of $\mathscr{Q}_{2}$. Again, $\mathscr{P}_{2}$ is non-singularly paired to $\mathscr{Q}_{2}$ by $B_{\omega}$ and hence there exist uniquely $p_{i}^{\prime} \in \mathscr{P}_{2}$, $i=l+1, \ldots, n$, such that $\left[q_{i}^{\prime}, p_{j}^{\prime}\right]=\delta_{i j} 1$. But now clearly $\mathscr{P}_{2}$ and $\mathscr{Q}_{1}$ are each other's orthocomplements with respect to $B_{\omega}$. By symmetry the same statement is true for $\mathscr{P}_{1}$ and $\mathscr{Q}_{2}$. Thus $p_{i}^{\prime}, q_{j}^{\prime}, i, j=1, \ldots, n$, is a coordinate system in $\mathbb{R}^{2 n}$ and

$$
\begin{equation*}
\omega=\sum_{i=1}^{l} d p_{i}^{\prime} \wedge d q_{i}^{\prime}+\sum_{j=l+1}^{n} d p_{j}^{\prime} \wedge d q_{j}^{\prime} . \tag{7.1.18}
\end{equation*}
$$

Now the symmetric $l \times l$ matrix $B=\left(b_{i j}^{\prime}\right)$ defined by

$$
\begin{equation*}
b_{i j}^{\prime}=B_{H}\left(\psi_{i}, \psi_{j}\right) \tag{7.1.19}
\end{equation*}
$$

will play an important role for us. We first observe

Proposition 7.1. Let $H$ be the Hamiltonian given by (7.1.11). Then with respect to the coordinates $p_{i}^{\prime}, q_{j}^{\prime}, i, j=1, \ldots, n$, one has

$$
\begin{equation*}
H=H_{1}+\sum_{i=l+1}^{n} \frac{\left(p_{i}^{\prime}\right)^{2}}{2} \tag{7.1.20}
\end{equation*}
$$

where

$$
\begin{equation*}
H_{1}=\frac{1}{2} \sum_{i, j=1}^{l} b_{i j}^{\prime} p_{i}^{\prime} p_{j}^{\prime}+\sum_{i=1}^{l} e^{q_{i}^{\prime}} . \tag{7.1.21}
\end{equation*}
$$

Proof. Obviously $e^{q_{i}^{\prime}}=r_{i} e^{L_{i}}, i=1, \ldots, l$. It suffices then to prove that if

$$
\begin{equation*}
K=\frac{1}{2} \sum_{i, j=1}^{l} b_{i j}^{\prime} p_{i}^{\prime} p_{j}^{\prime}+\sum_{i=l+1}^{n} \frac{\left(p_{i}^{\prime}\right)^{2}}{2} \tag{7.1.22}
\end{equation*}
$$

then $K$ is the kinetic energy in (7.1.11). Since both $K$ and the kinetic energy are in the symmetric subspace $S^{2}(\mathscr{P})$ it suffices to show that $B_{K}=B_{H}$, where $B_{K}$ is the symmetric bilinear form on $\mathscr{2}$ given by $B_{K}\left(q, q^{\prime}\right) 1=\left[q,\left[q^{\prime}, K\right]\right]$. But now evidently $B_{K}\left(\psi_{i}, \psi_{j}\right)=b_{i j}^{\prime}=B_{H}\left(\psi_{i}, \psi_{j}\right)$ by (7.1.15) and (7.1.16). Thus $B_{K}$ and $B_{H}$ agree on $\mathscr{Q}_{1}$. Also $\mathscr{Q}_{1}$ and $\mathscr{Q}_{2}$ are $B_{K}$ orthogonal and $q_{i}^{\prime}$, $i=l+1, \ldots, n$, are $B_{K}$ orthonormal. Thus $B_{H}=B_{K}$.
Q.E.D.
7.2. Now put $Z^{\prime}=\mathbb{R}^{2 l}$ and let $\rho_{i}^{\prime}, \varphi_{j}^{\prime} \in C^{\infty}\left(Z^{\prime}\right)$, where $i, j=1, \ldots, l$, be a linear coordinate system. Then ( $Z^{\prime}, \omega_{Z^{\prime}}$ ) is just $2 l$ phase space with the $\rho_{i}^{\prime}$ and $\varphi_{j}^{\prime}$ as canonical coordinates if we put

$$
\begin{equation*}
\omega_{Z^{\prime}}=\sum_{i=1}^{l} d \rho_{i}^{\prime} \wedge d \varphi_{i}^{\prime} \tag{7.2.1}
\end{equation*}
$$

Now let $I_{1}^{\prime} \in C^{\infty}\left(Z^{\prime}\right)$ be the function defined by putting

$$
\begin{equation*}
I_{1}^{\prime}=\frac{1}{2} \sum_{i, j=1}^{l} b_{i j}^{\prime} \rho_{i}^{\prime} \rho_{j}^{\prime}+\sum_{i=1}^{l} e^{\varphi_{i}^{\prime}}, \tag{7.2.2}
\end{equation*}
$$

where $b_{i j}^{\prime}$ is given by (7.1.18). Also let

$$
\begin{equation*}
\delta: \mathbb{R}^{2 n} \rightarrow Z^{\prime} \tag{7.2.3}
\end{equation*}
$$

be the unique smooth map such that $\rho_{i}^{\prime} \circ \delta=p_{i}^{\prime}$ and $\varphi_{i}^{\prime} \circ \delta=q_{i}^{\prime}, i=1, \ldots, l$. Finally let

$$
\begin{equation*}
\Gamma_{2}: \mathscr{P} \rightarrow \mathscr{P}_{2} \tag{7.2.4}
\end{equation*}
$$

be the projection defined by (7.1.17) so that $\mathscr{P}_{1}=\operatorname{Ker} \Gamma_{2}$.
Now, recalling the notation of Section $6.1, \xi_{H}$ and $\xi_{I_{1}^{\prime}}$ are respectively the Hamiltonian vector fields on $\mathbb{R}^{2 n}$ and $Z^{\prime}$ corresponding to the functions $H \in$ $C^{\infty}\left(\mathbb{R}^{2 n}\right)$ and $I_{1}^{\prime} \in C^{\infty}\left(Z^{\prime}\right)$. We recall also that $A$ is the $l \times n$ matrix ( $a_{i j}$ ) defined in (7.1.10). In addition one notes that since the $\psi_{i}$ are linearly independent the $l \times l$ matrix $B$ defined by (7.1.19) is invertible and hence if $A^{*}$ is the transpose of $A$ then $A^{*} B^{-1}$ is an $n \times l$ matrix.

Proposition 7.2. Let $H$ be the Hamiltomian function on $\mathbb{R}^{2 n}$ given by (7.1.11) and let $I_{1}^{\prime} \in C^{\infty}\left(Z^{\prime}\right)$ be given by (7.2.2). Then $\xi_{H}$ is globally integrable on $\mathbb{R}^{2 n}$ in case $\xi_{r_{1}^{\prime}}$ is globally integrable on $Z^{\prime}$. Furthermore in such a case if $x \in \mathbb{R}^{2 n}$ is an arbitrary initial point and, for notational convenience, we put $q_{j}(t)=q_{j}\left(\exp t \xi_{H} \cdot x\right)$, $j=1, \ldots, n$, then

$$
\begin{equation*}
q_{j}(t)-q_{j}(0)=\frac{1}{m_{j}}\left(c_{j} t+\sum_{k=1}^{l}\left(A^{*} B^{-1}\right)_{j k}\left(\varphi_{k}^{\prime}(t)-\varphi_{k}^{\prime}(0)\right)\right), \tag{7.2.5}
\end{equation*}
$$

where $c_{j}$ is the constant $\left(\Gamma_{2} p_{j}\right)(x)$ (see (7.2.4)) and letting $y^{\prime} \in Z^{\prime}$ be defined by $\delta x=y^{\prime}\left(\right.$ see (7.2.3)) one defines $\varphi_{k}^{\prime}(t)$ by putting $\varphi_{k}^{\prime}(t)=\varphi_{k}^{\prime}\left(\exp t \xi_{I_{1}^{\prime}} \cdot y^{\prime}\right)$.

Proof. Assume $\xi_{I_{1}^{\prime}}$ is globally integrable. Let $y(t)=\exp t \xi_{I_{1}^{\prime}} \cdot y^{\prime}$ and let $\mathbb{B} \rightarrow \mathbb{R}^{2 n}, t \rightarrow x(t)$ be the curve in $\mathbb{R}^{2 n}$ defined by the relation $p_{i}^{\prime}(x(t))=$ $\rho_{i}^{\prime}(y(t)), q_{i}^{\prime}(x(t))=\varphi_{i}^{\prime}(y(t)), i=1, \ldots, l, p_{k}^{\prime}(x(t))=p_{k}^{\prime}(x)$, and $q_{k}^{\prime}(x(t))=q_{k}^{\prime}(x)$, $k=l+1, \ldots, n$. But now recalling (7.2.3), (7.1.21), and (7.2.2), one obviously has $I_{1}^{\prime} \circ \delta=H_{1}$. But since Hamilton's equations are satisfied for the function $I_{1}^{\prime}$ along the curve $y(t)$ with respect to the $\varphi_{i}^{\prime}$ and $\rho_{j}^{\prime}$ it follows immediately that Hamilton's equations are satisfied for the function $H_{1} \in C^{\infty}\left(\mathbb{R}^{2 n}\right)$ along the curve $x(t)$ with respect to the $p_{i}^{\prime}$ and $q_{j}^{\prime}$. It follows then that $\xi_{H_{1}}$ is globally integrable and

$$
\begin{equation*}
x(t)=\exp t \xi_{H_{1}} \cdot x \tag{7.2.6}
\end{equation*}
$$

Now let $H_{2}=\sum_{i=l+1}^{n}\left(\left(p_{i}^{\prime}\right)^{2} / 2\right)$. It is obvious that $\xi_{H_{2}}$ is globally integrable where $p_{i}^{\prime}\left(\exp t \xi_{H_{2}} \cdot x\right)=p_{i}^{\prime}(x), i=1, \ldots, n, q_{i}^{\prime}\left(\exp t \xi_{H_{2}} \cdot x\right)=q_{i}^{\prime}(x)$ for $i \leqslant l$ and $q_{i}^{\prime}\left(\exp t \xi_{H_{2}} \cdot x\right)=q_{i}^{\prime}(x)+t p_{i}^{\prime}(x), i=l+1, \ldots, n$. But $H=H_{1}+H_{2}$ and clearly $\left[H_{1}, H_{2}\right]=0$. Thus the vector fields $\xi_{H_{1}}$ and $\xi_{H_{2}}$ commute and since $\xi_{H}=\xi_{H_{1}}+\xi_{H_{2}}$ it follows, as one knows, that $\xi_{H}$ is globally integrable where

$$
\begin{equation*}
\exp t \xi_{H} \cdot x=\exp t \xi_{H_{2}} \cdot\left(\exp t \xi_{H_{1}} \cdot x\right) \tag{7.2.7}
\end{equation*}
$$

Now for any $f \in C^{\infty}\left(\mathbb{R}^{2 n}\right)$ and $t \in \mathbb{R}$ put $f(t)=f\left(\exp t \xi_{H} \cdot x\right)$. From (7.1.7) one then clearly has the integral

$$
\begin{equation*}
q_{j}(t)-q_{j}(0)=\frac{1}{m_{j}} \int_{0}^{t} p_{j}(s) d s \tag{7.2.8}
\end{equation*}
$$

Now for any $p \in \mathscr{P}_{2}$ one has $\left[p, \psi_{i}\right]=0, i=1, \ldots, l$, since $\mathscr{P}_{2}$ is $B_{\omega}$ orthogonal to $\mathscr{Q}_{1}$. Thus clearly $[H, p]=\xi_{H} p=0$. Thus $p(t)$ is a constant function of $t$. But $\Gamma_{2} p_{j} \in \mathscr{P}_{2}$. Hence

$$
\begin{equation*}
\frac{1}{m_{j}} \int_{0}^{t}\left(\Gamma_{2} p_{j}\right)(s) d s=\frac{c_{j} t}{m_{j}} \tag{7.2.9}
\end{equation*}
$$

where $c_{j}=\left(\Gamma_{2} p_{j}\right)(x)$.

For any $g \in C^{\infty}\left(Z^{\prime}\right)$ and $t \in \mathbb{R}$ let $g(t)=g\left(\exp t \xi_{I_{1}^{\prime}} \cdot y^{\prime}\right)$, where recalling (7.2.3) one puts $\delta x=y^{\prime} \in Z^{\prime}$.

Now let $\Gamma_{1}=I-\Gamma_{2}$, where $I$ is the identity operator on $\mathscr{P}$ so that $\Gamma_{1}$ : $\mathscr{P} \rightarrow \mathscr{P}_{1}$ is the projection whose kernel is $\mathscr{P}_{2}$. But $p_{j}=\Gamma_{1} p_{j}+\Gamma_{2} p_{j}$. Thus to prove (7.2.5) it suffices by (7.2.8) and (7.2.9) to show that

$$
\begin{equation*}
\sum_{k=1}^{l}\left(A^{*} B^{-1}\right)_{j_{k}}\left(\varphi_{k}^{\prime}(t)-\varphi_{k}^{\prime}(0)\right)=\int_{0}^{t}\left(\Gamma_{\mathbf{1}} p_{j}\right)(s) d s \tag{7.2.10}
\end{equation*}
$$

Now complete the $l \times n$ matrix $A$ to an $n \times n$ matrix $\left(a_{i j}\right)$, where for $i>l$ one defines $a_{i j}$ by the relation $\sum_{j=1}^{n} a_{i j} q_{j}=q_{i}^{\prime}$. Now recalling (7.1.18) one then easily has (taking the transpose with respect to $B_{\omega}$ ) the relation

$$
\begin{equation*}
\sum_{i-1}^{n} a_{i j} p_{i}^{\prime}=p_{j} \tag{7.2.11}
\end{equation*}
$$

for all $1 \leqslant i \leqslant n$. But now $\Gamma_{1} p_{i}^{\prime}=0$ for $i>l$ and $\Gamma_{1} p_{i}^{\prime}=p_{i}^{\prime}$ for $i \leqslant l$. Thus applying $\Gamma_{1}$ to both sides of (7.2.11) one just replaces $n$ by $l$. That is,

$$
\begin{equation*}
\sum_{i=1}^{l} a_{i j} p_{i}^{\prime}=\Gamma_{1} p_{j} \tag{7.2.12}
\end{equation*}
$$

Now let $\mu_{i}, i=1, \ldots, l$, be the basis of $\mathscr{P}_{1}$ defined by putting $\mu_{i}=\beta_{H}\left(\psi_{i}\right)=$ $\left[\psi_{i}, H\right]=\left[q_{i}^{\prime}, H\right]$. Thus $\mu_{i}=-\xi_{H} q_{i}^{\prime}$. But now, recalling the first paragraph in this proof, one has $q_{i}^{\prime}(t)=\varphi_{i}^{\prime}(t)$. Thus

$$
\begin{equation*}
\varphi_{i}^{\prime}(t)-\varphi_{i}^{\prime}(0)=\int_{0}^{t} \mu_{i}(s) d s \tag{7.2.13}
\end{equation*}
$$

But now the $\mu_{i}$ and $\boldsymbol{p}_{j}^{\prime}$, where $i, j=1, \ldots, l$, are both bases of $\mathscr{P}_{1}$. We assert that

$$
\begin{equation*}
\mu_{i}=\sum_{j=1}^{l} b_{i j}^{\prime} p_{j}^{\prime} . \tag{7.2.14}
\end{equation*}
$$

Indeed if $s_{i j}$ is defined by the relation $\mu_{i}=\sum_{j} s_{i j} p_{j}^{\prime}$ then $B_{\omega}\left(\psi_{j}, \mu_{i}\right)=s_{i j}$ since $B_{\omega}\left(\psi_{j}, p_{i}^{\prime}\right)=\delta_{i j}$. However, $\mu_{i}=\beta_{H} \psi_{i}$. Thus $s_{i j}=B_{\omega}\left(\psi_{j}, \beta_{H} \psi_{i}\right)=$ $B_{H}\left(\psi_{j}, \psi_{i}\right)=b_{i j}^{\prime}$. This proves (7.2.14). But then since $B$ is the $l \times l$ matrix ( $b_{i j}^{\prime}$ ) one has $p_{i}^{\prime}=\sum_{k=1}^{l}\left(B^{-1}\right)_{i k} \mu_{k}$. Substituting in (7.2.12) one then has

$$
\begin{equation*}
\sum_{k=1}^{l}\left(A^{*} B^{-1}\right)_{j k} \mu_{k}=\Gamma_{\mathbf{1}} p_{j} \tag{7.2.15}
\end{equation*}
$$

But now evaluating both sides along the curve $\exp s \xi_{H} \cdot x$ and integrating from 0 to $t$ one obtains (7.2.10) from (7.2.13).
Q.E.D.
7.3. We now recall the definition of a Cartan matrix. (The definition will be made for the semi-simple and not just for the simple case.) See [11, Chap. 11] for more details. We remark first that the Cartan matrices classify the complex semi-simple Lie algebras or the compact semi-simple Lie algebra or, as we shall take it here, the real split semi-simple Lie algebras. Assume that $g$ is a real semi-simple Lie algebra, say, of rank $l$, which admits a Cartan subalgebra $\ell \subseteq g$ such that the eigenvalues of ad $x$ are real for all $x \in K$. Then $g$ is a real split semi-simple Lie algebra and $\hbar$ is a split Cartan subalgebra. This in fact has been our assumption in earlier sections. Let $\ell^{\prime}$ be the dual space to $h$ and let ( $\mu, \nu$ ) be a Weyl group invariant positive definite inner product on $\ell^{\prime}$. Let $\alpha_{i} \in \ell^{\prime}, i=1, \ldots, l$, be the set of simple positive roots relative to some lexicographical ordering in $\ell^{\prime}$. Then the $l \times l$ matrix $C$ defined by

$$
\begin{equation*}
C_{i j}=\frac{2\left(\alpha_{i}, \alpha_{j}\right)}{\left(\alpha_{j}, \alpha_{j}\right)} \tag{7.3.1}
\end{equation*}
$$

is called a Cartan matrix. One knows that the matrix is independent of the bilinear form (as long as it is Weyl group invariant) and the entries $C_{i j}$ are integers. Furthermore the matrix completely determines the structure of $g$. Now assume that $V$ is some finite-dimensional real vector space with a positive definite inner product $P$. Let $v_{1}, \ldots, v_{l} \in V$ be linearly independent vectors. Then if $2 P\left(v_{i}, v_{j}\right) / P\left(v_{j}, v_{j}\right)=C_{i j}^{\prime}$ is a Cartan matrix there thus exist uniquely up to isomorphism a split semi-simple Lie algebra $g$ of rank $l$, a split Cartan subalgebra $月$, and simple roots $\alpha_{i}$ such that one has $C_{i j}=C_{i j}^{\prime}$. Also if $C_{i j}^{\prime}$ is a Cartan matrix one may introduce a diagram, the Dynkin diagram, to describe the angles and relative lengths of the vectors $v_{i}$. The Dynkin diagram is based on the fact that if $v_{i}$ and $v_{j}$ are not orthogonal and $P\left(v_{i}, v_{i}\right) \geqslant P\left(v_{j}, v_{j}\right)$ then

$$
\begin{equation*}
P\left(v_{i}, v_{i}\right) / P\left(v_{j}, v_{j}\right)=1,2 \text {, or } 3 \tag{7.3.2}
\end{equation*}
$$

and accordingly the angle between $v_{i}$ and $v_{j}$ is 120,135 , or $150^{\circ}$. (See, e.g., [5].) With regard to terminology in this paper, where the situation is warranted we will either say that $C_{i j}^{\prime}$ is a Cartan matrix or that $v_{i}, i=1, \ldots, l$, defines a Dynkin diagram. In either case we will refer to $g, \hbar$, and the $\alpha_{i}$ as a corresponding split semi-simple Lie algebra, a split Cartan subalgebra, and a set of simple positive roots.

We recall in Section 1.2 we permitted some flexibility in the definition of the invariant bilinear form $Q$ on $g$. This may now be normalized according to $P$. As usual $Q$ is "carried" over to $\ell^{\prime}$ using the isomorphism $\ell \rightarrow \ell^{\prime}$ defined by $Q \mid h$.

Proposition 7.3. Assume that $v_{1}, \ldots, v_{l} \in V$ defines a Dynkin diagram with respect to some positive definite inner product $P$ on $V$. Let $g, h$, and $\alpha_{i}, \ldots, \alpha_{l}$
be a corresponding split semi-simple Lie algebra, a Cartan subalgebra, and a set of simple positive roots. Then there exists a unique invariant bilinear form $Q$ on $g$ which on each simple component of $g$ is a positive multiple of the Killing form such that

$$
\begin{equation*}
Q\left(\alpha_{i}, \alpha_{j}\right)=P\left(v_{i}, v_{j}\right) \tag{7.3.3}
\end{equation*}
$$

for $i, j=1, \ldots, l$.
Proof. Recalling the notion of connectedness of a Dynkin diagram (or Coxeter graphs-see, e.g., Section 11.3 in [11]) one notes from (7.3.2) that the ratios $r_{i j}=P\left(v_{i}, v_{i}\right) / P\left(v_{j}, v_{j}\right)$ are uniquely determined in case $\alpha_{i}$ and $\alpha_{j}$ are in the same connected component. In such a case one then clearly has $r_{i j}=$ $Q\left(\alpha_{i}, \alpha_{i}\right) / Q\left(\alpha_{j}, \alpha_{j}\right)$ for any choice of $Q$. However, as one knows, the connected components of the diagram correspond to the simple components of $g$. We may thus uniquely normalize $Q$ so that $Q\left(\alpha_{i}, \alpha_{i}\right)=P\left(\alpha_{i}, \alpha_{i}\right)$ for any $i$. One notes here that replacing $Q$ on a simple component of $g$ by a positive multiple $\lambda Q$ is equivalent to replacing $Q$ on the corresponding subspace of $\boldsymbol{h}^{\prime}$ by the positive multiple $Q / \lambda$.

One then has (7.3.3) from the equality $C_{i j}^{\prime}=2 Q\left(\alpha_{i}, \alpha_{j}\right) / Q\left(\alpha_{j}, \alpha_{j}\right)$. Q.E.D.
7.4. Now returning to the Hamiltonian $H$ of (7.1.11) we now assume that the $\psi_{i} \in \mathscr{Q}, i=1, \ldots, l$, define a Dynkin diagram with respect to $B_{H}$. Before proceeding we will give some examples where this is the case.
(1) $l=n-1$. All the masses are equal and $\psi_{i}=q_{i}-q_{i+1}, i=1, \ldots, n-1$. Then $H$ is the Hamiltonian of the usual nonperiodic Toda lattice. The potential here is that of "nearest-neighbor" particle interaction. The Dynkin diagram is that of $A_{n-1}$, using standard notation so that $g$ can be taken to be the Lie algebra of $S l(n, \mathbb{R})$ and $h$ is the space of traceless real diagonal matrices.
(2) $l=n$. All the masses are equal and $\psi_{i}=q_{i}-q_{i+1}, i=1, \ldots, n-1$. But $\psi_{n}=q_{n}$. The potential is similar to the one above except that the last particle may also be regarded as interacting with a fixed mass. The Dynkin diagram is of $B_{n}$ so that $g$ is the Lie algebra of $S O(n, n+1)$.
(3) Similar to (2) except that $\psi_{n}=2 q_{n}$. The last particle interacts even more strongly with a fixed mass. Hence the Dynkin diagram is that of $C_{n}$ and $g$ is the Lie algebra of the symplectic group $\operatorname{Sp}(2 n, \mathbb{R})$. In Section 7.8 we will apply Theorem 7.5 below to solve the three-body problem (the case $n=3$ here) given by

$$
\begin{equation*}
H=\sum_{i=1}^{3} \frac{p_{i}{ }^{2}}{2}+e^{q_{1}-q_{2}}+e^{q_{2}-q_{3}}+e^{2 q_{3}} \tag{7.4.1}
\end{equation*}
$$

(4) $l=n$. All the masses are equal. The case is that of (1) except that in
addition $\psi_{n}=q_{n-1}+q_{n}$. Thus the potential is similar to (1) except that the center of mass of the last two particles interacts with a fixed mass. Here the Dynkin diagram is that of $D_{n}$ and $g$ can be taken to be the Lie algebra of $S O(n, n)$.

We wish to emphasize that there are of course an infinite number of different dynamical systems which correspond to the same Dynkin diagram. In (1)-(4) above we have just given one system for each of the classical Lie algebras. A particularly interesting Dynkin diagram is that of $D_{4}$. The diagram is

and as the figure somewhat suggests one may give a system (different from that in (4) above) where $l=n=4$, so that there are four particles, three of which are interacting with a fixed body and whose center of gravity is interacting with the fourth particle. The Hamiltonian is

$$
\begin{equation*}
H=\sum_{i=1}^{4} \frac{p_{i}{ }^{2}}{2}+e^{q_{1}}+e^{q_{2}}+e^{q_{3}}+e^{\left(q_{4}-q_{1}-q_{2}-q_{3}\right) / 2} . \tag{7.4.3}
\end{equation*}
$$

Remark 7.4.1. For a classical Lie algebra one may solve the system by considering the usual Toda lattice for a larger number of variables. This is no longer true for the exceptional Lie algebras $F_{4}, E_{6}, E_{7}$, and $E_{8}$. In these cases if $z \in \mathcal{g}$ is a Jacobi element there is no representation $\pi$ such that $\pi(z)$ is a Jacobi operator in the usual sense, as one may show. An example of a dynamical system corresponding to $F_{4}$ is the four-body problem given by

$$
\begin{equation*}
H=\sum_{i=1}^{4} \frac{p_{i}^{2}}{2}+e^{q_{1}-q_{2}}+e^{\sigma_{2}-q_{3}}+e^{\left(\sigma_{4}-q_{1}-q_{2}-q_{3}\right) / 2} . \tag{7.4.4}
\end{equation*}
$$

Among all the others then Theorem 7.5 below will thus describe the trajectories of the systems (7.4.3) and (7.4.4) (in terms of the four fundamental representations of $D_{4}$ and $F_{4}$, respectively).

Remark 7.4.2. Even though our assumption concerning the $\psi$ 's is very special we wish to note that there is some latitude in satisfying the assumption. Namely, given the $\psi_{i}$ one varies the inner product $B_{H}$ within an $n$-dimensional variety of such inner products by varying the masses $m_{i}$. This raises the possibility of satisfying the assumption for a certain choice of the masses $m_{i}$. More precisely recalling that $b_{i j}^{\prime}=B_{H}\left(\psi_{i}, \psi_{j}\right)$ our assumption about the $\psi_{i}$ is that $C^{\prime}$ should be a Cartan matrix where $C_{i j}^{\prime}=b_{i j}^{\prime} \mid b_{i j}^{\prime}$. This is of course just a statement about the matrix $B$ where, as in Proposition $7.2, B=\left(b_{i j}^{\prime}\right)$. Thus given the $\psi_{i}$, that is, given the $l \times n$ matrix $A$ (see (7.1.10)) the following proposition describes how $B$ varies with the masses $m_{i}$.

Proposition 7.4.1. Let $A$ be respectively the $l \times n$ matrix defined as in
(7.1.10) and let $A^{*}$ be the transpose of $A$. Let $M$ be the $n \times n$ diagonal matrix where $M_{i i}=1 / m_{i}$ and let $B$, as in Proposition 7.2 , be the $l \times l$ matrix defined by putting $B_{i j}=B_{H}\left(\psi_{i}, \psi_{j}\right)$. Then

$$
\begin{equation*}
B=A M A^{*} \tag{7.4.5}
\end{equation*}
$$

Proof. One has for $i, j=1, \ldots, l$,

$$
\begin{equation*}
B_{H}\left(\psi_{i}, \psi_{j}\right)=B_{H}\left(\sum_{r=1}^{n} a_{i r} q_{r}, \sum_{s-1}^{l} a_{j s} q_{s}\right) . \tag{7.4.6}
\end{equation*}
$$

But then (7.4.5) follows from the relation $B_{H}\left(q_{r}, q_{s}\right)=\delta_{r s} / m_{r}$. See (7.1.14). Q.E.D.

Now let $g, h$, and $\alpha_{i}, i=1, \ldots, l$, correspond to the Dynkin diagram of the $\psi_{i}$. Thus $g$ is real split semi-simple Lie algebra, $h$ is a split Cartan subalgebra, and $\alpha_{i} \in \ell^{\prime}$ are the simple roots for a fixed system $\Delta_{+}$of positive roots relative to $(g, h)$. We use the notation and results of Sections 6 and 7.1.3. By Proposition 7.3 we may fix the bilinear form $Q$ in $g$ so that

$$
\begin{equation*}
Q\left(\alpha_{i}, \alpha_{j}\right)=B_{H}\left(\psi_{i}, \psi_{j}\right) \tag{7.4.7}
\end{equation*}
$$

for $i, j=1, \ldots, l$. That is, recalling (6.5.1) and (7.1.19) one has $b_{i j}^{\prime}=b_{i j}$ and hence the matrix $B=\left(b_{i j}^{\prime}\right)$ is just given by $B=\left(b_{i j}\right)=\left(Q\left(\alpha_{i}, \alpha_{j}\right)\right)$.

Now recall (see Sections 6.3 and 6.4 ) that $\left(Z, \omega_{Z}\right)$ is a $2 l$-dimensional symplectic manifold where $Z \subseteq g$ is the manifold of normalized Jacobi elements on $g$ and $\omega_{Z}$ is defined in Section 6.4. We recall also (see Proposition 6.5) that $\rho_{i}, \varphi_{j}$, $i, j=1, \ldots, l$, are a global coordinate system on $Z$ and $\omega_{Z}$ is explicitly given by (6.5.4). In addition we recall that the restriction $I_{1} \mid Z$ of the fundamental invariant $I_{1} \in S(g)^{G}$ is given by (6.5.5). See (6.5.2).

Now let $\left(Z^{\prime}, \omega_{Z}\right), \rho_{i}^{\prime}, \varphi_{3}^{\prime}$, and $I_{1}^{\prime}$ be as in Section 7.2.

Proposition 7.4.2. There exists a unique symplectic isomorphism

$$
\begin{equation*}
\sigma: Z^{\prime} \rightarrow Z \tag{7.4.8}
\end{equation*}
$$

of $\left(Z^{\prime}, \omega_{Z^{\prime}}\right)$ and $\left(Z, \omega_{Z}\right)$ such that $\rho_{i} \circ \sigma=\rho_{i}^{\prime}, \varphi_{j} \circ \sigma=\varphi_{j}^{\prime}$. Moreover one has $I_{1} \circ \sigma=I_{1}^{\prime}$.

Proof. Since $b_{i j}=b_{i j}^{\prime}$ this is immediate from a comparison of (7.2.1) and (7.2.2) with (6.8.1) and (6.8.2).
Q.E.D.

Now let $A, B$, and $M$ be respectively the $l \times n, l \times l$, and $n \times n$ matrices
defined, say, as in (7.1.10), (7.1.19), and Proposition 7.4.1. Obviously $B$ is invertible (since the $\psi_{i}$ are linearly independent). Put

$$
\begin{equation*}
D=B^{-\mathbf{1}} A M . \tag{7.4.9}
\end{equation*}
$$

Proposition 7.4.3. The $l \times n$ matrix $D$ is characterized by the properties (1)

$$
\begin{equation*}
D A^{*}=1_{l} \tag{7.4.10}
\end{equation*}
$$

where $1_{l}$ is the $l \times l$ identity matrix and (2) $A^{*} D$ is the transpose of the $n \times n$ matrix corresponding to the projection $\Gamma_{1}: \mathscr{P} \rightarrow \mathscr{P}_{1}$, relative to the basis $p_{i}$ of $\mathscr{P}$. That is, if $\Gamma_{2}: \mathscr{P} \rightarrow \mathscr{P}_{2}$ is the complementary projection, as in (7.2.4), then for $j=1, \ldots, n$

$$
\begin{equation*}
p_{j}-\sum_{k=1}^{n}\left(A^{*} D\right)_{j k} p_{k}=\Gamma_{2} p_{j} . \tag{7.4.11}
\end{equation*}
$$

Proof. One has (7.4.10) immediately from (7.4.5). Now recall (7.1.10). We recall also (see the proof of Proposition 7.2) that $\mu_{k}=\left[\psi_{k}, H\right]=$ $\left[\sum_{r=1}^{n} a_{k r} q_{r}, H\right]$. But $\left[q_{r}, H\right]=p_{r} / m_{r}$. Thus

$$
\begin{equation*}
\mu_{k}=\sum_{r=1}^{n} a_{k r} \frac{p_{r}}{m_{r}} \tag{7.4.12}
\end{equation*}
$$

Now substituting (7.4.12) in (7.2.15) one has $\Gamma_{1} p_{j}=\sum_{r=1}^{n}\left(A^{*} B^{-1} A M\right)_{j r} p_{r}$. But $B^{-1} A M=D$. This proves (7.4.11). Since $A^{*} D$ is then given the $l \times n$ matrix $D$ is uniquely characterized, as one knows, from the additional relation $D A^{*}=1_{l}$.
Q.E.D.

Now Proposition 7.2 enables us now to put into effect Theorem 6.8.1 to assert that $\xi_{H}$ is globally integrable and to determine the trajectory $\exp t \xi_{H} \cdot x$ for any $x$ in phase space $\mathbb{R}^{2 n}$. What of course is needed is the Jacobi element (see (7.4.8)) $\sigma \delta x=y \in Z$. This plays a critical role in more than one way. One notes that the corresponding trajectory in $Z$ is a curve on the isospectral leaf $Z(\gamma) \subseteq Z$, where $\gamma=\mathscr{I}(y)$. (See (6.8.3) and (3.6.2).) The following gives an explicit formula for $y$ in terms of the position $q_{i}(x)$ and momentum $p_{j}(x)$ coordinates of $x$.

Proposition 7.4.1. Let $x \in \mathbb{R}^{2 n}$ and let $Z \subseteq g$ be as in (2.2.3) and (2.2.4) the space of normalized Jacobi elements. Let $\delta: \mathbb{R}^{2 n} \rightarrow Z^{\prime}$ and let $\sigma: Z^{\prime} \rightarrow Z$ be defined as in (7.2.3) and (7.4.8), respectively. Put $y=\sigma \delta x$. Let $h_{i} \in h$ be the basis of $h$ defined as in (6.4.1) and let $f \in g$ be as in (1.5.4). Let $e_{\alpha_{i}}, i=1, \ldots, l$, be the simple positive root vectors defined as in Section 1.5 and let $\rho_{i}, \varphi_{j}, i, j=$
$1, \ldots, l$, be the coordinate system on $Z$ defined as in Proposition 6.5 so that y can be written

$$
\begin{equation*}
y=f+\sum_{i=1}^{l} \rho_{i}(y) h_{i}+\sum_{j=1}^{l} e^{\sigma_{j}(v)} e_{\alpha_{j}} . \tag{7.4.13}
\end{equation*}
$$

Then if $A=\left(a_{i j}\right)$ and $D=\left(d_{i j}\right)$ are the $l \times n$ matrices defined by (7.1.10) and (7.4.9), respectively, and $r_{j}$ is defined as in (7.1.11) one has

$$
\begin{equation*}
\varphi_{j}(y)=\sum_{k=1}^{n} a_{j k} q_{k}(x)+\log r_{j} \tag{7.4.14}
\end{equation*}
$$

and

$$
\begin{equation*}
\rho_{i}(y)=\sum_{k=1}^{n} d_{i k} p_{k}(x) . \tag{7.4.15}
\end{equation*}
$$

Proof. Recalling (7.2.3) and (7.4.8) one has $\varphi_{j} \circ \sigma \delta=q_{j}^{\prime}$ so that $\varphi_{j}(y)=$ $q_{i}^{\prime}(x)$. But $q_{j}^{\prime}=\psi_{j}+\log r_{j}$. The formula (7.4.14) then follows from (7.1.10). Now by (7.2.3) and (7.4.8) one has $p_{i}^{\prime}=p_{i} \circ \sigma \delta$. Thus $p_{i}^{\prime}(x)=\rho_{i}(y)$. But now by (7.2.14) one has $p_{i}^{\prime}=\sum_{k=1}^{i}\left(B^{-1}\right)_{i k} \mu_{k}$. Recalling (7.4.12) one then has $p_{i}^{\prime}=\sum_{k=1}^{n} d_{i k} p_{k}$ since $D=B^{-1} A M$. This proves 7.4.15.
Q.E.D.
7.5. Now recalling Section 3.3 let $h_{+} \subseteq h$ be the open Weyl chamber defined as in (3.3.1) and let $H$ be the Cartan subgroup of $G=\mathrm{Ad} g$ corresponding to $\hbar$. Also let $D \subseteq \hbar^{\prime}$ be defined in (5.1.5) so that $D$ parametrizes the equivalence classes of the finite-dimensional irreducible holomorphic representations of $g c$ or $G_{c}^{s}$. We recall that $\pi_{\lambda}$, for $\lambda \in D$, is a representation corresponding to $\lambda$. See (5.1.7). Now let $\nu_{k} \in D, k=1, \ldots, l$, be defined as in Section 5.1 so that $D=\Sigma_{k} \mathbb{Z}_{+} \nu_{k}$. See (5.1.6). The representations $\pi_{\nu_{k}}$, $k=1, \ldots, l$, are called the fundamental representations of $g \mathbb{c}$ or $G_{\mathbb{C}}{ }^{s}$. We will write $\pi_{k}$ for $\pi_{v_{k}}$, and in much of the earlier notation we will replace $\lambda$ by $k$ when $\lambda=\nu_{k} \in D$.
Let $g_{o} \in H$ and $w_{o} \in h_{+}$. For any $\lambda \in D$ and $t \in \mathbb{R}$ we defined the function $\Phi_{\lambda}\left(g_{o}, w_{0} ; t\right)$ in Section 5.10. See (5.10.6). We write $\Phi_{k}\left(g_{0}, w_{o} ; t\right)$ for this function when $\lambda=\nu_{k}$. If $\Delta^{k} \subseteq \hbar^{\prime}$ is the set of weights of the fundamental representation $\pi_{k}$ then by ( 5.10 .6 )

$$
\begin{equation*}
\Phi_{k}\left(g_{o}, w_{0} ; t\right)=(-1)^{\left.o v_{k}-k v_{k}\right)} \sum_{\nu \in \in^{t}} r_{\nu}\left(w_{o}\right) g_{o}^{\nu} e^{-t\left(v, w_{0}\right)}, \tag{7.5.1}
\end{equation*}
$$

where $r_{\nu}$ is the rational function of $w_{o}$ given by

$$
\begin{equation*}
r_{v}\left(w_{o}\right)=\sum_{s \in \in \mathscr{S}^{k}(v)} \frac{c_{s, k}}{p_{o\left(v_{k}-v\right)}\left(s, w_{o}\right)}, \tag{7.5.2}
\end{equation*}
$$

where $\mathscr{S}^{k}(\nu)=\mathscr{S}^{\lambda}(\nu)$ for $\lambda=\nu_{k}$ and $\mathscr{S}^{\lambda}(\nu)$ is defined in (5.10.1). Also, $c_{s, k}=c_{s, \nu_{k}}$, where the latter is defined by (5.9.3). Furthermore $p_{j}\left(s, w_{0}\right)$ is defined in (5.10.2). See also (5.10.3). The expression $o(\mu)$ is defined in (5.5.28). For later use we put $\mathscr{S}^{k}=\mathscr{S}^{\lambda}$ when $\lambda=\nu_{k}$ and $\mathscr{S}^{\lambda}$ is defined in Section 5.9.

Now as in (7.3.1) let $C$ be the Cartan matrix defined by $C_{i j}=2 Q\left(\alpha_{i}, \alpha_{j}\right) /$ $Q\left(\alpha_{j}, \alpha_{j}\right)$. If $R$ is the $l \times l$ diagonal matrix defined by putting $R_{j j}=2 / Q\left(\alpha_{j}, \alpha_{j}\right)$ one has

$$
\begin{equation*}
C=B R \tag{7.5.3}
\end{equation*}
$$

recalling (7.1.19). But also recalling (see Section 5.1) the definition of the highest weights $\nu_{k}$ of the fundamental representations one notes that the matrix expressing the simple roots $\alpha_{j}$ in terms of the $\nu_{k}$ is just the Cartan matrix. That is,

$$
\begin{equation*}
\alpha_{j}=\sum C_{j k} \nu_{k} \tag{7.5.4}
\end{equation*}
$$

since $2 Q\left(\nu_{i}, \alpha_{j}\right) / Q\left(\alpha_{j}, \alpha_{j}\right)-\delta_{i j}$.
Remark 7.5.1. With regard to formula (7.5.2) one first of all has $\mathscr{P}^{k}(\nu) \subseteq \mathscr{S}^{k}$. The Cartan matrix $C$ (or rather $C^{-1}$ ) which plays a critical role in describing the geometry of the $\psi_{i}$ in our Hamiltonian $H$ also may be used to describe the set $\mathscr{S}^{k}$ of sequences

$$
\begin{equation*}
s=\left(i_{1}, i_{2}, \ldots, i_{o\left(v_{k}-\kappa v_{k}\right.}\right), \tag{7.5.5}
\end{equation*}
$$

where $1 \leqslant i_{j} \leqslant l$ and, if $d_{k}=o\left(\nu_{k}-\kappa \nu_{k}\right), \sum_{j=1}^{d_{k}} \alpha_{i j}=v_{k}-\kappa \nu_{k}$. In fact inverting (7.5.4) one has $\nu_{k}=\Sigma\left(C^{-1}\right)_{k j} \alpha_{j}$. But $-\kappa \alpha_{j}$ is again a simple positive root and if $1 \leqslant j \leqslant l$ is defined by

$$
\begin{equation*}
\alpha_{J}=-\kappa \alpha_{j} \tag{7.5.6}
\end{equation*}
$$

one has

$$
\begin{equation*}
\nu_{k}-\kappa \nu_{k}=\sum_{j=1}^{l}\left(\left(C^{-1}\right)_{k j}+\left(C^{-1}\right)_{k J}\right) \alpha_{j} . \tag{7.5.7}
\end{equation*}
$$

Thus $\left(C^{-1}\right)_{k j}+\left(C^{-1}\right)_{k j}$ is a nonnegative integer and

$$
\begin{equation*}
o\left(\nu_{k}-\kappa \nu_{k}\right)=\sum_{j=1}^{l}\left(C^{-1}\right)_{k j}+\left(C^{-1}\right)_{k j} \tag{7.5.8}
\end{equation*}
$$

Also,

$$
\begin{equation*}
\operatorname{card} \mathscr{S}^{k}=\frac{o\left(\nu_{k}-\kappa \nu_{k}\right)!}{\prod_{j=1}^{l}\left(\left(C^{-1}\right)_{k j}+C_{k j}^{-1}\right)!} . \tag{7.5.9}
\end{equation*}
$$

In fact $\mathscr{S}^{k}$ is the set of all sequences $s$ such that any $1 \leqslant j \leqslant l$ occurs $\left(C^{-1}\right)_{k j}+$ $\left(C^{-1}\right)_{k j}$ times in $s$.

The following theorem is one of the main results of the paper. It determines the trajectories of the Hamiltonian vector field $\xi_{H}$. In the theorem which follows the next theorem we give the expression for the asymptotic, or scattering, behavior of the system. First, however, we introduce certain constants. Recalling Proposition 7.4 .3 let $E$ be the $n \times n$ projection matrix defined by putting $E=A^{*} D$. It reduces to the identity if $l=n$. By Proposition 7.4.3, $E$ may be given, in general, by

$$
\begin{equation*}
\Gamma_{1} p_{i}=\sum_{j=1}^{n} E_{i j} p_{j}, \quad i=1, \ldots, n . \tag{7.5.10}
\end{equation*}
$$

For any $x \in \mathbb{R}^{2 n}$ and $i=1, \ldots, n$ let

$$
\begin{equation*}
\bar{q}_{i}(x)=q_{i}(x)-\frac{1}{m_{i}} \sum_{j=1}^{n} E_{i j} m_{j} q_{j}(x) \tag{7.5.11}
\end{equation*}
$$

and

$$
\begin{equation*}
\bar{p}_{i}(x)=p_{i}(x)-\sum_{j=1}^{n} E_{i j} p_{j}(x) . \tag{7.5.12}
\end{equation*}
$$

Also, where $r_{j}>0, j=1, \ldots, l$, are the coefficients occurring in (7.1.11) let

$$
\begin{equation*}
b_{i}=-\sum_{j=1}^{l}\left(A^{*} B^{-1}\right)_{i j} \log r_{j} \tag{7.5.13}
\end{equation*}
$$

using the notation of Proposition 7.4.1, for $i=1, \ldots, n$.
Remark 7.5.2. One notes that $\bar{q}_{i}(x)$ and $\bar{p}_{i}(x)$ vanish if $l=n$. Also,

$$
\begin{equation*}
b_{i}=0 \quad \text { if all } \quad r_{k}=1 \tag{7.5.14}
\end{equation*}
$$

as in the usual Toda lattice and in the examples of Section 7.4.
Theorem 7.5. Let $\left(\mathbb{R}^{2 n}, \omega\right)$ be a classical phase space with canonical coordinates $p_{i}, q_{j} \in C^{\infty}\left(\mathbb{R}^{2 n}\right), i, j=1, \ldots, n$. Let $H \in C^{\infty}\left(\mathbb{R}^{2 n}\right)$ be a Hamiltonian function of the form

$$
\begin{equation*}
H=\sum_{j=1}^{n} \frac{p_{j}^{2}}{2 m_{j}}+r_{1} e^{t_{1}}+\cdots+r_{l} e^{\psi_{l}} \tag{7.5.15}
\end{equation*}
$$

where (the masses) $m_{j}>0, r_{i}>0$, and $l \leqslant n$,

$$
\begin{equation*}
\psi_{i}=\sum_{j=1}^{n} a_{i j} q_{j}, \quad i=1, \ldots, l \tag{7.5.16}
\end{equation*}
$$

are linearly independent linear combinations of the $q_{j}$ 's which define a Dynkin diagram with respect to the bilinear form $B_{H}$ given by (7.1.12). See (7.1.13) and (7.1.14). Then there exist, uniquely up to isomorphism, a real split semi-simple Lie algebra $g$ of rank $l$, a split Cartan subalgebra $\ell$, simple positive roots $\alpha_{1}, \ldots, \alpha_{l}$ in the dual $\hbar^{\prime}$, a nonsingular invariant bilinear form $Q$ on $g$ and correspondingly on $\ell^{\prime \prime}$ such that

$$
\begin{equation*}
B_{H}\left(\psi_{i}, \psi_{j}\right)=Q\left(\alpha_{i}, \alpha_{j}\right) . \tag{7.5.17}
\end{equation*}
$$

Let $\hbar_{+} \subseteq \hbar_{i}$ be the open Weyl chamber corresponding to the $\alpha_{i}$. Let $G$ be the adjoint group of $g$ and let $H \subseteq G$ be the subgroup corresponding to $f$.

Now let $\xi_{H}$ be the Hamiltonian vector field on $\mathbb{R}^{2 n}$ corresponding to the Hamiltonian function $H$. See (7.1.2). Then $\xi_{H}$ is globally integrable on $\mathbb{R}^{2 n}$. Let $x \in \mathbb{R}^{2 n}$ be arbitrary and let $x(t)=\exp \xi_{H} \cdot x$ be the trajectory of $\xi_{H}$, where $x(0)=x$. In particular then $q_{i}(x(t))$ is the position of the ith particle at time $t$ when the initial state of the system is $\left(q_{1}(x), \ldots, q_{n}(x), p_{1}(x), \ldots, p_{n}(x)\right)$. Let $y \in g$ be the normalized Jacobi element defined by (7.4.13)-(7.4.15) and let $\left(g_{o}, w_{o}\right) \in H \times \hbar_{+}$ be the "parameters" corresponding to $y$ in the sense of Section 3.7. That is, $w_{o} \in \ell_{+}$ (the "diagonal" representation of $y$ ) is defined by the relation $\mathscr{I}\left(w_{o}\right)=\mathscr{I}(y)$ (see (2.3.1)) and $g_{o} \in H$ is defined so that $\rho(g)=g_{o}$, where $g \in G_{0}{ }^{w}$ satisfies $\beta_{(w)}(g)=y$, using the notation of Theorem 3.6. Then if $\Phi_{k}\left(g_{o}, w_{o} ; t\right)$ is the finite sum of exponentials defined by (7.5.1) (and hence given in terms of the fundamental representation of $g_{\mathbb{C}}$ ) one has for $i=1, \ldots, n$

$$
\begin{equation*}
q_{i}(x(t))=\bar{q}_{i}(x)+\frac{1}{m_{i}}\left(b_{i}+\bar{p}_{i}(x) t-2 \sum_{k=1}^{l} \frac{a_{k i}}{Q\left(\alpha_{k}, \alpha_{k}\right)} \log \Phi_{k}\left(g_{o}, w_{o} ; t\right)\right), \tag{7.5.18}
\end{equation*}
$$

where $m_{i}$ is given in (7.5.15), $a_{k i}$ is given in (7.5.16) and $\bar{q}_{i}(x), \bar{p}_{i}(x)$, and $b_{i}$ are given respectively by (7.5.11), (7.5.12), and (7.5.13).

Proof. Recalling the notation and statement of Theorem 6.8.1 the Hamiltonian vector field $\xi$ on $Z$ corresponding to $I_{1} \mid Z$ is globally integrable. Furthermore if, for notational convenience, $\varphi_{k}(t)=\varphi_{k}(\exp t \xi \cdot y)$ then, by (6.8.4),

$$
\begin{equation*}
\varphi_{k}(t)=\log h(g \exp (-t) w)^{-\alpha_{k}} \tag{7.5.19}
\end{equation*}
$$

where, as usual, $w=f+w_{o}$. But then by Proposition 7.4.2 the Hamiltonian vector field $\xi_{r_{1}^{\prime}}^{\prime}$ on $Z^{\prime}$ is globally integrable and if one puts $\varphi_{k}^{\prime}(t)=\varphi_{k}^{\prime}\left(\exp t \xi_{r_{1}^{\prime}}^{\prime} y^{\prime}\right)$, where $\delta x=y^{\prime}$, one has

$$
\begin{equation*}
\varphi_{k}^{\prime}(t)=\varphi_{k}(t), \tag{7.5.20}
\end{equation*}
$$

where we choose $y=\sigma y^{\prime}$. One notes also that, by Proposition 7.4.4, $y$ is given by (7.4.13)-(7.4.15). But now $\xi_{H}$ is globally integrable by Proposition 7.2
and we may apply (7.2.5) to determine $q_{i}(t)$, which we have, more properly, written here as $q_{i}(x(t))$. But now by (7.1.15), $\varphi_{k}^{\prime}(0)=q_{k}^{\prime}(x)=\psi_{k}(x)+\log r_{k}$, recalling the definition (see (7.2.3)) of $\delta$. But $\psi_{k}=\sum_{j=1}^{n} a_{k j} q_{j}$. Thus for $s=1, \ldots, l$

$$
\begin{align*}
\sum_{k=1}^{l}\left(B^{-1}\right)_{s k} \varphi_{k}^{\prime}(0) & =\sum_{j=1}^{n}\left(B^{-1} A\right)_{s j} q_{j}(x)+\sum_{k=1}^{l}\left(B^{-1}\right)_{s k} \log r_{k}  \tag{7.5.21}\\
& =\sum_{j=1}^{n} D_{s j} m_{j} q_{j}(x)+\sum_{k=1}^{l}\left(B^{-1}\right)_{s k} \log r_{k}
\end{align*}
$$

since one has $D=B^{-1} A M$ and hence $B^{1} A=D M^{-1}$. See (7.4.9). But then for $i=1, \ldots, n$, since $E=A^{*} D$, one has, by (7.5.13) upon applying $A^{*}$ to (7.5.21),

$$
\sum_{k=1}^{l}\left(A^{*} B^{-1}\right)_{i k} \varphi_{k}^{\prime}(0)=\sum_{j=1}^{n} E_{i j} m_{j} q_{i}(x)-b_{i}
$$

But $q_{i}(t)$ equals $q_{i}(x)$ for $t=0$ and hence by (7.2.5) and (7.5.11) one has

$$
\begin{equation*}
q_{i}(t)=\bar{q}_{i}(x)+\frac{1}{m_{i}}\left(c_{i} t+b_{i}+\sum_{k=1}^{l}\left(A^{*} B^{-1}\right)_{i k} \varphi_{k}^{\prime}(t)\right), \tag{7.5.22}
\end{equation*}
$$

where $c_{i}=\left(\Gamma_{2} p_{i}\right)(x)$. But then by (7.5.10) and (7.5.12)

$$
\begin{equation*}
c_{i}=\bar{p}_{i}(x) . \tag{7.5.23}
\end{equation*}
$$

Now $\varphi_{j}^{\prime}(t)=\log h(g \exp (-t) w)^{-\alpha_{j}}$ for $j=1, \ldots, l$, by (7.5.19) and (7.5.20) and hence for $k=1, \ldots, l$

$$
\begin{equation*}
\sum_{j=1}^{l}\left(B^{-1}\right)_{k j} \varphi_{j}^{\prime}(t)=\sum_{j=1}^{l}\left(B^{-1}\right)_{k j} \log h(g \exp (-t) w)^{-\alpha_{j}} \tag{7.5.24}
\end{equation*}
$$

But $B^{-1}=R C^{-1}$ by (7.5.3) and hence $\sum_{j=1}^{l}\left(B^{-1}\right)_{k^{j}} \alpha_{j}=2 \nu_{k} / Q\left(\alpha_{k}, \alpha_{k}\right)$ by (7.5.4). But then if we put $\left(B^{-1}\right)_{j k}$ in the exponential on the right side of (7.5.24) one has

$$
\begin{aligned}
\sum_{j=1}^{l}\left(B^{-1}\right)_{k j} \varphi_{j}^{\prime}(t) & =-\frac{2}{Q\left(\alpha_{k}, \alpha_{k}\right)} \log h(g \exp (-t) w)^{\nu_{k}} \\
& =-\frac{2}{Q\left(\alpha_{k}, \alpha_{k}\right)} \log \Phi_{k}\left(g_{n}, w_{n} ; t\right)
\end{aligned}
$$

by (5.10.7) for $\lambda=\nu_{k}$. Hence applying $A^{*}$ one has for $i=1, \ldots, n$

$$
\begin{equation*}
\sum_{j=1}^{l}\left(A^{*} B^{-1}\right)_{i j} \varphi_{j}^{\prime}(t)=-2 \sum_{k=1}^{l} \frac{a_{k i}}{Q\left(\alpha_{k}, \alpha_{k}\right)} \log \Phi_{k}\left(g_{o}, w_{o} ; t\right) . \tag{7.5.25}
\end{equation*}
$$

The result (7.5.18) then follows from (7.5.22) and (7.5.23), recalling that $q_{i}(t)=q_{i}(x(t))$.
Q.E.D.

Remark 7.5.3. One notes that knowing the $q_{i}(x(t))$ determines the entire trajectory $x(t)$ in $\mathbb{R}^{2 n}$. This is clear from Hamilton's equations, (7.1.5), which for the case at hand asserts that

$$
p_{i}(x(t))=m_{i}\left(\frac{d q_{i}}{d t}\right)(t) .
$$

7.6. We now determine the asymptotic behavior of $q_{i}(x(t))$ as $t \rightarrow+\infty$ and $-\infty$. Since $q_{i}(x(t))$, by (7.5.18), is linear in $t$ plus the $\log$ of a finite sum of exponentials it follows easily that $q_{i}(x(t))$ is asymptotic to straight lines as $t \rightarrow \pm \infty$. We wish to determine both of these straight lines.

For any $i=1, \ldots, n$ let $\mu_{i}^{-} \in h^{\prime}$ be the linear form on the Cartan subalgebra $h$ defined by

$$
\begin{equation*}
\mu_{i}^{-}=2 \sum_{k=1}^{b} \frac{a_{k i} \nu_{k}}{Q\left(\alpha_{k}, \alpha_{k}\right)} \tag{7.6.1}
\end{equation*}
$$

and put $\mu_{i}{ }^{+}=\kappa \mu_{i}{ }^{-}$, where $\kappa$, as usual, is the element of the Weyl group which interchanges the positive and negative roots. Now let $x \in \mathbb{R}^{2 n}$ and let $w_{o} \in h_{+}$ be as in Theorem 7.5. Put

$$
\begin{equation*}
v_{i}^{+}(x)=\frac{1}{m_{i}}\left(\bar{p}_{i}(x)+\left\langle\mu_{i}^{+}, z v_{0}\right\rangle\right) \tag{7.6.2}
\end{equation*}
$$

and

$$
\begin{equation*}
v_{i}^{-}-(x)=\frac{1}{m_{i}}\left(\bar{p}_{i}(x)+\left\langle\mu_{i}^{-}, w_{a_{a}}\right\rangle\right) . \tag{7.6.3}
\end{equation*}
$$

Also let $d(w)$ be the element of the Cartan subgroup $H$ defined in (5.5.2). It depends only on the isospectral leaf $Z(\gamma) \subseteq Z$ containing $y=\sigma \delta x$. As usual $w=w_{0}+f$. Also, $\gamma=\mathscr{I}(y)$. For $i=1, \ldots, n$ and where $g_{o} \in H$ is defined as in Theorem 7.5 put

$$
\begin{equation*}
u_{i}^{+}(x)=\bar{q}_{i}(x)+\frac{1}{m_{i}}\left(b_{i}+\log \left(\left(g_{o}^{-1} d(w)\right)^{u_{i}^{+}}\right)\right) \tag{7.6.4}
\end{equation*}
$$

and

$$
\begin{equation*}
u_{i}-(x)=\bar{q}_{i}(x)+\frac{1}{m_{i}}\left(b_{i}+\log \left(\left(g_{o}^{-1} d(w)^{-1}\right)^{u_{i}^{-}}\right)\right) . \tag{7.6.5}
\end{equation*}
$$

Remark 7.6. Recall that $g_{o}^{\nu_{k}}$ and $d(w)^{\nu_{k}}$ are given in more explicit terms by (5.5.25) and (5.5.4). With regard to (5.5.25) one uses (3.7.4) to determine $h(g)^{\lambda}$. On the other hand note that, by (7.6.1),

$$
\begin{equation*}
\log \left(\left(g_{o}^{-1} d(w)^{-1}\right)^{u_{i}}\right)=-2 \sum_{k=1}^{l} \frac{a_{k i}}{Q\left(\alpha_{k}, \alpha_{k}\right)}\left(\log \left(g_{o}^{\nu_{k}}\right)+\log \left(d(w)^{\nu_{k}}\right)\right) . \tag{7.6.6}
\end{equation*}
$$

On the other hand if $1 \leqslant k \leqslant l$ is defined by (7.5.6) then

$$
\begin{equation*}
\log \left(\left(g_{o}^{-1} d(w)\right)^{\mu_{i}^{+}}\right)=-2 \sum_{k=1}^{l} \frac{a_{k i}}{Q\left(\alpha_{k}, \alpha_{k}\right)}\left(-\log \left(g_{o}^{\nu_{k}}\right)+\log \left(d(w)^{\nu_{k}}\right)\right) \tag{7.6.7}
\end{equation*}
$$

since $Q\left(\alpha_{k}, \alpha_{k}\right)=Q\left(\alpha_{k}, \alpha_{k}\right)$ and hence $-\kappa \nu_{k}=\nu_{k}$.
The point is of course that the two asymptotic lines are $v_{i}{ }^{+}(x) t+u_{i}{ }^{+}(x)$ and $v_{i}-(x) t+u_{i}-(x)$. They are thus given by the highest $\nu_{k}$, and lowest weights, $\kappa \nu_{k}$, of the fundamental representations of $g \mathbb{C}$.

Theorem 7.6. Let the notation be as in Theorem 7.5 and as in (7.6.2), (7.6.3) and (7.6.4), (7.6.5). Then for $i=1, \ldots, n$

$$
\begin{equation*}
\lim _{t \geqslant 1 \infty}\left(\left(v_{i}^{+}(x) t+u_{i}^{+}(x)\right)-q_{i}(x(t))=0\right. \tag{7.6.8}
\end{equation*}
$$

and

$$
\begin{equation*}
\lim _{t \rightarrow-\infty}\left(\left(v_{i}^{-}(x) t+u_{i}^{-}(x)\right)-q_{i}(x(t))=0\right. \tag{7.6.9}
\end{equation*}
$$

Proof. By (7.6.1) and (5.10.7) one has

$$
\begin{equation*}
\log \left(h(g \exp (-t) w)^{u_{i}-}\right)=2 \sum_{k=1}^{l} \frac{a_{k i}}{Q\left(\alpha_{k}, \alpha_{k}\right)} \log \Phi_{k}\left(g_{o}, w_{o} ; t\right) \tag{7.6.10}
\end{equation*}
$$

so that by (7.5.18)

$$
\begin{equation*}
q_{i}(x(t))=\bar{q}_{i}(x)+\frac{1}{m_{i}}\left(b_{i}+\bar{p}_{i}(x) t-\log \left(h(g \exp (-t) w)^{\mu_{i}}\right)\right) . \tag{7.6.11}
\end{equation*}
$$

But then if $\mu=\mu_{i}^{-}$in Theorem 5.11 one easily has

$$
\begin{equation*}
q_{i}(x(t))-\left(v_{i}^{+}(x) t+u_{i}^{+}(x)\right)=-R_{+}\left(g_{0}, w_{0} ; t\right) \tag{7.6.12}
\end{equation*}
$$

and

$$
\begin{equation*}
q_{i}(x(t))-\left(v_{i}^{-}(x) t+u_{i}^{-}(x)\right)=-R_{-}\left(g_{0}, w_{0} ; t\right) . \tag{7.6.13}
\end{equation*}
$$

However, $\lim _{t \rightarrow+\infty} R_{+}\left(g_{o}, w_{o} ; t\right)==\lim _{t \rightarrow-\infty} R_{-}\left(g_{o}, w_{o} ; t\right)=0$ by Theorem 5.11. This proves (7.6.8) and (7.6.9).
Q.E.D.

We refer respectively to $v_{i}{ }^{+}(x)$ and $v_{i}-(x)$ as the $+\infty$ and $-\infty$ limiting velocities of the $i$ th particle when it occupies the classical state $x \in \mathbb{R}^{2 n}$ at time $t=0$. Also (with J. Moser) under the same conditions we refer to $u_{i}^{+}(x)$ and $u_{i}^{-}(x)$, respectively, as the $+\infty$ and $-\infty$ limiting phases of the $i$ th particle. Graphically plotting the position $q_{i}$ of the $i$ th particle against time one has


Of the two "parameters" $\left(g_{o}, w_{o}\right)$ of $y=\sigma \delta x \in Z$ (see Section 3.7) we recall that $w_{o} \in h_{+}$picks out the isospectral leaf $Z(\gamma), \gamma=\mathscr{I}(y)$, containing $y$ and $g_{o}$ (with its action angle coordinates $\log g_{o}{ }^{\nu}$ ) picks out $y$ in that leaf. See Theorems 3.7 and 4.3.

Proposition 7.6.1. Let the notation be as in Theorem 7.6.1. Then for any $x \in \mathbb{R}^{2 n}$ the difference $v_{i}{ }^{+}(x)-v_{i}-(x)$ of the limiting velocities of the ith particle depends only on the "spectrum" of the Jacobi element $y=\sigma \delta x \in g$. In fact one has

$$
\begin{equation*}
v_{i}^{+}(x)-v_{i}^{-}(x)=\left\langle\kappa \mu_{i}^{-}-\mu_{i}^{-}, w_{o}\right\rangle . \tag{7.6.14}
\end{equation*}
$$

The phase change, on the other hand, depends upon both $g_{o}$ and $w_{o}$. In fact recalling that d(w) depends only on $w_{0}$ one has

$$
\begin{equation*}
u_{i}^{+}(x)-u_{i}^{-}(x)=\log g_{o^{\mu^{-}}-\kappa \mu_{i}^{-}}+\log d(w)^{\mu_{i}^{-}+\kappa \mu_{i}^{-}} . \tag{7.6.15}
\end{equation*}
$$

Proof. This is immediate from definitions (7.6.2), (7.6.3) and (7.6.4), (7.6.5).

Proposition 7.6.2. Let the notation be as in Theorem 7.6.1. Assume that every one of the simple components of $g$ is of type $A_{1}, B_{l}, C_{l}, D_{2 l}, G_{2}, F_{4}$,
$E_{7}$ or $E_{8}$. Then the phase change depends only on the action angle coordinates. That is,

$$
\begin{equation*}
\mu_{i}^{+}(x)-\mu_{i}^{-}(x)=2 \log g_{o^{\mu^{-}}} . \tag{7.6.16}
\end{equation*}
$$

Moreover if in addition $l=n$ then one has a reversal of velocities. That is,

$$
\begin{equation*}
v_{i}^{+}(x)=-v_{i} \sim(x) . \tag{7.6.17}
\end{equation*}
$$

Proof. If $g$ satisfies the condition as stated then one knows that $\kappa$ equals minus the identity on $h^{\prime}$. (See, e.g., Theorem 0.16 in $\left[7\right.$, p. 335].) Thus $\mu_{i}-(x)=$ $-\mu_{i}{ }^{+}(x)$. But then (7.6.15) implies (7.6.16). On the other hand if $l=n$ then as noted in Remark 7.5.2 one has $\bar{p}_{i}(x)=0$. Thus (7.6.17) follows from definition (7.6.2), (7.6.3).
Q.E.D.
7.7. For the usual Toda lattice $l=n-1$ and $\psi_{i}=q_{i}-q_{i+1}$, $i=1, \ldots, n-1$. Also, $r_{i}=m_{i}=1, i=1, \ldots, n$. The corresponding Dynkin diagram is then of type $A_{l}$ so that we can take $g$ to be the space of all real $n \times n$ matrices of trace zero and $h$ to be the space of all diagonal matrices in $g$. If $w_{o} \in \hbar$ and

$$
\begin{equation*}
w_{o}=\operatorname{diag}\left(w_{1}, w_{2}, \ldots, w_{n}\right), \quad \text { where } \quad w_{i} \in \mathbb{R}, \quad i \geqslant 1, \tag{7.7.1}
\end{equation*}
$$

then we may choose the simple roots $\alpha_{i}, i==1, \ldots, l$, so that

$$
\begin{equation*}
\left\langle\alpha_{i}, w_{o}\right\rangle=w_{i}-w_{i+1} . \tag{7.7.2}
\end{equation*}
$$

One then has

$$
\begin{equation*}
w_{o} \in h_{+} \quad \text { if and only if } \quad w_{i}>w_{i+1} \tag{7.7.3}
\end{equation*}
$$

for all $i=1, \ldots, l$. Also, the Weyl group element $\kappa$ is given by

$$
\begin{equation*}
\kappa w_{o}=\operatorname{diag}\left(w_{n}, w_{n-1}, \ldots, w_{1}\right) . \tag{7.7.4}
\end{equation*}
$$

One notes that (recalling that the masses $m_{i}$ all equal 1) $B_{H}\left(\psi_{i}, \psi_{i}\right)=2$ so that

$$
\begin{equation*}
Q\left(\alpha_{i}, \alpha_{i}\right)=2 \tag{7.7.5}
\end{equation*}
$$

for the invariant bilinear form $Q$. It follows easily then that $Q$ is given by

$$
\begin{equation*}
Q(u, v)=\operatorname{tr} u v \tag{7.7.6}
\end{equation*}
$$

for $u, v \in g$, where the multiplication $u v$ is as operators on $\mathbb{R}^{n}$.

Now the action of $g$ on $\mathbb{R}^{n}$ extends to an action of $g$ as a Lie algebra of derivations of the exterior algebra $\Lambda \mathbb{R}^{n}$. The homogeneous subspace $\Lambda^{k} \mathbb{R}^{n}$ is stable under $g$ and as one knows, the fundamental representations $\pi_{k}, k=1, \ldots$, $n-1$, may be taken so that

$$
\begin{equation*}
\pi_{k}: g \rightarrow \operatorname{End} \Lambda^{k} \mathbb{R}^{n} \tag{7.7.7}
\end{equation*}
$$

and (7.7.7) is the restriction of this action of $g$ to $\Lambda^{k} \mathbb{R}^{n}$. The set $\Delta^{k}$ of weights of $\pi_{k}$ is in natural correspondence with the set of all $\binom{n}{k}$ subsets $1 \leqslant i_{1}<$ $i_{2}<\cdots<i_{k} \leqslant n$ of $k$ integers from 1 to $n$. In fact using such subsets as parameters, $\nu=\nu\left(i_{1}, \ldots, i_{k}\right) \in \Delta^{k}$, the correspondence is given by

$$
\begin{equation*}
\left\langle\nu, w_{o}\right\rangle=w_{i_{1}}+w_{i_{2}}+\cdots+w_{i_{k}} \tag{7.7.8}
\end{equation*}
$$

for $w_{o} \in \ell$. Now also write $r_{i_{1}, \ldots, i_{k}}\left(w_{o}\right)=r_{\nu}\left(w_{o}\right)$, using the notation of (7.5.2) so that $r_{i_{1} \ldots . . i_{k}}\left(w_{o}\right)$ is a rational function of the eigenvalues $w_{i}$ of $w_{o}$. Now $H$ can be taken to be the group of all diagonal positive matrices of determinant 1. Thus if $g_{o} \in H$ then we can write for $g_{i}>0, i=1, \ldots, n$,

$$
\begin{equation*}
g_{o}=\operatorname{diag}\left(g_{1}, \ldots, g_{n}\right) . \tag{7.7.9}
\end{equation*}
$$

We then note that for $g_{o} \in H, w_{o} \in h_{+}$the function $\Phi_{k}\left(g_{o}, w_{o} ; t\right)$ takes the form
$\Phi_{k}\left(g_{o}, w_{o} ; t\right)=c \sum_{1 \leqslant i_{1}<\cdots<i_{k} \leqslant n} r_{i_{1}, \ldots, i_{k}}\left(w_{o}\right) g_{i_{1}} \cdots g_{i_{k}} \exp \left(-t\left(w_{i_{1}}+\cdots+w_{i_{k}}\right)\right)$,
where $c=(-1)^{o\left(\nu_{k}-\kappa \nu_{k}\right)}$. See (7.5.1).
Remark 7.7.1. Recall that $\Phi_{k}\left(g_{o}, w_{o} ; t\right)=h(g \exp (-t) w)^{v_{k}}$. For computational purposes, by (7.7.7), one may then determine $\Phi_{k}\left(g_{0}, w_{o} ; t\right)$ by taking the $k \times k$ principal minor of the standard representation on $\mathbb{R}^{n}$ of the element $s_{o}(\kappa)^{-1} \bar{n}_{-f}(w) g_{o} \exp (-t) w_{o}\left(\bar{n}_{f}(w)\right)^{-1}=a(t)$, (see (5.3.6)), that is, the $k \times k$ principal minor of $\pi_{1}(a(t))$. The matrices $\pi_{1}\left(s_{0}(\kappa)^{-1}\right)$ and $\pi_{1}\left(g_{0} \exp (-t) w_{o}\right)$ are easy to write down. The matrices $\pi_{1}\left(\bar{n}_{-f}(w)\right)$ and $\pi_{1}\left(\bar{n}_{f}(w)^{-1}\right)$ are computable from (5.8.6) and (5.8.7). In fact the matrix $\pi_{1}\left(\bar{n}_{-f}(w)\right)$, using Proposition 5.8.2, is written down in Proposition 7.7.3.

If one does use (7.7.10) to determine $\Phi_{k}\left(g_{o}, w_{o} ; t\right)$ then first of all with regard to the constant $c$.

Proposition 7.7.1. One has

$$
\begin{equation*}
o\left(\nu_{k}-\kappa \nu_{k}\right)=k(n-k) \tag{7.7.11}
\end{equation*}
$$

so that if $n$ is odd then $c=1$ for all $k$ and if $n$ is even then $c$ alternates in sign with $k$.

Proof. It is clear from (7.7.2) and (7.7.6) that if $w_{0} \in \ell$ is arbitrary and is given by (7.7.1) then

$$
\begin{equation*}
\left\langle\nu_{k}, w_{o}\right\rangle=w_{1}+w_{2}+\cdots+w_{k} \tag{7.7.12}
\end{equation*}
$$

Furthermore by (7.7.4) one has

$$
\begin{equation*}
-\kappa \nu_{k}=\nu_{n-k} \tag{7.7.13}
\end{equation*}
$$

But since $w_{1}+\cdots+w_{n}=0$ one then easily has

$$
\begin{equation*}
\left\langle\nu_{k}-\kappa \nu_{k}, w_{o}\right\rangle=\left(w_{1}+\cdots+w_{k}\right)-\left(w_{n-k}+w_{n-k+1}+\cdots+w_{n}\right) . \tag{7.7.14}
\end{equation*}
$$

On the other hand if $x_{o} \in h$ is defined as in Section 2.1, then, by (7.7.2),

$$
\begin{equation*}
x_{0}=\operatorname{diag}\left(\frac{n-1}{2}, \frac{n-3}{2}, \ldots, \frac{1-n}{2}\right) . \tag{7.7.15}
\end{equation*}
$$

But $o\left(\nu_{k}-\kappa \nu_{k}\right)=\left\langle\nu_{k}-\kappa \nu_{k}, x_{o}\right\rangle$. But then by (7.7.14) and (7.7.15) one has $o\left(\nu_{k}-\kappa \nu_{k}\right)=(n-1)+(n-3)+\cdots+(n-2 k+1)=k(n-k) . \quad$ Q.E.D.

Now by (1.5.2) and (7.7.6) we may clearly choose root vectors $e_{\alpha_{i}}, e_{-\alpha_{i}} \in g$ so that in terms of the usual matrix units $e_{i j}$ one has, for $i=1, \ldots, n-1=l$,

$$
\begin{equation*}
e_{\alpha_{i}}=e_{i i+1}, \quad e_{-\alpha_{i}}=e_{i+1 i} \tag{7.7.16}
\end{equation*}
$$

It follows then from (1.5.1) (see Remark 1.5.1) that with regard to the Cartan decomposition of $g$ one has that $k$ is the Lie algebra of all $n \times n$ skew-symmetric matrices. The inner product (see Section 5.1) in $V^{v_{k}}=\Lambda^{k} \mathbb{R}^{n}$, using the notation of Section 5.1, may then be taken to be the standard inner product on $\Lambda^{k} \mathbb{R}^{n}$. In particular if $\epsilon_{i} \in \mathbb{R}^{n}, i=1, \ldots, n$, is the standard basis of $\mathbb{R}^{n}$ (i.e., the $j$ th coordinate of $\epsilon_{i}$ is $\delta_{i j}$ ), then the elements $\left\{\epsilon_{i_{1}} \wedge \cdots \wedge \epsilon_{i_{k}}\right\}$, where $1 \leqslant$ $i_{1}<\cdots<i_{k} \leqslant n$, are an orthonormal weight basis of $\Lambda^{k} \mathbb{R}^{n}$. In particular the highest weight vector, now written $v^{k}=v^{\nu k}$, can be taken to be given as

$$
\begin{equation*}
v^{k}=\epsilon_{1} \wedge \cdots \wedge \epsilon_{k} \tag{7.7.17}
\end{equation*}
$$

Recall that $\mathscr{S}$ for the present case is the set of all finite sequences

$$
\begin{equation*}
s=\left(i_{1}, \ldots, i_{d}\right) \tag{7.7.18}
\end{equation*}
$$

where $1 \leqslant i_{j} \leqslant l=n-1$.
Lemma 7.7. Let $1 \leqslant j \leqslant n$. Then for any $s \in \mathscr{S}$ one has $e_{-s} \epsilon_{j}=0$ unless
either $s$ is trivial or is of the form $s=(j, j+1, j+2, \ldots, i-1)$, where $j<i \leqslant n$, in which case $e_{-s} \epsilon_{j}=\epsilon_{i}$.

Proof. This is immediate from (7.7.16). Indeed by (7.7.16) $e_{-\alpha_{k}} \epsilon_{i}=0$ unless $k=i$ in which case $e_{-\alpha_{k}} \epsilon_{i}=\epsilon_{i+1}$.
Q.E.D.

Let $m_{k i} \in \mathbb{Z}_{+}$be defined so that $\nu_{k}-\kappa \nu_{k}=\sum_{i=\mathbf{1}}^{l} m_{k i} \alpha_{i}$.
Remark 7.7.2. Note that $\nu_{k}-\kappa \nu_{k}=\nu_{n-k}-\kappa \nu_{n-k}$ by (7.7.13) so that to determine the $m_{k i}$ one need consider only the case where $k \leqslant[n / 2]$. In this case by (7.7.14) and (7.7.2) one easily has that $m_{k i}=i$ for $i \leqslant k, m_{k i}=k$ for $k \leqslant i \leqslant n-k$ and $m_{k i}=n-i$ for $i \geqslant n-k$.

By definition $\mathscr{S}^{k}$ is the set of all $s \in \mathscr{S}$ so that if $s$ is given by (7.7.18) then, by (7.7.11), $d=k(n-k)$ and $i$ occurs $m_{k i}$ times among the $i_{j}$. The definition of $r_{i_{1} \ldots i_{k}}\left(w_{o}\right)$ in (7.7.10) depends upon a sum over a subset of $\mathscr{S}^{k}$. By (7.5.2) what is needed then, to use (7.7.10) to determine $\Phi_{k}\left(g_{o}, w_{0} ; t\right)$, is a knowledge of $c_{s, k}$ for $s \in \mathscr{S}^{k}$.

Proposition 7.7.2. One has $c_{s, k}=0$ or $(-1)^{k(n-k)}$ for any $s \in \mathscr{S}^{k}$. Furthermore for any $s \in \mathscr{S}$ and $1 \leqslant i_{1}<\cdots<i_{k} \leqslant n$ one has $e_{-8} \epsilon_{i_{1}} \wedge \cdots \wedge \epsilon_{i_{k}}=0$ or there exists $1 \leqslant j_{1}<\cdots<j_{k} \leqslant n$ so that

$$
\begin{equation*}
e_{-s} \epsilon_{i_{1}} \wedge \cdots \wedge \epsilon_{i_{k}}=e_{j_{1}} \wedge \cdots \wedge e_{j_{k}} \tag{7.7.19}
\end{equation*}
$$

Finally using the notation of (5.2.10) the lowest weight vector $v^{\kappa \nu \nu_{k}}$ is given by

$$
\begin{equation*}
v^{k \nu_{k}}=(-1)^{k(n-k)} \epsilon_{n-k} \wedge \epsilon_{n-k+1} \wedge \cdots \wedge \epsilon_{n} . \tag{7.7.20}
\end{equation*}
$$

Proof. For any $1 \leqslant i \leqslant l$ note that $e_{-\alpha_{i}} \epsilon_{i_{1}} \wedge \cdots \wedge \epsilon_{i_{k}}=0$ unless there exists a $j$ such that $i=i_{j}$ and $i_{j+1} \geqslant 2+i_{j}$ (putting $i_{k+1} \stackrel{ }{=}=n+1$ ), in which case

$$
\begin{equation*}
e_{-\alpha_{i}} \epsilon_{i_{i}} \wedge \cdots \wedge \epsilon_{i_{j}} \wedge \cdots \wedge \epsilon_{i_{k}}=\epsilon_{i_{1}} \wedge \cdots \wedge \epsilon_{i_{j}+1} \wedge \cdots \wedge \epsilon_{i_{k}} \tag{7.7.21}
\end{equation*}
$$

This proves the second statement and (7.7.19) in particular. Thus if $s \in \mathscr{S}^{k}$ one has either $e{ }_{s} \epsilon_{1} \wedge \cdots \wedge \epsilon_{k}=0$ or

$$
\begin{equation*}
e_{-s} \epsilon_{1} \wedge \cdots \wedge \epsilon_{k}=\epsilon_{n-k} \wedge \cdots \wedge \epsilon_{n} . \tag{7.7.22}
\end{equation*}
$$

Let $\mathscr{S}_{*}^{k}$ be the set of all $s \in \mathscr{S}^{k}$ satisfying (7.7.22). But now since $v^{k v_{k}}$ has unit length it follows that ${v^{\kappa \nu} v_{k}}=a \epsilon_{n-k} \wedge \cdots \wedge \epsilon_{n}$, where $a$ is either 1 or -1 . But then by (7.7.22) one has $c_{s, k}-0$ if $s \notin \mathscr{P}_{*}^{k}$ and $c_{s, k}=a$ if $s \in \mathscr{S}_{*}^{k}$. But then $a=(-1)^{k(n-k)}$ by (5.9.7) and (7.7.11).
Q.E.D.

Remark 7.7.3. The question as to whether $c_{s, k}=(-1)^{k(n-k)}$ or 0 , i.e.,
whether $s \in \mathscr{S}_{*}^{k}$ or not, is of course a readily resolved combinatorial question using (7.7.21).

For the standard Toda lattice Theorem 7.5 becomes

Theorem 7.7.1. Let the notation be as in Theorem 7.5, where $l=n-1$, $\psi_{i}=q_{i}-q_{i+1}$, and $r_{j}=m_{j}=1$ for $j=1, \ldots, n$ so that we can take $g$ to be the set of all real $n \times n$ matrices of trace zero, $h$ to be the set of diagonal matrices in $g, \alpha_{1}, \ldots, \alpha_{l}$ to be the set of simple roots given by (7.7.2), and $Q$ equal to the bilinear form on $g$ given by (7.7.6). Then for any $x$ in phase space $\mathbb{R}^{2 n}$ one has

$$
\begin{align*}
q_{i}(x(t))= & \frac{1}{n}\left(\sum_{j=1}^{n}\left(q_{j}(x)+t p_{j}(x)\right)\right) \\
& +\log \Phi_{i-1}\left(g_{o}, w_{o} ; t\right)-\log \Phi_{i}\left(g_{o}, w_{o} ; t\right) \tag{7.7.23}
\end{align*}
$$

where $\log \Phi_{j}\left(g_{o}, w_{0} ; t\right)=0$ for $j=0$ or $n$ and is otherwise given by taking the $\log$ of (7.7.10).

Proof. Applying Theorem 7.5 for the present case we note first of all that $b_{i}=0$ by (7.5.14). Next recalling (7.1.17) one easily has that $\mathscr{P}_{2}$ is the onedimensional space spanned by $p_{1}+\cdots+p_{n}$ and $\mathscr{P}_{1}$ is the $(n-1)$-dimensional space of all $p \in \mathscr{P}$ of the form $p=\sum_{i=1}^{n} c_{i} p_{i}$, where $\sum c_{i}=0$. It follows easily that $\Gamma_{2} p_{j}=(1 / n) \sum_{i=1}^{n} p_{i}$ for any $j$. But then by (7.5.11) and (7.5.12) one has $\bar{q}_{i}(x)=(1 / n) \sum_{j=1}^{n} q_{j}(x)$ and $\bar{p}_{i}(x)=(1 / n) \sum_{j=}^{n} p_{j}(x)$. Next recalling (7.1.10) one has $a_{i j}=0$ unless $j=i, i+1$, and

$$
\begin{equation*}
a_{i i}=1, \quad a_{i i+1}=-1 \tag{7.7.24}
\end{equation*}
$$

The result (7.7.23) then follows from (7.5.18) and (7.7.5).
Q.E.D.

We now wish to recover the results of Moser (see 4.3 in [19, p. 481]) on the scattering of the Toda lattice.

By (7.6.1) and (7.7.24) one has

$$
\begin{equation*}
\mu_{i}^{-}=v_{i}-v_{i-1} \tag{7.7.25}
\end{equation*}
$$

for $i=1, \ldots, n-1$, where we put $\nu_{0}=0$. But $-\kappa \nu_{i}=\nu_{n-i}$ by (7.7.13) and hence

$$
\begin{equation*}
\mu_{i}^{+}=\nu_{n+1-i}-v_{n-i} \tag{7.7.26}
\end{equation*}
$$

where $\nu_{n}=0$.

Proposition 7.7.3. Let the notation be as in Proposition 5.8.2 and above
where $w_{o}$ is given by (7.7.1). Let $m=\left\{m_{i j}\right\}$ be the $n \times n$ matrix $m=\pi_{1}\left(\bar{n}_{-f}(w)\right)$. Then $m_{i j}=0$ for $i<j, m_{i i}=1$ and for $i>j$ one has

$$
\begin{equation*}
m_{i j}=(-1)^{i+j}\left(\left(w_{j}-w_{j+1}\right)\left(w_{j}-w_{j+2}\right) \cdots\left(w_{j}-w_{i}\right)\right)^{-1} . \tag{7.7.27}
\end{equation*}
$$

Proof. As in Lemma 7.7 let $s \in \mathscr{S}$ be of the form $s=(j, j+1, \ldots, i-1)$, where $1 \leqslant j<i \leqslant n$. For such a sequence $s$ one clearly has

$$
\begin{equation*}
p\left(s, w_{o}\right)=\left(w_{j}-w_{j+1}\right)\left(w_{j}-w_{j+2}\right) \cdots\left(w_{j}-w_{i}\right), \tag{7.7.28}
\end{equation*}
$$

recalling (5.8.4) and (7.7.2). By Lemma 7.7 one has $e_{-s} \epsilon_{j}=\epsilon_{i}$ if $s$ is of this form and $e_{-s} \epsilon_{j}=0$ if $s$ is not of this form. The result then follows immediately from the computation of $\pi_{1}\left(\bar{n}_{-f}(w)\right) \epsilon_{j}$ using (5.8.6).
Q.E.D.

Now let $\delta_{k}\left(w_{o}\right)$ be the product of the determinant of the $k \times k$ minor in the lower left-hand corner of the matrix $m$ with $(-1)^{k(n-k)}$. That is, where $m$ is given by Proposition 7.7.3 put

$$
\begin{equation*}
\delta_{k}\left(w_{o}\right)=(-1)^{k(n-k)} \operatorname{det}_{i, j=1, \ldots, k} m_{n-k+i j} \tag{7.7.29}
\end{equation*}
$$

Remark 7.7.4. One may show inductively that

$$
\begin{equation*}
\delta_{k}\left(w_{o}\right)=\frac{\prod_{i<j \leqslant k} w_{i}-w_{j}}{\prod_{i<s \leqslant n, i \leqslant k} w_{i}-w_{s}} . \tag{7.7.30}
\end{equation*}
$$

The agreement of the following application of Theorem 7.7.1 with the results of Moser in [19] on the scattering of the Toda lattice will be clarified in Remark 7.7.5.

Theorem 7.7.2. Let the notation be as in Theorem 7.6 and as above. Thus for any point $x$ in phase space $v_{i}{ }^{+}(x) t+u_{i}{ }^{+}(x)$ and $v_{i}{ }^{-}(x) t+u_{i}{ }^{-}(x)$ respectively, are the asymptotic lines of the position curve $\left(t, q_{i}(x(t))\right.$, as $t \rightarrow+\infty$ and $t \rightarrow-\infty$, of the ith particle in the standard Toda lattice when the system occupies $x$ at $t=0$. Then if $w_{o}$ is defined as in Theorem 7.5 (so that $w_{o}$ picks out the isospectral leaf of the Jacobi matrix $y=\sigma \delta x$ ) and is described as a diagonal matrix by (7.7.1) one has, in terms of the eigenvalues of $w_{o}$,

$$
\begin{equation*}
v_{k}^{+}(x)=w_{n+1-k}+\bar{p}(x) \tag{7.7.31}
\end{equation*}
$$

and

$$
\begin{equation*}
v_{k}-(x)=w_{k}+\bar{p}(x), \tag{7.7.32}
\end{equation*}
$$

where $\bar{p}(x)=(1 / n) \sum_{i=1}^{n} p_{i}(x)$. In particular then

$$
\begin{equation*}
v_{k}^{+}(x)=\widetilde{v_{n+1-k}}(x) . \tag{7.7.33}
\end{equation*}
$$

Furthermore the phase difference $u_{n+1-i}^{+}-u_{i^{-}}^{-}$also depends only on $w_{o}$. In fact if $\delta_{k}\left(w_{0}\right)$ is defined by (7.7.29) one has

$$
\begin{equation*}
u_{n+1-k}^{+}(x)-u_{k}^{-}(x)=\log \delta_{k}\left(w_{o}\right)^{2}-\log \delta_{k-1}\left(w_{o}\right)^{2} . \tag{7.7.34}
\end{equation*}
$$

Proof. Recalling (7.7.12) onc has

$$
\left\langle v_{k}, w_{o}\right\rangle=w_{1}+w_{2}+\cdots+w_{k} .
$$

Thus $\left\langle\mu_{k}{ }^{+}, w_{o}\right\rangle=w_{n+1-k}$ and $\left\langle\mu_{k}^{-}, w_{o}\right\rangle=w_{k}$ by (7.7.25) and (7.7.26). But then (7.7.31) and (7.7.32) follow from (7.6.2) and (7.6.3).

But now computing $u_{n+1-k}^{+}(x)-u_{k}^{-}(x)$ using (7.6.4) and (7.6.5) note that $\bar{q}_{k}(x), b_{k} / m_{k}$, and also the dependence on $g_{o}$ cancel out. That is, by (7.6.4) and (7.6.5)

$$
\begin{equation*}
u_{n+1-k}^{+}(x)-u_{k}^{-}(x)=\log d(w)^{2\left(v_{k}-v_{k-1}\right)} \tag{7.7.35}
\end{equation*}
$$

But now by (5.5.4) one has $d(w)^{\nu_{k}}=\left\{\bar{n}_{-f}(w) v^{\nu_{k}}, v^{\kappa_{k} k}\right\}$. But then recalling (7.7.17) and (7.7.20), $d(w)^{\nu_{k}}$ is just $(-1)^{k(n-k)}$ times the determinant of the $k \times k$ minor in the lower left-hand corner of $\pi_{1}\left(\bar{n}_{-f}(w)\right)$. That is,

$$
\begin{equation*}
d(w)^{\nu_{k}}=\delta_{k}\left(w_{o}\right) . \tag{7.7.36}
\end{equation*}
$$

Then (7.7.34) follows from (7.7.35) and (7.7.36).
Q.E.D.

Remark 7.7.5. We now align the notation here with that in [19]. Let $L$ be the symmetric Jacobi matrix given in [19, p. 473]. Now let $D$ be the diagonal matrix

$$
\begin{equation*}
D=\operatorname{diag}\left(e^{x_{1} / 2}, \ldots, e^{x_{n} / 2}\right) \tag{7.7.37}
\end{equation*}
$$

Then

$$
\begin{equation*}
2\left(D L D^{-1}-\frac{1}{n}(\operatorname{tr} L) I\right)=y \tag{7.7.38}
\end{equation*}
$$

where $y$ is the normalized Jacobi element defined in Theorem 7.5 of this paper and $I$ is the $n \times n$ identity matrix. Since $2 b_{k}=-y_{k}$ by 2.1 in the notation of [19] ( $y_{k}$ in [19] is $p_{k}$ here) the initial state of the system $u$ has the negative of the momentum considered in this paper. Thus if $x_{k}(t)$ is defined as in [19] and $q_{k}(x(t))$ is defined as in this paper one has

$$
\begin{equation*}
x_{k}(-t)=q_{k}(x(t)) \tag{7.7.39}
\end{equation*}
$$

In particular

$$
\begin{equation*}
2 b_{k}=p_{k}(x) \tag{7.7.40}
\end{equation*}
$$

where in the notation of Theorem 7.7.2 here $p_{k}(x)$ is the momentum of the $k$ th particle in the initial state $x$. Now the eigenvalues of $y$ in decreasing order are given here, by (7.7.1) and (7.7.3), as $w_{1}, \ldots, w_{n}$. The eigenvalues of $L$ in increasing order are given as $\lambda_{1}, \ldots, \lambda_{n}$ in [19]. Thus by (7.7.38) one has

$$
\begin{equation*}
2 \lambda_{n-k+1}=\frac{2 \operatorname{tr} L}{n}+w_{k} . \tag{7.7.41}
\end{equation*}
$$

The asymptotic lines in [19] for $+\infty$ and $-\infty$ are written as $\alpha_{k}{ }^{+} t+\beta_{k}{ }^{+}$and $\alpha_{k}^{-} t+\beta_{k}^{-}$. However, by (7.7.39) one must have

$$
\begin{equation*}
\alpha_{k}^{+}=-v_{k}^{-(x)}, \quad \alpha_{k}^{-}=-v_{k}^{+}(x) . \tag{7.7.42}
\end{equation*}
$$

Also,

$$
\begin{equation*}
\beta_{k}^{+}=u_{k}^{-}(x), \quad \beta_{k}^{-}=u_{k}^{+}(x) \tag{7.7.43}
\end{equation*}
$$

But now (see Section 4 in [19]) one has $\alpha_{k}{ }^{+}=-2 \lambda_{n-k+1}$ and $\alpha_{k}^{-}=-2 \lambda_{k}$. Noting (7.7.41) this checks with our results (7.7.31) and (7.7.32) since, by (7.7.40), $\bar{p}(x)=(2 / n) \operatorname{tr} L$. Also, the result 4.2 in [19] is given here as (7.7.33). Now by (7.7.43)

$$
\begin{equation*}
\beta_{n-k+1}^{+}-\beta_{k}^{-}=u_{n-k+1}^{-}(x)-u_{k}^{+}(x) . \tag{7.7.44}
\end{equation*}
$$

By (7.7.34) the right side of (7.7.44) is given as $\log \delta_{n-k}\left(w_{o}\right)^{2}-\log \delta_{n-k+1}\left(w_{o}\right)^{2}$. However, noting Remark 7.7 .4 this yields

$$
\begin{equation*}
\beta_{n-k+1}^{+}-\beta_{k}^{-}=2 \log \frac{\prod_{j=1}^{k-1}\left(w_{n-k+1}-w_{n-j+1}\right)}{\prod_{i=k+1}^{n}\left(w_{n-i+1}-w_{n-k+1}\right)} . \tag{7.7.45}
\end{equation*}
$$

But by (7.7.31), (7.7.32), and (7.7.42)

$$
w_{n-k+1}-w_{n-r+1}=\alpha_{r}^{-}-\alpha_{k}^{-} .
$$

Thus

$$
\begin{equation*}
\beta_{n-k+1}^{+}-\beta_{k}^{-}=2 \log \frac{\prod_{j=1}^{k-1}\left(\alpha_{j}^{-}-\alpha_{k}^{-}\right)}{\prod_{i=k+1}^{n}\left(\alpha_{k}^{-}-\alpha_{i}^{-}\right)} . \tag{7.7.46}
\end{equation*}
$$

But this is exactly the result 4.3 in [19, p. 481].
7.8. As a further application of Theorem 7.5 consider the three-body
problem mentioned in (7.4.1). Here $n=3$ and $l=3$ and we recall the Hamiltonian $H$ given by

$$
\begin{equation*}
H=\sum_{i=1}^{3} \frac{p_{i}^{2}}{2}+e^{a_{1}-q_{2}}+e^{q_{2}-q_{3}}+e^{2 q_{3}} \tag{7.8.1}
\end{equation*}
$$

This then differs from the usual Toda lattice in that the particle whose position is given by $q_{3}$ is interacting also exponentially with some fixed mass. The Dynkin diagram in question, as one easily sees, is so that $g_{\mathrm{C}}$ is of type $C_{3}$ and hence $g$ is isomorphic to the Lie algebra of $\operatorname{Sp}(6, \mathbb{R})$.

We shall be dealing here with $6 \times 6$ matrices $u$ and we will often write

$$
u=\left(\begin{array}{ll}
A & B  \tag{7.8.2}\\
C & D
\end{array}\right)
$$

where $A=A(u), B=B(u), C=C(u)$, and $D=D(u)$ are $3 \times 3$ matrices. If $E$ is a $3 \times 3$ matrix let $E^{\prime}$ be the $3 \times 3$ matrix obtained from $E$ by transposing with respect to the diagonal which runs from the lower left-hand corner to the upper right-hand corner. Thus $\left(E^{\prime}\right)_{y}=E_{4-i 4-j}$. We find it particularly convenient to identify $g$ with the set of all matrices $u$ in (7.8.2), where

$$
\begin{equation*}
D=-A^{\prime}, \quad C=C^{\prime}, \quad \text { and } \quad B=B^{\prime} \tag{7.8.3}
\end{equation*}
$$

Remark 7.8.1. It is hoped that the reader is not mislead by the fact that the 6 in question describing the sizes of the matrices is also the dimension of the phase space of our system. If the $e^{2 q_{3}}$ term in (7.8.1) were replaced by $e^{q_{3}}$ then $g \cong S O(3,4)$ and we would be dealing with $7 \times 7$ matrices. The present case was chosen to simplify the computations.

The Cartan subalgebra $\hbar$ is the set of all $w_{o} \in g$ such that $C\left(w_{o}\right)=B\left(w_{o}\right)=0$ and

$$
\dot{A}\left(w_{0}\right)=\left(\begin{array}{lll}
w_{1} & 0 & 0  \tag{7.8.4}\\
0 & w_{2} & 0 \\
0 & 0 & w_{3}
\end{array}\right)
$$

The open Weyl chamber $\ell_{+} \subseteq \ell_{\text {may }}$ men be defined by the condition that $w_{0} \in h_{+}$if and only if $w_{1}>w_{2}>w_{3}>0$. One notes that the simple roots are given by
$\left\langle\alpha_{1}, w_{o}\right\rangle=w_{1}-w_{2}, \quad\left\langle\alpha_{2}, w_{o}\right\rangle=w_{2}-w_{3}, \quad\left\langle\alpha_{3}, w_{o}\right\rangle=2 w_{3}$.
The bilinear form $Q$ on $g$ is also easily seen to be given by $Q(y, z)=\frac{1}{2} \operatorname{tr} y z$. The condition that $k$ is the set of all skew-symmetric elements in $g$ then normalizes the root vectors up to sign. The choice of sign for the simple root
vectors will be made so that if $w_{0} \in \hbar_{+}$, where $A\left(w_{0}\right)$ is given by (7.8.4), then

$$
w=\left(\begin{array}{cccccc}
w_{1} & 0 & 0 & 0 & 0 & 0  \tag{7.8.6}\\
1 & w_{2} & 0 & 0 & 0 & 0 \\
0 & 1 & w_{3} & 0 & 0 & 0 \\
0 & 0 & \sqrt{2} & -w_{3} & 0 & 0 \\
0 & 0 & 0 & -1 & -w_{2} & 0 \\
0 & 0 & 0 & 0 & -1 & -w_{1}
\end{array}\right)
$$

recalling that $w=f+w_{o}$ and $f$ is given by (1.5.4).
Now the Lie group $G^{s}$ of Section 5.1 may be taken to be the subgroup of $G l(6, \mathbb{R})$ corresponding to $g$. The unipotent subgroups $\bar{N}$, and $N$ of $G^{s}$ are easily seen to be respectively groups of lower and upper triangular matrices and $H \subseteq G^{s}$ is a group of diagonal matrices. The element $\bar{n}_{f}(w) \in \bar{N}$ we recall is uniquely characterized by the conditions that $\bar{n}_{f}(w) w_{o} \bar{n}_{f}(w)^{-1}=w$. Using (5.8.7), as in the case of Proposition 7.7.3, one easily determines $\bar{n}_{f}(w)^{-1}$. In fact, explicitly

$$
\begin{align*}
A\left(\bar{n}_{f}(w)^{-1}\right) & =\left(\begin{array}{ccc}
1 & 0 & 0 \\
\frac{1}{w_{2}-w_{1}} & 1 & 0 \\
\frac{1}{\left(w_{1}-w_{3}\right)\left(w_{2}-w_{3}\right)} & \frac{1}{w_{3}-w_{2}} & 1
\end{array}\right),  \tag{7.8.7}\\
D\left(\bar{n}_{f}(w)^{-1}\right) & =\left(\begin{array}{ccc}
1 & 0 & 0 \\
\frac{1}{w_{2}-w_{3}} & 1 & 0 \\
\frac{1}{\left(w_{1}-w_{2}\right)\left(w_{1}-w_{3}\right)} & \frac{1}{w_{1}-w_{2}} & 1
\end{array}\right) . \tag{7.8.8}
\end{align*}
$$

Of course $B\left(\bar{n}_{f}(w)^{-1}\right)=0$. However, $(\sqrt{2} / 2) C\left(\bar{n}_{f}(w)^{-1}\right)$ equals

$$
\left[\begin{array}{ccc}
\frac{-1}{2 w_{3}\left(w_{2}+w_{3}\right)} & \frac{1}{2 w_{3}\left(w_{2}+w_{3}\right)} & \frac{-1}{2 w_{3}} \\
\times\left(w_{1}+w_{3}\right) & \frac{-1}{\left(w_{2}-w_{3}\right)\left(w_{2}+w_{3}\right)} & \frac{1}{\left(w_{2}-w_{3}\right)\left(w_{2}+w_{3}\right)} \\
\frac{1}{\left(w_{2}-w_{3}\right)\left(w_{2}+w_{3}\right)} & \times 2 w_{2} & \\
\times 2 w_{2}\left(w_{1}+w_{2}\right) & \frac{1}{-1} & \frac{-1}{\left(w_{1}-w_{2}\right)\left(w_{1}-w_{3}\right)} \\
\frac{-1}{\left(w_{1}-w_{2}\right)\left(w_{1}-v_{3}\right)} & \frac{1}{\left(w_{1}-w_{2}\right)\left(w_{1}-w_{3}\right)} & \frac{1}{\left(w_{1}\right)}
\end{array}\right) .
$$

Remark 7.8.2. One recognizes the product of roots occurring in the denominators. One knows that the most general root $\varphi$ is given by $\left\langle\varphi, w_{o}\right\rangle= \pm w_{i} \pm w_{j}$, $1 \leqslant i \leqslant j \leqslant 3$.

The element $\vec{n}_{-f}(w) \in \bar{N}$ is given by (5.8.6). One has $A\left(\bar{n}_{-f}(w)\right)_{i j}=$ $\left(D\left(\bar{n}_{f}(w)^{-1}\right)^{\prime}\right)_{i j}(-1)^{i+j}$. The same is true if $A$ and $D$ are interchanged. Furthermore

$$
\begin{equation*}
C\left(\bar{n}_{-f}(w)\right)=C\left(\bar{n}_{f}(w)^{-1}\right)^{\prime} . \tag{7.8.10}
\end{equation*}
$$

One can now be very explicit about the clement $d(w) \in H$ which we recall enters into the scattering of the mechanical system. An element $g_{o} \in H$ is determined by $A\left(g_{o}\right)$ since one easily has

$$
\begin{equation*}
D\left(g_{o}\right)=\left(A\left(g_{o}\right)^{-1}\right)^{\prime} \tag{7.8.11}
\end{equation*}
$$

## Proposition 7.8.1. One has

$A(d(w))$

$$
=\left[\begin{array}{ccc}
\frac{\sqrt{2}}{\left(w_{1}{ }^{2}-w_{2}{ }^{2}\right)\left(w_{1}{ }^{2}-w_{3}{ }^{2}\right) 2 w_{1}} & 0 & 0 \\
0 & \frac{\sqrt{2}\left(w_{1}{ }^{2}-w_{2}{ }^{2}\right)}{\left(w_{2}^{2}-v_{3}{ }^{2}\right)\left(w_{1}+w_{2}\right)^{2} 2 w_{2}} & 0 \\
0 & 0 & \frac{\sqrt{2}\left(w_{1}{ }^{2}-w_{3}{ }^{2}\right)\left(w_{2}{ }^{2}-w_{3}^{2}\right)}{\left(w_{1}+w_{3}\right)^{2}\left(w_{2}+w_{3}\right)^{2} 2 w_{3}}
\end{array}\right)
$$

Proof. One sees easily that the highest weight $\nu_{i}, i=1,2,3$, of the fundamental representation $\pi_{i}$ is given by

$$
\begin{equation*}
\left\langle v_{i}, w_{o}\right\rangle=\sum_{j=1}^{i} w_{j} \tag{7.8.13}
\end{equation*}
$$

It follows then that if we extend the action of $G^{8}$ on $\mathbb{R}^{6}$ to the exterior algebra $\Lambda \mathbb{R}^{6}$, such that $G^{s}$ operates as a group of automorphisms then we can take $V^{v_{i}}$ to be the subspace of $\Lambda^{i} \mathbb{R}^{6}$ spanned by $G^{s} \epsilon_{1} \wedge \cdots \wedge \epsilon_{i}$. We are using the notation of Section 7.7. Also, the inner product on $V^{\nu_{i}}$ is just the restriction to $V^{v_{i}}$ of the inner product on $\Lambda^{i} \mathbb{R}^{l}$ defined in Section 7.7. This is clear since $k$ is just the space of skew-symmetric elements in $g$. One then has

$$
\begin{equation*}
v^{\nu} i=\epsilon_{1} \wedge \cdots \wedge \epsilon_{i} \tag{7.8.14}
\end{equation*}
$$

and $\boldsymbol{v}^{\kappa \nu_{i}}= \pm \epsilon_{7-i} \wedge \cdots \wedge \epsilon_{6}$. But now by (5.5.4)

$$
d(w)^{v_{i}}=\left\{\bar{n}_{-f}(w) v^{\nu_{i}}, v^{k v_{i}}\right\} .
$$

Thus $d(w)^{\nu_{i}}$ is given by the absolute value (since it is necessarily positive)
of the determinant of the $i \times i$ minor in the lower left-hand corner of $C\left(\bar{n}_{f}(w)^{-1}\right)$. See (7.8.9). This fixes

$$
\begin{equation*}
v^{\kappa \nu_{i}}=-\epsilon_{7-i} \wedge \cdots \wedge \epsilon_{6} \tag{7.8.15}
\end{equation*}
$$

and one has

$$
\begin{align*}
& d(w)^{\nu_{1}}=\frac{\sqrt{2}}{\left(w_{1}^{2}-w_{2}^{2}\right)\left(w_{1}^{2}-w_{3}^{2}\right) 2 w_{1}}, \\
& d(w)^{\nu_{2}}=\frac{1}{\left(w_{1}^{2}-w_{3}^{2}\right)\left(w_{2}^{2}-w_{3}^{2}\right) w_{1} 2 w_{2}},  \tag{7.8.16}\\
& d(w)^{\nu_{3}}=\frac{\sqrt{2}}{\left(w_{1}+w_{2}\right)^{2}\left(w_{1}+w_{3}\right)^{2}\left(w_{2}+w_{3}\right)^{2} 2 w_{2} 2 w_{3}} .
\end{align*}
$$

But then by (7.8.14) the diagonal entries in $A(d(w))$ are just the relative quotients $d(w)^{\nu_{1}}, d(w)^{\nu_{2}-\nu_{1}}$, and $d(w)^{\nu_{3}-\nu_{2}}$. But then (7.8.12) follows from (7.8.16).
Q.E.D.

Now let $g_{o} \in H$ and write

$$
\begin{equation*}
A\left(g_{o}\right)=\operatorname{diag}\left(g_{1}, g_{2}, g_{3}\right) \tag{7.8.17}
\end{equation*}
$$

Thus

$$
A\left(g_{o} \exp t w_{o}\right)=\left(\begin{array}{ccc}
g_{1} e^{t w_{1}} & 0 & 0  \tag{7.8.18}\\
0 & g_{2} e^{t w_{2}} & 0 \\
0 & 0 & g_{3} e^{t w_{3}}
\end{array}\right) .
$$

Now by (5.2.10), (7.8.14), and (7.8.15) the element $s_{o}(\kappa)^{-1} \in G$ is given by $A\left(s_{o}(\kappa)^{-1}\right)=0=D\left(s_{o}(\kappa)^{-1}\right)$,

$$
C\left(s_{o}(\kappa)^{-1}\right)=\left(\begin{array}{ccc}
0 & 0 & -1  \tag{7.8.19}\\
0 & 1 & 0 \\
-1 & 0 & 0
\end{array}\right)
$$

and $B\left(s_{o}(\kappa)^{-1}\right)=-C\left(s_{o}(\kappa)^{-1}\right)$. Put $a=s_{o}(\kappa)^{-1} \bar{n}_{-f}(w)\left(g_{o} \exp t w_{o}\right) \bar{n}_{f}(w)^{-1}$ so that recalling (5.3.6), $a$ can be written

$$
\begin{equation*}
a=\bar{n} h n \tag{7.8.20}
\end{equation*}
$$

for $\bar{n} \in N, h \in H, n \in N$. In fact in the notation of Section 5.3, where $g$ is the element in $G_{0}{ }^{w} \subseteq \operatorname{Ad} g_{\mathbb{C}}$ corresponding to $g_{o}$ by (5.3.1), one has $\bar{n}=\bar{n}(g \exp t w)$, $h=h(g \exp t w)$, and $n=n(g \exp t w)$. We are particularly interested in
$h(g \exp t w)$ since, recalling Theorem 7.5 and (5.10.7), this determines the solution of the mechanical system.

Now let $M_{1}$ be the $3 \times 3$ matrix given by putting

$$
\begin{align*}
M_{1}= & C\left(\bar{n}_{-f}(w)\right) A\left(g_{o} \exp t w_{o}\right) A\left(\bar{n}_{f}(w)^{-1}\right) \\
& +D\left(\bar{n}_{-f}(w)\right) A\left(\left(g_{o} \exp t w_{o}\right)^{-1}\right)^{\prime} C\left(\bar{n}_{f}(w)\right)^{-1} \tag{7.8.21}
\end{align*}
$$

and put

$$
\begin{equation*}
M\left(g_{0}, w_{o} ; t\right)=F M_{1}, \tag{7.8.22}
\end{equation*}
$$

where

$$
F=\left(\begin{array}{rrr}
-1 & 0 & 0 \\
0 & 1 & 0 \\
0 & 0 & -1
\end{array}\right)
$$

Remark 7.8.3. Note that the $3 \times 3$ matrix $M\left(g_{o}, w_{o} ; t\right)$ has been determined here since all the components in (7.8.21) have been written down.

Proposition 7.8.2. For any $g_{o} \in H, w_{o} \in h_{+}$and $t \in \mathbb{R}, h(g \exp t w)^{v_{i}}$ for $i=1,2,3$ is equal to the determinant of the principal $i \times i$ minor of $M\left(g_{a}, w_{o} ; t\right)$.

Proof. By (7.8.20) one has

$$
\begin{equation*}
h(g \exp t w)^{v_{i}}=\left\{a v^{\nu_{i}}, v^{\nu_{i}}\right\} . \tag{7.8.23}
\end{equation*}
$$

But then the proposition follows from (7.8.14) and the definition of $a$.
Q.E.D.

Now let $x$ be a point in the phase space $\mathbb{R}^{6}$ and regard $x$ as the initial point for the time development of the system whose Hamiltonian is given by (7.8.1). If we put $b_{i}=p_{i}(x), i=1,2,3$, and $a_{j}=\exp \left(q_{j}(x)-q_{j+1}(x)\right), j=1,2$, and $a_{3}=\exp \left(2 q_{3}(x)\right)$ then the Jacobi element $y \in g$ corresponding to $x$ as defined in Theorem 7.5 is clearly given by

$$
y=\left(\begin{array}{cccccc}
b_{1} & a_{1} & 0 & 0 & 0 & 0  \tag{7.8.24}\\
1 & b_{2} & a_{2} & 0 & 0 & 0 \\
0 & 1 & b_{3} & \sqrt{2} a_{3} & 0 & 0 \\
0 & 0 & \sqrt{2} & -b_{3} & -a_{2} & 0 \\
0 & 0 & 0 & -1 & -b_{2} & -a_{2} \\
0 & 0 & 0 & 0 & -1 & -b_{1}
\end{array}\right)
$$

Now let $g_{o} \in H$ and $w_{o} \in h_{+}$be the "parameters" for $y$ as defined at the end of Section 3.7, recalling Theorem 3.7. Thus if $w_{o}$ is given by (7.8.4) then the eigenvalues of $y$ are $\pm w_{i}, i=1,2,3$, where $w_{1}>w_{2}>w_{3}>0$. To determine $g_{o}$ first let $n \in N$ be the unique element (see Proposition 2.3.2) such that
$n w n^{-1}=y$. In the notation of Theorem 3.7 one has $n=n(g)$. Thus $w$ and $y$ are explicitly given by (7.8.6) and (7.8.24), respectively, and $n$ is the unique upper triangular unipotent $6 \times 6$ matrix such that

$$
\begin{equation*}
n w=y n . \tag{7.8.25}
\end{equation*}
$$

Since $n$ is unipotent one inductively solves for $n_{i j}$. In fact what is needed here is the $3 \times 3$ matrix $B(n)$. Then one finds

$$
\begin{aligned}
& B(n)_{11}=\frac{a_{1}}{\sqrt{2}}\left(2 b_{1}+b_{2}-w_{1}-w_{2}-w_{3}\right)+\frac{\left(b_{1}-w_{1}\right)\left(b_{1}-w_{2}\right)\left(b_{1}-w_{3}\right)}{\sqrt{2}} \\
& B(n)_{21}=\frac{1}{\sqrt{2}}\left(\left(b_{2}-w_{3}\right)\left(b_{1}-w_{1}\right)+\left(b_{2}-w_{3}\right)\left(b_{2}-w_{2}\right)\right. \\
& \left.+\left(b_{1}-w_{1}\right)\left(b_{1}-w_{2}\right)+a_{1}+a_{2}\right), \\
& B(n)_{31}=\frac{1}{\sqrt{2}}\left(\left(b_{1}-w_{1}\right)+\left(b_{2}-w_{2}\right)+\left(b_{3}-w_{3}\right)\right), \\
& B(n)_{12}=-\frac{1}{\sqrt{2}}\left(a _ { 1 } \left[\left(b_{1}+w_{3}\right)\left(2 b_{1}+b_{2}-w_{1}-w_{2}-w_{3}\right)\right.\right. \\
& +\left(b_{2}-w_{3}\right)\left(b_{1}+b_{2}-w_{1}-w_{2}\right) \\
& \left.+\left(b_{1}-w_{1}\right)\left(b_{1}-w_{2}\right)+a_{1}+a_{2}\right] \\
& \left.+\left(b_{1}{ }^{2}-w_{1}{ }^{2}\right)\left(b_{1}-w_{2}\right)\left(b_{1}-w_{3}\right)\right), \\
& B(n)_{22}=-\frac{1}{\sqrt{2}}\left(a_{1}\left(2 b_{1}+2 b_{2}-w_{1}-w_{2}\right)+a_{2}\left(b_{1}+2 b_{2}+b_{3}-w_{1}-w_{2}\right)\right. \\
& +\left(b_{1}-w_{1}\right)\left(b_{1}-w_{2}\right)\left(b_{1}-w_{3}\right) \\
& +\left(b_{2}{ }^{2}-w_{3}{ }^{2}\right)\left(b_{1}+b_{2}-w_{1}-w_{2}\right) \\
& \left.+\left(b_{2}+w_{3}\right)\left(b_{1}-w_{1}\right) b_{1}-w_{2}\right), \\
& B(n)_{32}=-\frac{1}{\sqrt{2}}\left(a_{1}+a_{2}+2 a_{3}+\left(b_{2}-w_{3}\right)\left(b_{1}+b_{2}-w_{1}-w_{2}\right)\right. \\
& +\left(b_{1}-w_{1}\right)\left(b_{1}-w_{2}\right) \\
& \left.+\left(b_{3}+w_{3}\right)\left(b_{1}+b_{2}+b_{3}-w_{1}-w_{2}-w_{3}\right)\right), \\
& (B(n))_{13}=-\frac{1}{\sqrt{2}}\left(a_{1} a_{2}\left(w_{1}-b_{3}\right)+2 a_{1} a_{3}\left(w_{1}+b_{2}\right)+a_{2}^{2}\left(w_{1}-b_{1}\right)\right. \\
& -a_{1}\left(w_{1}{ }^{2}-b_{3}{ }^{2}\right)\left(w_{1}+b_{2}\right)-a_{2}\left(w_{1}-b_{1}\right)\left(w_{1}+b_{2}\right)\left(w_{1}+b_{3}\right) \\
& -2 a_{3}\left(\left(w_{1}-b_{1}\right)\left(w_{1}{ }^{2}-b_{2}{ }^{2}\right)\right. \\
& \left.+\left(w_{1}-b_{1}\right)\left(w_{1}{ }^{2}-b_{2}{ }^{2}\right)\left(w_{1}{ }^{2}-b_{3}{ }^{2}\right)\right),
\end{aligned}
$$

$$
\begin{aligned}
B(n)_{23}=-\frac{1}{\sqrt{2}}( & -2 a_{1} a_{3}+a_{1}\left(w_{1}^{2}-b_{3}^{2}\right)+a_{2}\left(w_{1}-b_{1}\right)\left(w_{1}+b_{3}\right) \\
& \left.+2 a_{3}\left(w_{1}-b_{1}\right)\left(w_{1}-b_{2}\right)-\left(w_{1}-b_{1}\right)\left(w_{1}-b_{2}\right)\left(w_{1}^{2}-b_{3}^{2}\right)\right) \\
B(n)_{33}=-\frac{1}{\sqrt{2}}( & -a_{1}\left(w_{1}-b_{3}\right)-a_{2}\left(w_{2}-b_{2}\right) \\
& \left.+\left(w_{1}-b_{1}\right)\left(w_{1}-b_{2}\right)\left(w_{1}-b_{3}\right)\right)
\end{aligned}
$$

As a consequence one has a formula for $g_{o}$ in terms of $y$ and $w_{o}$.
Proposition 7.8.3. Let $B_{i}(n)$, for $i=1,2,3$, be the $i \times i$ minor of the $3 \times 3$ matrix $B(n)$. (The matrix $B(n)$ is explicitly written down above.) Then if $g_{i}, i=1,2,3$, are the diagonal entries of $A\left(g_{o}\right)$ as in (7.8.17) one has

$$
\begin{align*}
& g_{1}=a_{1} a_{2} a_{3}^{1 / 2}\left(-\operatorname{det} B_{1}(n)\right) \\
& g_{2}=a_{2} a_{3}^{1 / 2} \operatorname{det} B_{2}(n)\left(\operatorname{det} B_{1}(n)\right)^{-1}  \tag{7.8.26}\\
& g_{3}=a_{3}^{1 / 2} \operatorname{det} B_{3}(n) \operatorname{det}\left(B_{2}(n)\right)^{-1}
\end{align*}
$$

Proof. Since $g$ is of type $C_{3}$ one knows that the Weyl group element $\kappa$ is given by

$$
\begin{equation*}
\kappa=-1 \tag{7.8.27}
\end{equation*}
$$

where 1 here is the identity element. Thus in the notation of Theorem 5.5 one has $g_{0}^{-\kappa}=g_{0}$. But then by (5.5.25)

$$
\begin{align*}
\left(h(g)^{-1} g_{o}\right)^{\nu_{i}} & =\left\{v^{\nu_{i}}, n(g) v^{k \nu_{i}}\right\} \\
& =-\operatorname{det} B_{i}(n) \tag{7.8.28}
\end{align*}
$$

by (7.8.14) and (7.8.16) since $n(g)=n$. This implies that $B_{i}(n)$ is invertible. But now by (7.8.5) and (7.8.13) one has

$$
\begin{align*}
\alpha_{1}+\alpha_{2}+\alpha_{3} / 2 & =\nu_{1}, \\
\alpha_{1}+2 \alpha_{2}+\alpha_{3} & =\nu_{2},  \tag{7.8.29}\\
\alpha_{1}+2 \alpha_{2}+3 \alpha_{3} / 2 & =\nu_{3},
\end{align*}
$$

On the other hand by (3.7.4) one has $a_{i}=h(g)^{-\alpha_{i}}, i=1,2,3$. Thus by (7.8.28) and (7.8.29)

$$
\begin{aligned}
& g_{o}^{v_{1}}=a_{1} a_{2} a_{3}^{1 / 2}\left(-\operatorname{det} B_{1}(n)\right), \\
& g_{o}^{y_{2}}=a_{1} a_{2}^{2} a_{3}\left(-\operatorname{det} B_{2}(n)\right), \\
& g_{o}^{y_{3}}=a_{1} a_{2} a_{3}^{3 / 2}\left(-\operatorname{det} B_{3}(n)\right) .
\end{aligned}
$$

But $g_{1}=g_{0}^{\nu_{1}}, g_{2}=g_{0}^{\nu_{2}-\nu_{1}}$, and $g_{3}=g_{0}^{\nu_{3}-\nu_{2}}$. This proves (7.8.26).
Q.E.D.

We can now integrate Hamilton's equations for the Hamiltonian (7.8.1). The only terms entering into the solution below which we have not written down here are the eigenvalues $\pm w_{i}, i=1,2,3$, for the matrix $y$. This of course reduces to solving a cubic equation.

Theorem 7.8.1. Let $x$ be a point in the phase space $\mathbb{R}^{6}$. Let $q_{i}(x(t)), i=1,2,3$, be the position of the ith particle at time $t$ of the mechanical system whose Hamiltonian is

$$
H=\sum_{i=1}^{3} \frac{p_{i}{ }^{2}}{2}+e^{q_{1}-q_{2}}+e^{q_{2}-q_{3}}+e^{2 q_{3}}
$$

and which at $t=0$ occupies the state $x$. Let $y$ be the $6 \times 6$ matrix given by (7.8.24), where $b_{i}=p_{i}(x), i=1,2,3$, and $a_{i}=\exp \left(q_{i}(x)-q_{i+1}(x)\right), i=1,2$, and $a_{3}=e^{2 q_{3}(x)}$. Then there uniquely exists $w_{1}>w_{2}>w_{3}>0$ such that $\pm w_{i}$, $i=1,2,3$, are the eigenvalues of $y$. Let $w_{o}$ be the corresponding $6 \times 6$ diagonal matrix defined by (7.8.3) and (7.8.4) and let $g_{o}$ be the $6 \times 6$ diagonal matrix given by (7.8.11), (7.8.17), and (7.8.26). Let $M\left(g_{o}, w_{o} ; t\right)$ be the $3 \times 3$ matrix defined by (7.8.21) and (7.8.22) and for $i=1,2,3$, let $M_{i}\left(g_{o}, w_{o} ; t\right)$ be the principal $i \times i$ minor of $M\left(g_{o}, w_{o} ; t\right)$. Then $\operatorname{det} M_{i}\left(g_{o}, w_{o} ; t\right)>0$ and one has

$$
\begin{aligned}
& q_{1}(x(t))=-\log \operatorname{det} M_{1}\left(g_{o}, w_{o} ;-t\right), \\
& q_{2}(x(t))=\log \left(\operatorname{det} M_{1}\left(g_{o}, w_{o} ;-t\right) / \operatorname{det} M_{2}\left(g_{o}, w_{o} ;-t\right)\right),
\end{aligned}
$$

and

$$
q_{3}(x(t))=\log \left(\operatorname{det} M_{2}\left(g_{o}, w_{o} ;-t\right) / \operatorname{det} M_{3}\left(g_{o}, w_{o} ;-t\right)\right) .
$$

Proof. The statement concerning the eigenvalues of $y$ follows from (7.8.5) and (3.3.4). We use formula (7.5.18) for $q_{i}(x(t))$. Since clearly $\mathscr{P}_{1}=\mathscr{P}$ in the present case one has $\bar{q}_{i}(x)=\bar{p}_{i}(x)=0$. Also, all $b_{i}$ (as defined in (7.5.13) and not as above) vanish by (7.5.14) since all $r_{i}=1$ ). But now the matrix $A=\left(a_{i j}\right)$ defined as in (7.1.10) is just

$$
\left(\begin{array}{rrr}
1 & -1 & 0  \tag{7.8.30}\\
0 & 1 & -1 \\
0 & 0 & 2
\end{array}\right)
$$

On the other hand $Q\left(\alpha_{i}, \alpha_{i}\right)=2, i=1,2$, and $Q\left(\alpha_{3}, \alpha_{3}\right)=4$. Thus by (7.5.18)

$$
\begin{aligned}
& q_{1}(x(t))=-\log \Phi_{1}\left(g_{o}, w_{a} ; t\right), \\
& q_{2}(x(t))=\log \left(\Phi_{1}\left(g_{o}, w_{o} ; t\right) / \Phi_{2}\left(g_{o}, w_{o} ; t\right)\right), \\
& q_{3}(x(t))=\log \left(\Phi_{2}\left(g_{o}, w_{a} ; t\right) / \Phi_{3}\left(g_{o}, w_{o} ; t\right)\right) .
\end{aligned}
$$

But now recalling (5.10.7) and the notation of Section 7.5, $\Phi_{i}\left(g_{0}, w_{o} ; t\right)=$ $h(g \exp (-t) w)^{\nu_{i}}$. But then the result follows from Proposition (7.8.2). Q.E.D.

The scattering of our mechanical system is explicitly given in
Theorem 7.8.2. For $i=1,2,3$ let $\boldsymbol{v}_{i}{ }^{+}(x) t+u_{i}{ }^{+}(x)$ and let $v_{i}-(x) t+u_{i}-(x)$ be the asymptotic lines defined by $\left(t, q_{i}(x(t))\right.$ as $t \rightarrow+\infty$ and $t \rightarrow-\infty$, respectively. Then the asymptotic velocities are given by $v_{i}^{+}(x)=-w_{i}$ and $v_{i}-(x)=w_{i}$. The corresponding phases are given by

$$
\begin{align*}
& u_{1}^{+}(x)=\log g_{1}+\log \frac{\left(w_{1}^{2}-w_{2}{ }^{2}\right)\left(w_{1}{ }^{2}-w_{3}{ }^{2}\right) 2 w_{1}}{\sqrt{2}}, \\
& u_{2}^{+}(x)=\log g_{2}+\log \frac{\left(w_{2}{ }^{2}-w_{3}{ }^{2}\right)\left(w_{1}+w_{2}\right) 2 w_{2}}{\sqrt{2}\left(w_{1}-w_{2}\right)},  \tag{7.8.31}\\
& u_{3}^{+}(x)=\log g_{3}+\log \frac{\left(w_{1}+w_{3}\right)\left(w_{2}+w_{3}\right) 2 w_{3}}{\sqrt{2}\left(w_{1}-w_{3}\right)\left(w_{2}-w_{3}\right)}
\end{align*}
$$

where $g_{1}, g_{2}$, and $g_{3}$ are given by (7.8.26) and one has

$$
\begin{equation*}
u_{i}^{-}(x)=u_{i}^{+}(x)-2 \log g_{i} \tag{7.8.32}
\end{equation*}
$$

for $i=1,2,3$.
Proof. We use Theorem 7.6 and formulas (7.6.1)-(7.6.3). By (7.8.30) and the relations which follow it one has

$$
\begin{equation*}
\mu_{1}^{-}=\nu_{1}, \quad \mu_{2}^{-}=\nu_{2}-\nu_{1}, \quad \mu_{3}^{-}=\nu_{3}-\nu_{2} \tag{7.8.33}
\end{equation*}
$$

and hence

$$
\begin{equation*}
\mu_{1}^{+}=-\nu_{1}, \quad \mu_{2}^{+}=\nu_{1}-\nu_{2}, \quad \mu_{3}^{+}=\nu_{3}-\nu_{2} . \tag{7.8.34}
\end{equation*}
$$

But then $\left\langle\mu_{i}^{-}, w_{o}\right\rangle=w_{i}$ and $\left\langle\mu_{i}{ }^{+}, w_{o}\right\rangle=-w_{i}$. The relation $v_{i}{ }^{+}(x)=-w_{i}$ and $v_{i}-(x)=w_{i}$ then follow from (7.6.2) and (7.6.3) since as noted in the proof of Theorem 7.8.1 one has $\bar{p}_{i}(x)=0$. But also, as noted there, $\bar{q}_{i}(x)=b_{i}=0$. Thus

$$
\begin{equation*}
u_{i}^{+}(x)=\log \left(\frac{d(w)}{g_{o}}\right)^{u_{i}^{+}} \tag{7.8.35}
\end{equation*}
$$

and

$$
\begin{equation*}
u_{i}^{-}(x)=\log \left(\frac{1}{g_{0} d(w)}\right)^{u_{i}^{-}} \tag{7.8.36}
\end{equation*}
$$

But then if $d(w)_{i}$ are the diagonal entrics of $A(d(w)), i=1,2,3$, as defined by (7.8.16) one has, by (7.8.5), (7.8.33), and (7.8.34)

$$
\begin{equation*}
u_{i}^{+}(x)=\log \left(\frac{g_{i}}{d(w)_{i}}\right) \tag{7.8.37}
\end{equation*}
$$

and

$$
\begin{equation*}
u_{i}-(x)=\log \left(\frac{1}{g_{i} d(w)_{i}}\right) . \tag{7.8.38}
\end{equation*}
$$

But then (7.8.31) follows from Proposition 7.8 .1 and (7.8.37). Relation (7.8.32) is immediate from (7.8.37) and (7.8.38). Q.E.D.

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[^1]:    ${ }^{1}$ Added in proof. A solution of the Toda lattice using representation theory has appeared in [29]. However, the solution in [29] is expressed in terms of integrals. See (18) in [29]. In effect, the results in Section 5 of our paper determine these integrals.

