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## HAMILTONIAN SYSTEMS WITH NILPOTENT STRUCTURES

## A thesis presented by <br> Malcolm Irving

to

The Department of Engineering<br>Control Theory Centre<br>in fulfillment of the requirements<br>for the degree of<br>Doctor of Philosophy<br>University of Warwick<br>Coventry

July 1983

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To my parents.

## SUMMARY

Symplectic Geometry has proved a powerful method in extending the knowledge of the classical theory of Hamiltonian mechanics without external variables. In this thesis these methods are applied to a class of Hamiltonian systems with controls in order to answer fundamental questions arising from Systems Theory and Classical Mechanics.

The theoretical aspects of this thesis deal with the extension of the Lie algebraic results of Engel and Lie on nilpotent and solvable Lie algebras, respectively, by the introduction of symplectic structures. It provides revealing results on the internal structure of symplectic vector spaces acted on by nilpotent or solvable Lie algebras. Then, using the methods of Kostant and Kirillov, these results are globalized to look at nilpotent transitive actions on simply connected symplectic manifolds and the consequent internal structures.

This theory is then applied to realizations of finite Volterra series with the additional property that the realization is Hamiltonian. These realizations are known to have an underlying nilpotent structure, A canonical realization is found and then shown to be closely linked with the theory of interconnections.

Finally, the concepts of complete integrability on free Hamiltonian systems is put into a feasible framework for Hamiltonian systems with controls which have an associated nilpotent Lie algebra. It is found that it is still possible to integrate these systems by quadratures but the structure is now much more complex.

## Chapter 1

## INTRODUCTION

Hamiltonian systems have been studied since early in the $19^{\text {th }}$ century starting with the equations of planetary motion. In more recent times there has been a great revival of interest in this area of Classical Mechanics with the advent of Symplectic Geometry, see for example [1], [2].

However, the systems investigated in the above texts are free systems i.e. with no external variables. Hamiltonian systems with controls and their corresponding properties of controllability and observability are a much more recent development, [3], [4]. This thesis aims to develop further the control theoretic aspects of this research to include an extension of the realization results of Crouch [5] to Hamiltonian systems. Also, it introduces a generalization of the classical complete integrability results of Arnold [2] to allow for the greater complexity of structure present in control systems.

## Realization Theory

In the linear case it is well known that if the input-output map has the form

$$
y(t)=\int_{0}^{t} W\left(t-\sigma_{1}\right) u\left(\sigma_{1}\right) d \sigma_{1}
$$

where $u(\cdot)$ is a $\mathbb{R}^{m}$ - valued input function and $y(\cdot)$ is a $R^{q}$ - valued output function and the Volterra kernel W satisfies

$$
W\left(t-\sigma_{1}\right)=H(t) G\left(\sigma_{1}\right)
$$

with $W$, $H$ and $G$ continuously differentiable maps, then it has a realization given by the linear system

$$
\begin{aligned}
& \dot{x}=A x+B u \quad x(0)=x_{0} \\
& y=C x
\end{aligned}
$$

where $x \in \mathbb{R}^{n}, \quad A \in \mathbb{R}^{\text {nxn }}, \quad B \in B^{n x m}, C \in B^{\text {qxn }}$ and conversely.

In Brockett and Rahimi [6] it is shown that if the Volterra kernel $W\left(t-\sigma_{1}\right)=-W\left(\sigma_{1}-t\right)$ then the realization is of the form

$$
\begin{array}{r}
\binom{\dot{q}}{\dot{p}}=\left(\begin{array}{cc}
A & B \\
-C & -A^{\prime}
\end{array}\right)\binom{q}{p}+\binom{\alpha}{-B} \quad \begin{array}{l}
B=B^{\prime} \\
C=C^{\prime}
\end{array}  \tag{1.1.1}\\
y=p^{\prime} \alpha+q^{\prime} \beta \quad(q, p)^{\prime} \in \mathbb{R}^{2 n} \\
q(0)=q_{0}, p(0)=p_{0}
\end{array}
$$

where $A, B, C, \alpha, B$ are appropriately dimensioned matrices.

This is called a Hamiltonian system as it consists of Hamiltonian vector fields on ( $\mathcal{B}^{2 n}, J$ ) where $J$ is the usual symplectic form

$$
J=\left(\begin{array}{cc}
0 & I_{n} \\
-I_{n} & 0
\end{array}\right)
$$

If Hamiltonian functions are given by
$\mathrm{H}=\frac{1}{2} q^{\prime} \mathrm{Cq}+\frac{1}{2} \mathrm{p}^{\prime} \mathrm{Bp}+\mathrm{p}^{\prime} \mathrm{Aq}$
$H_{u}=p^{\prime} \alpha+q^{\prime} \beta$
using standard notation (1.1.1) can be rewritten as

$$
\begin{aligned}
& \dot{q}=\left(\partial /_{\partial p}\left(H+u H_{u}\right)\right)^{\prime}, \dot{p}=-\left(\partial / \partial q\left(H+u H_{u}\right)\right), \\
& y=H_{u}
\end{aligned}
$$

Notice that if $u \equiv 0$ then this is exactly the form of Hamilton's equations.

The realizations considered in this thesis are a generalization of this linear case to that given by linear analytic systems of the form

$$
\begin{aligned}
\dot{x}=X_{H}(x)+\sum_{i=1}^{r} u_{i} X_{H_{u_{i}}}(x) \quad & x(0)=x_{0} \\
& x \in(M, \omega)
\end{aligned}
$$

$$
\begin{equation*}
y^{i}(t)=H_{u_{i}}(x) \tag{1.1.2}
\end{equation*}
$$

where each vector field is Hamiltonian and all the data is real analytic. In Brockett [7] conditions are given for finite Volterra Series of the form

$$
\begin{aligned}
& y(t)=W_{0}(t)+\int_{0}^{t} W_{1}\left(t, \sigma_{1}\right) u\left(\sigma_{1}\right) d \sigma_{1} \\
& +\sum_{i=1}^{n} \int_{0}^{t} \int_{0 \ldots}^{\sigma} \int_{0}^{\sigma_{i}} W_{i+1}\left(t, \sigma_{1}, \ldots, \sigma_{i+1}\right) u\left(\sigma_{1}\right) \ldots u\left(\sigma_{i+1}\right) d \sigma_{i+1} \ldots d \sigma_{1}
\end{aligned}
$$

to have linear analytic realizations.

Work done jointly by the author and Dr. P.E. Crouch give the further necessary and sufficient conditions for the realization to be of the form (1.1.2).

## Liouville's Theorem

As has been already mentioned the problem of planetary motion or the so called Kepler problem was one of the earliest known examples of a

Hamiltonian system. This is the problem of the motion of a particle in the gravitational field of a fixed point mass. It is modelled by the following system

$$
\begin{equation*}
\dot{x}=X_{H}(x) \quad x \in\left({\underset{\sim}{R}}^{6}, J\right) \tag{1.1.3}
\end{equation*}
$$

where $H(q, p)=\frac{1}{2} p^{2}-\alpha_{q}, \quad \alpha>0 \quad$ some constant
$x=\left(q_{1}, q_{2}, q_{3}, p_{1}, p_{2}, p_{3}\right)^{\prime}$ are the canonical coordinates of position and momentum on ${\underset{\sim}{R}}^{6}$ with

$$
p^{2}=\sum_{i=1}^{3} p_{i}^{2}, \quad q=\left(\sum_{i=1}^{3} q_{i}^{2}\right)^{\frac{1}{2}}
$$

Liouville's theorem, as presented by Arnold [2] gives conditions for the existence of a structure preserving map $\psi:{\underset{\sim}{R}}^{2 n}+{\underset{\sim}{R}}^{2 n}$ such that the Hamiltonian function $H$ in the new coordinates ( $I_{1}, \ldots, I_{n}, \phi_{1} \ldots, \phi_{n}$ ) on ${\underset{\sim}{R}}^{2 n}$ is of the form $H(I, \phi)=H(I)$ so the Hamiltonian vector field is

$$
\begin{equation*}
\dot{I}_{i}=\frac{\partial H}{\partial \phi_{i}}=0, \quad \dot{\phi}_{i}=-\frac{\partial H}{\partial \bar{I}_{i}}=\text { constant } \tag{1.1.4}
\end{equation*}
$$

These are the action-angle coordinates and in particular it is the case that the system (1.1.3) satisfies these conditions and thus under $\psi(1.1 .3)$ is integrable by quadratures. This high-lights the importance of choice of canonical coordinates for determining the solution of a system.

In Liouville's theorem on free systems the underlying structure is abelian. With the introduction of contrals this structure is normally lost. As a first step to find a generalization of Liouville's theorem with controls, nilpotent structures are considered and conditions sought for the existence of canonical coordinates for which the control system is integrable by quadratures.

Notation
The following notations will be used throughout this thesis assuming real analytic differentiability unless otherwise stated.

| R | The real numbers |
| :---: | :---: |
| $\sim^{R}$ | n - dimensional Euclidean space |
| $\mathrm{M}\left(\mathrm{M}^{\mathrm{n}}\right)$ | n - dimensional connected manifold |
| TM | Tangent bundle of $M$ |
| T*M | Cotangent bundle of M |
| $\mathrm{T}_{\mathrm{x}} \mathrm{M}$ | Tangent space to $M$ at x |
| $\mathrm{T}_{\mathrm{x}}{ }^{*} \mathrm{M}$ | Cotangent space to M at x |
| $v(M)$ | Set of all vector fields on M |
| $\Omega^{k}(M)$ | Set of all k -forms on M |
| $\left(M^{2 n}, \omega\right)$ | 2 n - dimensional connected symplectic manifold with symplectic form $\omega$ |
| C(M) | Ring of real valued functions on $M$ |
| Diff (M) | Group of diffeomorphisms of M |
| $l{ }_{l}$ p.c. | Class of piecewise constant functions on $[0, \infty)$ caking values in ${\underset{\sim}{\mathrm{R}}}^{\mathrm{m}}$. |

## Chapter 2

## SYSTEMS THEORY AND MECHANICS

This chapter contains the fundamental concepts and results that are needed in this thesis. Amplification of these can be found in Crouch [8] for systems theory and Abraham and Marsden [1] for the appropriate areas of Symplectic geometry.

### 2.1 Nonlinear Systems Theory

The class of system considered throughout this thesis is defined by the following equations.

$$
\begin{align*}
& \dot{x}=f(x)+\sum_{i=1}^{r} u_{i} g_{i}(x) \quad x(0)=x_{0}, x \in M^{n} \\
& y=h(x)
\end{align*}
$$

where the associated vector fields $f+\sum_{i=1}^{r} \alpha_{i} g_{i}$ for any $\left(\alpha_{1}, \ldots, \alpha_{r}\right) \in \mathbb{R}^{r}$, are complete analytic vector fields on $M^{n}$, an $n$-dimensional real analytic connected manifold, and $h$ is $a \underset{\sim}{R^{q}}$ valued analytic function on $M$.

Since the associated vector fields are complete,solutions to the above equations are defined on $[0, T]$ for all piecewise constant controls $u$ on $[0, T]$, and all positive times $T$. Sussmann [9] enables an extension to include measurable controls. However, piecewise constant controls, uc $\ell \ell_{p . c .}$, shall be used throughout.

For the case where the Lie algebra of (2.1.1) is finite dimensional, denoted by

$$
\mathcal{L}=\left\{f, g_{1}, \ldots, g_{r}\right\}_{L . A}
$$

Palais [10] guarantees that $\mathcal{L}$ consists of complete analytic vector fields on $M$.

All Lie algebras throughout this thesis associated to systems of the form (2.1.1) are assumed to be finite dimensional unless otherwise specified.

Definition 2.1.1 A distribution $\Delta$ on a manifold $M$ is the assigment to each $x \in M$ of a subspace $\Delta(x)$ of $T_{X} M, \Delta$ is said to be a $C^{K}(K=\infty, \omega)$ distribution if it is spanned by a family $D$ of $C^{K}$ vector fields i.e. if $\Delta(x)$ is the linear subspace of $T_{X} M$ spanned by the family $\{X(x)\}$ with $X \in D$ for every $x \in M$.
$\Delta$ is said to be m-dimensional if $\Delta(x)$ has dimension $m$ for each $x \in M$.

Definition 2.1.2 $A C^{K}$ - distribution $\Delta$ is involutive if the vector fields $X, Y \in \Delta$ then $[X, Y] \in \Delta$ i.e. $X(x), Y(x) \in \Delta(x) \Rightarrow[X, Y](x) \in \Delta(x)$, $x \in M$.

Definition 2.1.3 A submanifold $N$ of $M$ is said to be an integrable submanifold of $\Delta$ if $T_{x} N=\Delta(x) \forall x \in N$, and a maximal integral submanifold of $\Delta$ is a connected integral stbmanifold of $\Delta$ maximal for the relation of inclusion.

Further, $\Delta$ is integrable if through every point of $M$ there passes a maximal integral submanifold of $\Delta$.

Theorem 2.1.4 (Nagano [11]). An analytic involutive distribution is integrable.

Theorem 2.1.5 (Frobenius [12]). A $C^{\infty}$ constant dimensional involutive distribution is integrable.

The associated vector fields for system (2.1.1) are the vector fields
$x \longmapsto f(x)+\sum_{i=1}^{r} u_{i} g_{i}(x)$ for each fixed $u \in \ell \ell_{p . c \text {. }}$, assumed
throughout to be complete, and denote the set of associated vector fields by F .

The group $G_{F}$ (semi-group $S_{F}$ ) of $F$ is defined to be the sub (semi -) group of diffeomorphisms of $M$ of the form

$$
\begin{aligned}
& G_{F}=\left\{g \in \operatorname{Diff}(M): g=\gamma_{X_{1}}\left(t_{1}\right)_{0 \ldots} \ldots o_{X_{n}}\left(t_{n}\right),\right. \\
&\left.X_{i} \in F, t_{i} \in \underset{\sim}{R}, n \in{\underset{\sim}{z}}^{+}\right\} \\
& S_{F}=\left\{g \in \operatorname{Diff}(M): g=\gamma_{X_{1}}\left(t_{1}\right)_{0 \ldots 0^{\gamma} X_{n}}\left(t_{n}\right),\right. \\
&\left.X_{i} \in F, t_{i} \in{\underset{\sim}{R}}^{+}, n \in{\underset{\sim}{z}}^{+}\right\},
\end{aligned}
$$

where $\gamma_{X_{i}}\left(t_{i}\right)$ is the flow of $X_{i} \in F$.

Definition 2.1.6 A linear analytic system is said to be accessible if $S_{F}(x)$ has nonempty interior in $M, \forall x \in M$.

The following theorem, Krener and Hermann [13], gives algebraic conditions for accessibility as well as a proof of importance in a later chapter.

Theorem 2.1.7 A linear analytic system is accessible if and only if $\operatorname{dim} L(x)=n=\operatorname{dim} M, \forall x \in M$.

Proof Let $x \in M$ and $U$ a neighbourhood of $x$. An open set $V$ of $M$ is constructed such that $V \subset U$ and each point $V \in V \subset S_{F}(x)-S_{F}$ system semi-group- is joined to $x$ by a trajectory of the system remaining in $U$. In fact a sequence of submanifolds $V_{j}$ of dimension $j$ are constructed with these properties as follows:

For $j=1$ choose $X_{1} \in F$ such that $X_{1}(x) \neq 0$ and let $(t, x)+\gamma_{1}(t) x$ be the flow of $X_{1}$. Choose $\varepsilon>0$ such that $\gamma_{1}\left(r_{1}\right) x \in U$ for $0 \leq t_{1}<\varepsilon$ then $V_{1}=\left\{\gamma_{1}\left(t_{1}\right) x: 0<t_{1}<\varepsilon\right\}$ is a submanifold of dimension one which satisfies the requirements. $V_{j}$ is now constructed by induction.

Suppose $V_{j-1}$ satisfies the above requirements and is defined by
$v_{j-1}=\left\{\gamma_{j-1}\left(t_{j-1}\right)_{0 \ldots 0_{1}}\left(t_{1}\right) x:\left(t_{\left.j-1, \ldots, t_{1}\right)} \in\right.\right.$
an open subset of the positive orthant in ${\underset{\sim}{R}}^{\mathrm{j}-1}$ \} where $(t, x) \rightarrow \gamma_{k}\left(t_{k}\right) x$ is the flow of a vector field $X_{k} \in F$

Now if $j<m$ find $X_{j} \in F$ and $v \in V_{j-1}$ such that $X_{j}(v) \notin T_{v} V_{j-1}$. If this was not possible $F_{v} \& T_{v} V_{j-1}, \Psi v \in V_{j-1}$ and so $\hat{L}(v) \subset T_{v} V_{j-1}$, $\forall v \in V_{j-1}$ which implies that the distribution defined by $\mathcal{L}$ has dimension $j-1$ on $V_{j-1} \subset U$ which is a contradiction. It follows that
the map $\left(t_{j}, \ldots, t_{1}\right) \rightarrow \gamma_{j}\left(t_{j}\right) \quad \ldots o^{\gamma_{1}}\left(t_{1}\right) x$ is of full rank on some open subset of the positive orthant of $\mathrm{R}^{j}$. Thus by restriction if necessary it can be assumed that the range $V_{j}$ of this map is a submanifold of $M$ of dimension $j$ with the required properties. The desired set $V$ can now be taken to be $\mathrm{V}^{\text {m. }}$.

Conversely, if for $x \in M$ it is assumed that $r=\operatorname{dim} \mathcal{L}(x)<\operatorname{dim} T_{x} M=n$ then there exists a maximal integral submanifold $N$ of dimension $r$ passing through $x$. Using Nagano's theorem 2.1 .4 clearly the reachable set $S_{F}(x)$ of the system from $x$ is contained in $N$ and so cannot have non-empty interior in $M$ with respect to the topology on $M$. Therefore the system cannot be accessible, which is a contradiction.
Q.E.D.

Denote the reachable set from $x_{0}$ ar time $T \geqslant 0$ by $R\left(T, x_{0}\right)$, it is defined to be the set of all points $x_{1} \in M$ such that there exists a piecewise constant control $u(\cdot)$ on $[0, T]$ and corresponding solution $x(\cdot)$ on $[0, T]$ such that $x(T)=x_{1}$.

The reachable set from $x_{0}$, denoted by $R\left(x_{0}\right)$ is defined by

$$
\begin{gathered}
R\left(x_{0}\right)=U_{R\left(T, x_{0}\right)} \\
T \geqslant 0
\end{gathered}
$$

Notice that $S_{\bar{r}}\left(x_{0}\right)=R\left(x_{0}\right)$.

Definition 2.1.8 A linear analytic system is said to be strongly accessible if for each $x \in M$ there exists a $T>0$ such that $R(x, T)$ has a nonempty interior in $M$.

Let $S$ denote the ideal of $\mathcal{L}$ generated by $\left\{g_{i}\right\}_{\text {L.A. }}{ }^{i}$ $s$ has codimension one or zero in $\mathcal{L}$.

Corollary 3.4 Sussmann [14], shows that $S(x)$ has constant dimension on $M$ if the system is accessible and the following theorem, Sussmann [14], is a characterization of strong accessibility.

Theorem 2.1.9 A linear analytic system is strongly accessible if and only if $\operatorname{dim} S(x)=n=\operatorname{dim} M, \forall x \in M$. In this case $R(x, T)$ has nonempty interior in $M \forall T>0$.

An immediate corollary from the proof of theorem 4.3 in Crouch [5] for systems which are strongly accessible and such that $S=L$ is the following

Theorem 2.1.10 Given a strongly accessible linear analytic system of the form

$$
\begin{equation*}
\hat{x}=f(x)+\sum_{i=1}^{r} u_{i} g_{i}(x) \quad x(0)=x_{0}, x \in M^{n} \tag{2.1.2}
\end{equation*}
$$

such that $S$ is a finite dimensional nilpotent Lie algebra and $S=\mathcal{L}$. Then there exists an analytic diffeomorphic change of coordinates such that (2.1.2) takes the form

$$
\begin{array}{lll}
i_{1}=f_{1} & +\sum_{i=1}^{r} u_{i} g_{1}^{i} & z_{1}(0)=0 \\
i_{2}=f_{2}\left(z_{1}\right) & +\sum_{i=1}^{r} u_{i} g_{2}^{i}\left(z_{1}\right) & z_{2}(0)=0 \\
i_{j}=f_{j}\left(z_{\left.1, \ldots,{ }_{j-1}\right)}\right. & +\sum_{i=1}^{r} u_{i} g_{j}^{i}\left(z_{\left.1, \ldots, z_{j-1}\right)} z_{j}(0)=0\right. \\
i_{m}=f_{m}\left(z_{\left.1, \ldots, r_{m-1}\right)}\right. & \quad+\sum_{i=1}^{r} u_{i} g_{m}^{i}\left(z_{\left.1, \ldots, z_{w-1}\right)} z_{m}(0)=0\right.
\end{array}
$$

where $f_{j}, g_{j}^{i}$ are vector valued polynomials of the components of the vectors $z_{i} \in \stackrel{R}{\sim}^{n_{i}}, \sum_{i=1}^{m} n_{i}=n$, being the dimension of the state space.

If a linear analytic system is accessible then $G$, the connected Lie group corresponding to $L$ the Lie algebra generated by the system, acts transitively on $M$, the state space. By a standard result, Hel gason [15], $M$ is therefore analytically diffeomorphic to the homogeneous space $G / G_{x_{0}}$ where $G_{X_{0}}$ is the isotropy group

$$
G_{x_{0}}=\left\{g \varepsilon G: g \cdot x_{0}=x_{0}\right\}
$$

If the system is strongly accessible the connected Lie subgroup $N$ of $G$, corresponding to $S$ also acts transitively on $M$, and so $M$ can be expressed as a homogeneous space $N / N_{x_{0}}$.

A number of concepts of observability will also be required throughout this thesis.

Definition 2.1.11 Two points $x_{0}, x_{1} \in M$ are said to be indistinguishable if the input-output maps for the system initialized at $x_{0}$ and $x_{1}$ are equal for all $u(\cdot) \subset \ell_{p . c .}$.

Definition 2.1.12 $A$ system is said to be observable if $x_{0}$ and $x_{1} \in M$ being indistinguishable implies $x_{0}=x_{1}$.

Definition 2.1.13 A system is said to be weakly observable if for all states $x_{0} \in M$ there exists a neighbourhood $U$ of $x_{0}$, such that if $x_{1} \in U$ is indistinguishable from $x_{0}$ then $x_{0}=x_{1}$.

Let $\mathcal{H}$ denote the smallest linear subspace of $C(M)$ containing the functions $h_{i}, i=1, \ldots, q$ and closed under Lie differentiation by elements of $\mathcal{L}$. Thus $\mathscr{K}$ consists of all linear combinations of the functions.

$$
L_{x_{1}}\left(L _ { x _ { 2 } } \left(\ldots\left(L_{x_{n}}\left(h_{i}\right) \ldots\right), x_{i} \in \mathcal{L}\right.\right.
$$

The following theorem is found in Krener and Hermann [13],

Theorem 2.1.14 A linear analytic accessible system is weakly observable if and only if

$$
T_{x}{ }^{*} M=d \mathscr{X}(x)=\{d \tau(x): \tau<\forall\}, \forall x \in M \text { where } T_{x}{ }^{*} M \text { is the }
$$ cotangent space to $M$ at $x$.

Given two systems $\sum_{1}$ and $\sum_{2}$, with the same control class and output space, $x_{1} \in M_{1}$ and $x_{2} \in M_{2}$ are said to be indistinguishable if $h_{1} o^{g}\left(x_{1}\right)=h_{2} o_{2}\left(x_{2}\right)$ where $g_{1} \in \quad G_{F_{1}}$ and $g_{2} \in G_{F_{2}}$ is obtained using the same piecewise constant control as the one involved in $g_{1}$.

Definition 2.1.15 $\Sigma_{1}$ and $\Sigma_{2}$ are strongly equivalent if every state in $\Sigma_{1}$ is indistinguishable from some state in $\Sigma_{2}$ and conversely.

For further results on strong equivalence refer to Goncalves [16].

### 2.2 Basic theory of Classical Mechanics

Most of the work in this thesis depends on a knowledge of Symplectic Geometry. This section presents the basic concepts. Good references for a more extensive survey are Abraham and Marsden [1] and Arnold [2].

Definition 2.2.1 Let $V$ be a vector space and $\omega$ a non-degenerate two form on $V$. Then the pair ( $V, \omega$ ) is said to be a symplectic vector space.

A globalization of this gives the following definition.

Definition 2.2.2 Let $M$ be a manifold with a nondegenerate closed two form $\omega$ defined on $M$. The pair $(M, \omega$ ) is said to be a symplectic manifold. $\omega$ is often called a symplectic form.

The non-degeneracy of $\omega$ guarantees that $M$ is even dimensional. One of the most fundamental examples of a symplectic manifold is the cotangent bundle of a manifold i.e. $M=T^{*} Q$. This has a naturally defined symplectic structure as follows:

Let $\pi: \mathbb{I}^{*} Q \longmapsto Q$ be the cotangent bundle projection. Define the 1 - form $\theta$ in $T^{*} Q$ by

$$
\theta(\alpha)(x)=\alpha_{0} \pi_{*}(x)
$$

where $x \in T \alpha^{T^{*}} 0$ and $\pi_{*}$ denotes the push forward.
Then the natural symplectic form $\omega$ is given by $\omega=-\mathrm{d} \theta$.

The symplectic manifolds throughout this thesis are assumed to be finite dimensional unless otherwise specified.

Locally, the above forms can be written as
$\theta=\sum_{i=1}^{n} p_{i} d q_{i}$
$\omega$

$$
=\sum_{i=1}^{n} d q_{i} \cdot d p_{i}
$$

where $q_{1}, \ldots, q_{n}$ are coordinates on $Q$ and $\left(q_{1}, \ldots, q_{n}, p_{1}, \ldots, p_{n}\right)$ are canonical coordinates on $M=I^{*} Q$.

Definition 2.2.3 Let $\left(M_{1}, \omega_{1}\right)$ and $\left(M_{2}, \omega_{2}\right)$ be symplectic manifolds. A $C^{K}$ - mapping $f: M_{1} \longmapsto M_{2}$ is called either symplectic or a symplectomorphism if $f^{\star} \omega_{2}=\omega_{1}$.

The definition of locally symplectic follows obviously from the above.

Definition 2.2.4 Let $(M, \omega)$ be a symplectic manifold and $H \in C^{K}(M)$. The Hamiltonian vector field $X_{H}$ on $M$ is defined by

$$
i\left(X_{H}\right) w=d H
$$

where $i\left(X_{H}\right) \omega$ is the inner product of $X_{H}$ and $\omega$. Locally, this corresponds to

$$
\begin{array}{ll}
\dot{q}_{i}=\frac{\partial H}{\partial p_{i}} & i=1, \ldots, n \\
\dot{p}_{i}=-\frac{\partial H}{\partial q_{i}} & i=1, \ldots, n
\end{array}
$$

where $\left(q_{1}, \ldots, q_{n}, p_{1, \ldots,} P_{n}\right)$ are canonical coordinates on $M$ and
$\omega=\sum_{i=1}^{n} d q_{i} \cdot d p_{i}$.

These are Hamilton's equations.

Note that $H$ is called the Hamitonian function of the Hamilconian vector field $\mathrm{X}_{\mathrm{H}}$.

Hamiltonian vector fields yield an inportant example of a symplectomorphism. For if $X_{H}$ is a Hamiltonian vector field on ( $M$, w) then its flow $\gamma(t)$ is a symplectomorphism.

A weaker version of Hamiltonian vector field is a locally Hamiltonian vector field on ( $M, \omega$ ) which is a vector field $X$ on ( $M, \omega$ ) such that for every $x \in M$ there exists a neighbourhood $U$ of $x$ such that $X$ restricted to $U$ is Hamiltonian. An equivalent form of this is that $L_{X}{ }^{\omega}=0$, where $L_{X} \omega$ is the Lie derivative of $\omega$ with respect to $X$.

Definition 2.2.5 Let $(M, W)$ be a symplectic manifold and $f, g \in C^{K}(M)$ then the Poisson bracket of $f$ and $g$ is given by
$\{f, g\}=\omega\left(X_{f}, X_{g}\right)$
In canonical coordinates ( $q_{1, \ldots, q_{n}}, P_{i, \ldots,} P_{n}$ ) on ( $M, \omega$ ) this is given by

$$
\{f, g\}=\sum_{i=1}^{n} \frac{\partial f}{\partial q_{i}} \frac{\partial g}{\partial p_{i}}-\frac{\partial f}{\partial p_{i}} \frac{\partial g}{\partial q_{i}}
$$

Proposition 2.2.6 The real vector space $C(M)$, together with the Poisson bracket \{, \} forms a Lie algebra.

A number of mathematical identities will be required repeatedly. These are listed here,
(i) $\quad d \mathrm{~d} \alpha=0, \mathrm{~d}(\alpha, \beta)=\mathrm{d} \alpha, \beta+(-1)^{\mathrm{k}} \alpha_{\wedge}, \mathrm{d} \beta, \alpha-\mathrm{k}$-form
(ii) $\quad f^{*} d a=d f^{*} \alpha, f: M \longmapsto M$.
(iii) $L_{X} \alpha=d i(X) \alpha+i(X) d \alpha, X \in V(M)$
(iv) fadiffeomorphism, $f^{*} L_{X} \alpha=L_{f_{*}} I_{X} f^{*} \alpha$
(v) $\quad\left[X_{f}, X_{g}\right]=X_{\{f, g\}}$ on $(M, \omega)$.

Proofs and further details may be found in chapters 2 and 3 of Abraham and Marsden [1].

One of the main theorems in Classical Mechanics, and certainly one which has attracted much attention in recent years, is one associated with Liouville and Arnold. The form it is stated in here follows Arnold [2].

Theorem 2.2.7 Let $(M, \omega)$ be a $2 n$ - dimensional symplectic manifold, and consider the system

$$
\begin{equation*}
\dot{x}=X_{H}(x) \quad x(0)=x_{0}, x \in M \tag{2.2.1}
\end{equation*}
$$

Suppose there exists $n$ - functions $f_{1, \ldots,} f_{n}$ on $M$ in involution i.e. $\left\{f_{i}, f_{j}\right\}=0,1 \leqslant i, j \leqslant n$ with $f_{1}=H$ and let

$$
M_{f(c)}=\left\{x: f_{i}(x)=c_{i}, \quad i=1, \ldots, n\right\}
$$

with the $n$ - one - forms $d f_{i}$ linearly independent at each point of

$$
M_{f(c)} \text {. Then }
$$

(i) $\quad M_{f(c)}$ is a $n$ - dimensional submanifold of $M$ invariant under the flow of the Hamiltonian vector field $X_{H}$ on ( $M, \omega$ ).
(ii) If $M_{f(c)}$ is compact and connected, then it is diffeomorphic to the $n$ - dimensional torus $\mathrm{T}^{\mathrm{n}}$.
(iii) There exists canonical coordinates ( $I_{1}, \ldots, I_{n}, \phi_{1}, \ldots, \phi_{n}$ ) on ( $M, \omega$ ) such that $H=H\left(I_{1}, \ldots, I_{n}\right)$ called the action angle coordinates.
(iv) The canonical equations with Hamiltonian function $H$ can be integrated by quadratures.

An excellent proof of this theorem may be found in Vinogradov and Kuperschmidt [17].

Throughout che field of Classical Mechanics and Symplectic Geometry much work has been done on finding systems of the form (2.2.1) which have $n$ - functions in involution or perhaps fewer functions in involution but n - functions which commute with $H$ under Poisson bracket see Nehorosev [18]. Much work has also been done in this area using the methodology of Kostant [19], Souriau [20] and Kirillov [21], see for instance Guillemin and Sternberg [22] or Marsden and Weinstein [23].

The basic difference between this and the work presented in later chapters is that conmuting Hamiltonian functions under Poisson bracket are no - longer considered. The following example shows that for control
systems the underlying structure is much more complex.

Example 2.2.8 Consider the linear analytic dynamical system on $\left({\underset{\sim}{R}}^{4}, \sum_{i=1}^{2} \operatorname{dq}_{i}{d p_{i}}\right.$ ) with Hamiltonian functions defined by

$$
\mathrm{H}=\mathrm{p}_{1}+\frac{1}{2} p_{2} q_{1}^{2}
$$

$$
H_{u}=p_{1}
$$

$$
\text { so } \dot{x}=X_{H}(x)+u X_{H}(x) \quad x(0)=x_{0} \quad x \in\left({\underset{\sim}{R}}^{4}, \omega\right)
$$

$$
x=\left(q_{1}, q_{2}, p_{1}, p_{2}\right)^{\prime}
$$

That is $\dot{\mathbf{q}}_{1}=1+u$

$$
q_{1}(0)=q_{1}^{0}
$$

$$
\begin{array}{ll}
\dot{q}_{2}=\frac{1}{2} q_{1}^{2} & q_{2}(0)=q_{2}^{0}  \tag{2,2,2}\\
\dot{p}_{1}=-p_{2} q_{1} & p_{1}(0)=p_{1}^{0} \\
\dot{p}_{2}=0 & p_{2}(0)=p_{2}^{0}
\end{array}
$$

the Poisson brackets of the Hamiltonian functions are found as follows,

$$
\left\{H_{u}, H\right\}=\sum_{i=1}^{2} \frac{\partial H_{u}}{\partial q_{i}} \frac{\partial H}{\partial p_{i}}-\frac{\partial H_{u}}{\partial p_{i}} \frac{\partial H}{\partial q_{i}}=-p_{2} q_{1}
$$

$$
\left\{H_{u}\left\{H_{u}, H\right\}\right\}=P_{2}=\left\{H\left\{H_{u}, H\right\}\right\}
$$

All other Poisson brackets vanish.

Notice that although $H$ and $H_{u}$ do not commute under Poisson brackets the system (2.2.2) can still be integrated by quadratures and further the level surface

$$
N=\left\{x \in \underline{R}^{4}: f(x)=0 \text { for all } f\right. \text { such that }
$$ $\left.X_{f} \in \mathcal{L}, \mathcal{L}=\left\{X_{H}, X_{H_{u}}\right\}_{L . A}\right\}$ is a 2 - dimensional submanifold of the symplectic manifold $\left(\underset{\sim}{R^{4}}, \omega\right)$ described by $p_{1}=p_{2}=0$.

Also, on calculation of the Lie brackets of $\mathcal{l}, \mathcal{l}$ spans the tangent space of $N$ at each point of $N$ since

$$
\begin{aligned}
& {\left[X_{H_{u}}, X_{H}\right] }=\left[\frac{\partial}{\partial q_{1}}, \partial / \partial q_{1}+\frac{1}{2} q_{1}{ }^{2} / \partial q_{2}-p_{2} q_{1}{ }^{\partial} / \partial p_{1}\right] \\
&=q_{1}{ }^{\partial} / \partial q_{2} \\
& {\left[p_{2}{ }^{\partial} / \partial p_{1}\right.} \\
& {\left[X_{H_{u}}\left[X_{H_{u}} X_{H}\right]\right]=\left[\frac{\partial}{\partial q_{1}}, q_{1}{ }^{2} / \partial q_{2}-p_{2}{ }^{\partial} / \partial p_{1}\right] } \\
&=\partial / \partial q_{2}=\left[X_{H}\left[X_{H_{u}}, X_{H}\right]\right]
\end{aligned}
$$

So $X_{H_{u}}(x)$ and $\left[X_{H_{u}}\left[X_{H_{u}} X_{H}\right]\right](x)$ span $T_{X} N, \forall x \in N$, all other Lie brackets vanish, so by theorem 2.1.7 this system is not accessible on $\mathbb{R}^{4}$ but it is on the level surface $N$.

It is this type of generalization of Liouville's theorem that will come under investigation in Chapter 5.

## Chapter 3

## GENERAL THEORY OF NILPOTENT STRUCTURES

## ON SYMPLECTIC MANIFOLDS

This chapter begins with a section investigating the presence of nilpotent and solvable structures on a symplectic vector space extending the classical Lie algebraic results of Angel and Lie. The second section then considers a globalization to symplectic manifolds of the above. The appropriate conditions are found on a sequence of involutive distributions for the existence of canonical coordinates which have the same properties as found in section one. Finally, in the third section the existence of such sequences of involutive distributions is investigated using the techniques of Kostant [19] and Kirillov [21].

### 3.1 Hamiltonian Endomorphisms

This section investigates the new properties which arise when a symplectic structure is placed on the classical theorems of Angel and Lie, see Humphreys [24].

Let $L$ be a Lie algebra. Define a sequence of ideals by

$$
\mathcal{L}^{1}=L, L^{2}=[L, L], \ldots, L^{k}=\left[L, L^{k-1}\right]
$$

Recall that $\mathcal{L}$ is said to be nilpotent if $\mathcal{L}^{n}=0$ for some $n$.
If $L$ is any Lie algebra $X \in \mathcal{R}$ is ad-nilpotent if ad $_{X}$ is a nilpotent endomorphism, ie. $\left(\mathrm{ad}_{\mathrm{X}}\right)^{n}=0$ for some $n$. If each element of a Lie algebra $L$ is ad-nilpotent it is said to consist of nilpotent endomorphisms.

Proposition 3.1.1 Let $\mathcal{L}$ be a Lie algebra.
(i) If $\mathcal{L}$ is nilpotent then so are all sub and quotient algebras (ii) If $L / Z(\mathcal{L})$ is nilpotent, so is $L$, where $Z(L)$ is the centre of $L$. (iii) $\mathcal{L}$ nilpotent, $L \notin 0 \Rightarrow z(\mathcal{L}) \neq 0$.

Theorem 3.1.2 (Engel) Let $V$ be a non-zero finite dimensional vector space over $\underset{\sim}{R}$, and let $\mathcal{L}$ be a subalgebra of $g 1(V)$ consisting of nilpotent endomorphisms. Then
(i) $\mathcal{L}$ is nilpotent
(ii) There exists a vector $v \neq 0$ in $V$ such that $X, v=0$ for all $x \in \mathcal{L}$.
(iii) There exists a basis $e_{1, \ldots,} \mathbf{e}_{\mathbf{n}}$ in $V$ such that all the endomorphisms $\mathrm{X} \in \mathcal{L}$ are expressed by matrices with zeros on and above the diagonal.

It is part (iii) of Engel's theorem with vector space replaced by symplectic vector space that will be of particular interest. First some notation,

Definition 3.1.3 If $V$ is a vector space then a flag in $V$ is a sequence of subspaces

$$
0=v_{0}=v_{1} \subset \ldots \subset v_{n-1} \subset v_{n}=v
$$

such that $\operatorname{dim} V_{i}=\mathbf{i}$.
If $X \in \operatorname{End}(V)$ and $X, V_{i} \subset V_{i}$ then $X$ is said to stabilize the flag.

In terms of flags theorem 3.1.2(iii) could be rewritten as, there exists a flag in $v$ such that $\mathcal{L} \cdot v_{i} \subset v_{i-1}$. Since $v_{i-1} \subset v_{i}, \mathcal{L}$ stabilizes the flag.

Definition 3.1.4 Let $(V, w)$ be a symplectic vector space, then a linear mapping $X \in L(V, V)$ is infinitesimally symplectic with respect to the symplectic form $\omega$ if

$$
\omega\left(X v, v^{\prime}\right)+\omega\left(v, X v^{\prime}\right)=0
$$

for all $v, v^{\prime} \in V$.

Denote the set of all linear mappings in $L(V, V)$ that are infinitesimally symplectic with respect to $\omega$ by $\operatorname{sp}(V, \omega)$.

Infinitesimal symplectic linear mappings are often referred to as linear Hamiltonian mappings.

An ordered basis can be chosen for $V$ such that the matrix of $\omega$ is

$$
J=\left(\begin{array}{cc}
0 & I \\
-I & 0
\end{array}\right)
$$

where, if $\operatorname{dim} V=2 n, I$ is the $n \times n$ identity matrix. Thus definition 3.1.4 says that $X \in \operatorname{sp}(V, J)$ if

```
X'J + JX = 0
```



The matrix $\left(\begin{array}{ll}x_{1} & x_{2} \\ x_{3} & x_{4}\end{array}\right)$ is said to be a Hamiltonian matrix.

If $X_{1, \ldots} X_{n}$ are infinitesimally symplectic their Lie algebra
$\mathcal{L}=\left\{X_{1, \ldots, X_{n}}\right\}_{L . A}$ is said to be a Lie algebra of Hamiltonian endomorphisms i.e. it is clear that each element of $\mathcal{L}$ is infinitesimally symplectic. In fact $s p(V, w)$ is a Lie subalgebra of $g l(V)$.

Definition 3.1.5 Let $(V, \omega)$ be a symplectic vector space and $U \subset V$ a subspace. The w-orthogonal complement of $U$ is the subspace defined by

$$
U^{\perp}=\{v \in V: w(v, u)=0 \quad \downarrow u \in U\}
$$

U is said to be
(i) isotropic if $U \subset U^{\perp} \quad$ i.e. $\omega\left(u, u^{\prime}\right)=0 \quad \psi u, u^{\prime} \in U$
(ii) co-isotropic if $U^{\perp} \in U$ i.e. $\omega\left(u, u^{\prime}\right)=0 \forall u^{\prime} \in U \Rightarrow u \in U$
(iii) Lagrangian if $U$ is isotropic and has an isotropic complement, i.e. $V=U \oplus U^{\prime}$ where $U^{\prime}$ is isotropic.

The following two propositions the proofs of which appear for example, in Abraham and Marsden [1], are of continual use.

Proposition 3.1.6 Suppose $U_{1}, U_{2} \simeq V$ are subspaces
(i) $\mathrm{U}_{1} \subset \mathrm{U}_{2} \Rightarrow \mathrm{U}_{2}^{\perp} \subset \mathrm{U}_{1}^{+}$
(ii) $U_{1}^{\perp} \cap u_{2}^{\perp}=\left(U_{1}+U_{2}\right)^{\perp}$
(iii) $\operatorname{dim} V=\operatorname{dim} U_{1}+\operatorname{dim} U_{1}^{\perp}$
(iv) $\mathrm{U}_{1}=\mathrm{U}_{1}^{\perp \perp}$
(v) $\left(U_{1} \cap U_{2}\right)^{\perp}=U_{1}^{\perp}+U_{2}^{\perp}$

Proposition 3.1.7 Let $(V, \omega)$ be a symplectic space and $U \subset V$ a subspace. Then the following statements are equivalent:
(i) U is Lagrangian
(ii) $U=U^{\perp}$
(iii) $U$ is isotropic and $\operatorname{dim} U=\frac{1}{2} \operatorname{dim} V$.

Note that (iii) can be rephrased by saying that Lagrangian subspaces are maximal isotropic subspaces.

Lemma 3.1.8 Let $\mathcal{L}$ be a Lie algebra of Hamiltonian endomorphisms on a symplectic vector space $(V, \omega)$, and let $U, W$ be subspaces of $V$ such that $\mathcal{L}: \mathrm{U} \longmapsto \mathrm{W}$. Then,

$$
\mathcal{L}: W^{\perp} \longmapsto U^{\perp} .
$$

Proof. Let $x \in \mathcal{L}$. Since $\mathcal{L}$ consists of Hamiltonian endomorphisms $X$ is infinitesimally symplectic, and therefore satisfies

$$
\omega\left(X v_{1}, v_{2}\right)+\omega\left(v_{1}, X v_{2}\right)=0 \psi v_{1}, v_{2} \in v
$$

```
Let v & W', +},u\in
\omega(X v,u)=-\omega(v,Xu)=0
```

since $X . u \in W$ which implies $X . v \in U^{\perp}$.
Q.E.D.

Theorem 3.1.9 Let $\mathcal{L}$ be a nilpotent Lie algebra of Hamiltonian endomorphisms on a $2 n$ - dimensional symplectic vector space (V, $\omega$ ). Then there exists a chain of isotropic subspaces $V_{i}$ of dimension $i$

$$
\{0\}=v_{0} \subset v_{1}=v_{2} \subset \ldots \subset v_{n-1} \subset v_{n}
$$

where $V_{n}$ is a Lagrangian subspace, and

$$
\mathcal{L} \cdot v_{i} \subset v_{i-1}, \quad 1 \leqslant i \leqslant n
$$

Moreover, $\mathcal{L} \cdot V_{i-1}^{\perp} \subset V_{i}^{\perp}$, and so there exists a flag of the form

$$
\{0\}=v_{0} \subset v_{1} \subset \ldots \in v_{n-1}=v_{n}=v_{n}^{\perp} \subset v_{n-1}^{\perp} \subset \ldots \in v_{1}^{1} \subset v_{0}^{\perp}=v
$$

with $\mathcal{L}$ stabilizing the flag.
Proof. By Engel's theorem there exists an $\alpha_{1} \neq 0$ in $V$ such that $\mathcal{L}, \alpha_{1}=0 . \operatorname{Set} V_{1}=\underset{\sim}{R} \alpha_{1}$.

Thus $\mathcal{L}: V_{1} \longrightarrow 0 \subset V_{1}$. Since $V_{1}$ is one dimensional it is obviously isotropic with respect to $(V, \omega)$ i.e. $V_{1} \subset V_{1}^{\perp}$, and from lemma 3.1.8

$$
\mathcal{L}: v_{1}^{\perp} \longmapsto v_{1}^{\perp}
$$

Proceed by induction. Assume that $V_{i}$ is an i-dimensional isotropic subspace of $v$ with $L: v_{i} \longmapsto v_{i-1} \Leftarrow v_{i}$ then $L: v_{i}^{\perp} \longmapsto v_{i}{ }^{\perp}$

Consider the reduced space $V_{i}^{\perp} / V_{i}=W_{i}$. Now, since $L: v_{i} \longleftrightarrow v_{i}$ and $L: v_{i}^{\perp} \longmapsto V_{i}^{\perp}, \quad$ induces an endomorphism $L^{+}$on $W_{i}$. Moreover, since $\mathcal{L}$ is nilpotent, by proposition 3.1.1., $\mathcal{L}^{+}$is nilpotent on $W_{i}$.

Thus, there exists $0 \neq \beta_{i+1} \in W_{i}$ such that $L^{+} \cdot \beta_{i+1}=0$. Let $\alpha_{i+1}$ be a representative for $\beta_{i+1}$ in $v_{i}^{\perp} \subset v$. Then

$$
\begin{gathered}
\mathcal{L} \cdot \alpha_{i+1} \subset v_{i} \\
\text { Define } v_{i+1}=v_{i}+\underset{\sim}{R} \alpha_{i+1}, \text { so } \mathcal{L}: v_{i+1} \longrightarrow v_{i} \in v_{i+1} \text { and } \mathcal{L}: v_{i+1}^{\perp} \longrightarrow v_{i+1}^{\perp}
\end{gathered}
$$ by lemma. Claim $V_{i+1}$ is isotropic ie.

$$
\begin{aligned}
& v_{i}+\underset{\sim}{R} \alpha_{i+1} \subset\left(V_{i}+\underset{\sim}{R} \alpha_{i+1}\right)^{\perp} \\
& \text { But }\left(v_{i}+\underset{\sim}{R} \alpha_{i+1}\right)^{\perp}=v_{i}^{\perp} \cap\left(\underset{\sim}{R} \alpha_{i+1}\right)^{\perp} \text { by proposition } 3.1 .6 \text { (ii). } \\
& \text { By construction } \underset{\sim}{R} \alpha_{i+1} \subset v_{i}^{\perp} \text { so } v_{i} \subset\left(\underset{\sim}{R} a_{i+1}\right)^{\perp} \\
& \text { In particular, since } v_{i} \subset v_{i}^{\perp}, v_{i} \subset v_{i}^{\perp} \cap\left(\underset{\sim}{R} \alpha_{i+1}\right)^{\perp} \text { and } \\
& \text { since } \underset{\sim}{R} a_{i+1} \subset\left(\underset{\sim}{R} \alpha_{i+1}\right), \underset{\sim}{R} a_{i+1} \subset v_{i}^{\perp} \cap\left(\underset{\sim}{R} a_{i+1}\right)^{\perp}
\end{aligned}
$$

Thus,

$$
\begin{aligned}
v_{i+1}=v_{i}+\underset{\sim}{R}{\underset{a}{i+1}} & =v_{i}^{\perp} \cap\left(\underset{\sim}{R} \alpha_{i+1}\right)^{\perp} \\
& =\left(v_{i}+\underset{\sim}{R}{\underset{i}{i+1}}\right)^{\perp} \\
& =v_{i+1}^{\perp}
\end{aligned}
$$

Hence the claim is proved and $V_{i+1}$ is isotropic and $i+1$ - dimensional. By lemma 3.1 .8 since $\mathcal{L}: v_{i} \longmapsto v_{i-1}$ then $\mathcal{L}: v_{i-1}^{\perp} \longmapsto v_{i}^{\perp}$.

Proceed until $V_{n}$ is reached, which will be maximal isotropic since $\operatorname{dim} V_{n}=\frac{1}{2} \operatorname{dim} V$, and therefore by proposition 3.1.7 is a Lagrangrian subspace. Also since $v_{i} \subset v_{n} \Rightarrow v_{n}^{\perp} \subset v_{i}^{\perp}$ there exists a flag
$\{0\}=v_{0} \subset v_{1} \subset \ldots \subset v_{n-1} \subset v_{n}=v_{n}^{\perp} \subset v_{n-1}^{\perp} \subset \ldots c v_{1}^{\perp} \subset v_{0}^{\perp}=v$
with dim $V_{i}=i \quad 0 \leqslant i \leqslant n$

$$
\operatorname{dim} v_{i}^{\perp}=2 n-i \quad 0 \leqslant i \leqslant n
$$

such that $\mathcal{L}$ stabilizes the flag.
Q.E.D.

This theorem shows that there exists an ordered basis on $v,{ }^{e} q_{q_{1}}, \ldots{ }^{e} q_{n},{ }^{e} p_{1, \ldots},{ }^{e} p_{n}$ say, such that

$$
e_{p_{1}}, \ldots, e_{p_{i}} \text { spans } V_{i} \quad l \leqslant i \leqslant n
$$

and
$e_{p_{1}}, \ldots,{ }^{e} p_{n}, e_{q_{n}}, \ldots, e^{e} q_{i+1}$ spans $V_{i}^{1}, 1 \leqslant i \leqslant n-1$
In coordinates relative to this basis theorem 3.1 .9 says that each $X \in \mathcal{L}$ can be expressed by

$$
X=\left(\begin{array}{cc}
N & 0 \\
-Q & -N^{\prime}
\end{array}\right)
$$

where $Q=Q^{\prime}$ and $N$ is a strictly lower triangular matrix, $N, Q \in R^{n \times n}$. For example, consider a free system of the form

$\left(q_{1}, \ldots, q_{n}, p_{1}, \ldots, p_{n}\right)$ canonical coordinates with respect to the above basis.

It is a Hamiltonian system with Hamiltonian $H$ given by,

$$
\begin{aligned}
& H=p^{\prime} N q+\frac{1}{2} q^{\prime} Q q \\
& q=\left(q_{1}, \ldots, q_{n}\right)^{\prime}, \quad p=\left(p_{1}, \ldots, p_{n}\right)^{\prime}
\end{aligned}
$$

So

$$
\dot{x}=X_{H}(x) \quad x(0)=x_{0}, \quad x \in\left(\mathbb{R}^{2 n}, J\right)
$$

Equally, it is possible to choose an ordered basis $e_{q_{1}}, \ldots, e_{q_{n}}$, ${ }^{e_{p_{1}}}, \ldots,{ }^{e_{p_{n}}}$ on $V$ such that $e_{q_{1}}, \ldots,{ }^{e}{ }_{q_{i}}$ spans $v_{i}$ $e_{q_{1}}, \ldots, e_{q_{n}}, e_{p_{n}}, \ldots, e_{p_{i+1}}$ spans $v_{i}^{\perp}$ and express each element $X$ in coordinates relative to this basis as,

$$
x=\left(\begin{array}{cc}
N & P \\
0 & -N
\end{array}\right)
$$

where $N$ is strictly upper triangular matrix, $P=P^{\prime}, N, P \in \mathbb{R}^{n x n}$. then for a typical free system the Hamiltonian function is

$$
H=\frac{1}{2} p^{\prime} P p+p^{\prime} N q
$$

where $(q, p)^{\prime}=\left(q_{1}, \ldots, q_{n}, p_{1}, \ldots, p_{n}\right)$ are the canonical coordinates relative to the symplectic basis $e_{q_{1}}, \ldots,{ }^{e}{ }_{q_{n}},{ }^{e}{ }_{p_{1}}, \ldots,{ }^{e}{ }_{p_{n}}$.

Example 3.1.10 Suppose the state space is the symplectic vector space $\left(R^{4}, \omega\right)$ with canonical coordinates $\left(q_{1}, q_{2}, p_{1}, p_{2}\right)$ and $\omega=\sum_{i=1}^{2} d q_{i} \wedge \mathrm{dp}_{i}$. Consider the following bilinear system

$$
\dot{\mathrm{q}}_{1}=\mathrm{u} \quad \mathrm{q}_{1}(0)=0
$$

$$
\dot{q}_{2}=q_{1}+u q_{1} \quad q_{2}(0)=0
$$

$$
\begin{array}{ll}
\dot{p}_{1}=-p_{2}-u p_{2} & p_{1}(0)=0  \tag{3.1.1}\\
\dot{p}_{2}=-q_{2} & p_{2}(0)=0
\end{array}
$$

This system is obviously Hamiltonian, if

$$
\dot{x}=X_{H}+u X_{H} \quad x(0)=0
$$

$x=\left(q_{1}, q_{2}, p_{1}, p_{2}\right)^{\prime}$
then $H=p_{2} q_{1}+\frac{1}{2} q_{2}^{2}, \quad H_{u}=p_{1}+p_{2} q_{1}$
So $\quad X_{H}=q_{1} \partial / \partial q_{2}-p_{2} \partial / \partial p_{1}-q_{2}^{\partial} / \partial p_{2}$

$$
x_{H_{u}}=\partial / \partial q_{1}+q_{1}^{\partial / \partial q_{2}}-p_{2}^{\partial / \partial p_{1}}
$$

(3.1.1) can be rewritten as

Introduce new coordinates

$$
\begin{array}{ll}
\dot{q}_{0}=0 & \dot{p}_{0}=-u p_{1} \\
q_{0}(0)=1 & p_{0}(0)=1
\end{array}
$$

$$
\text { let } \bar{x}=\left(q_{0}, q_{1}, q_{2}, p_{0}, p_{1}, p_{2}\right)^{6} \text { then (3.1.1) becomes }
$$

$$
\dot{\bar{x}}=\left[\begin{array}{cccccc}
0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & -1 \\
0 & 0 & -1 & 0 & 0 & 0
\end{array}\right] \bar{x}+u\left[\begin{array}{cccccc}
0 & 0 & 0 & 0 & 0 & 0 \\
1 & 0 & 0 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & -1 & 0 \\
0 & 0 & 0 & 0 & 0 & -1 \\
0 & 0 & 0 & 0 & 0 & 0
\end{array}\right] \bar{x}
$$

A
B

The Lie algebra $\mathcal{L}=\{A, B\}_{\text {L.A. }}$ is nilpotent and the subspaces in theorem 3.1 .9 correspond to

$$
\begin{aligned}
& V_{1} \text { is spanned by }\left[\begin{array}{l}
0 \\
0 \\
0 \\
1 \\
0 \\
0
\end{array}\right]=\left(e_{p_{1}}\right) \\
& \mathrm{V}_{2} \text { is spanned by }\left[\begin{array}{l}
0 \\
0 \\
0 \\
1 \\
0 \\
0
\end{array}\right],\left[\begin{array}{l}
0 \\
0 \\
0 \\
0 \\
1 \\
0
\end{array}\right]=\left(e_{p_{1}}, e_{p_{2}}\right) \\
& \mathrm{V}_{3} \text { is spanned by }\left[\begin{array}{l}
0 \\
0 \\
0 \\
1 \\
0 \\
0
\end{array}\right],\left[\begin{array}{l}
0 \\
0 \\
0 \\
0 \\
1 \\
0
\end{array}\right],\left[\begin{array}{l}
0 \\
0 \\
0 \\
0 \\
0 \\
1
\end{array}\right] \quad=\left(e_{\mathrm{p}_{1}}, e_{\mathrm{P}_{2}}, \mathrm{e}_{\mathrm{p}_{3}}\right) \\
& \mathrm{V}_{2} \text { is spanned by }\left[\begin{array}{l}
0 \\
0 \\
0 \\
1 \\
0 \\
0
\end{array}\right],\left[\begin{array}{l}
0 \\
0 \\
0 \\
0 \\
1 \\
0
\end{array}\right],\left[\begin{array}{l}
0 \\
0 \\
0 \\
0 \\
0 \\
1
\end{array}\right],\left[\begin{array}{l}
0 \\
0 \\
1 \\
0 \\
0 \\
0
\end{array}\right] \\
& =\left(e_{p_{1}}, e_{p_{2}}, e_{p_{3}}, e_{q_{3}}\right) \\
& \mathrm{v}_{1}^{\perp} \text { is spanned by }\left[\begin{array}{l}
0 \\
0 \\
0 \\
1 \\
0 \\
0
\end{array}\right],\left[\begin{array}{l}
0 \\
0 \\
0 \\
0 \\
1 \\
0
\end{array}\right],\left[\begin{array}{l}
0 \\
0 \\
0 \\
0 \\
0 \\
1
\end{array}\right],\left[\begin{array}{l}
0 \\
0 \\
1 \\
0 \\
0 \\
0
\end{array}\right],\left[\begin{array}{l}
0 \\
1 \\
0 \\
0 \\
0 \\
0
\end{array}\right] \\
& =\left(e_{p_{1}}, e_{p_{2}}, e_{p_{3}}, e_{q_{3}}, e_{q_{2}}\right)
\end{aligned}
$$

In theorem 3.1.9 the dimension of each subspace was one greater than the previous subspace. The question now addressed is to whether a similar result is obtained when this is not necessarily so.

Theorem 3.1.11 Let $(V, \omega)$ be a symplectic vector space and

$$
\{0\}=v_{0} \subset v_{1}=\ldots \subset v_{n}=v_{n}^{\perp} \subset v_{n-1}^{\perp} \ldots \subset v_{0}^{\perp}=v
$$

a sequence of subspaces contained in $V$. Then there exists syoplectic basis of the form,

$$
\begin{aligned}
& e_{q_{1}}^{1}, \ldots,{ }^{e} q_{r_{1}},{ }^{e} q_{1}, \ldots,{ }^{e} q_{r_{2}}^{2}, \ldots,{ }^{e} q_{q_{1}^{\overline{1}}}, \ldots,{ }^{e} q_{r_{n}^{n}}^{n}
\end{aligned}
$$

such that

$$
e_{p_{1}}^{1}, \ldots,{ }_{p_{r_{i}}}^{i} \text { spans } v_{i} \quad l \leqslant i \leqslant n
$$

and

$$
e_{p_{1}^{1}}, \ldots,{ }^{e}{ }_{p_{r_{n}}^{n}},{ }_{q_{q_{n}}^{n}}, \ldots,{ }^{e} q_{1}^{n}, \ldots,{ }^{e} q_{r_{i+1}}^{i+1}, \ldots .,{ }^{e} q_{1}^{i+1} \text { spans } v_{i} \quad 1 \leqslant i \leqslant n-1
$$

Proof. By definition 3.1.5(iii), since $V_{n}$ is Lagrangian there exists a $V_{n}^{\prime}$ such that $V=V_{n} \oplus V_{n}^{\prime}$ and $V_{n}^{\prime}$ is a Lagrangian subspace. Furthermore let $v^{\prime} \subset v_{n}^{\prime}$ then the map $\left.i\left(v^{\prime}\right) \omega\right|_{V_{n}}$ defined by $\hat{\omega}, \hat{\omega}: v_{n}{ }^{\prime} \longmapsto v_{n}{ }^{*}$, is obviously an isomorphism, since ker $\hat{\omega}=\{0\}$ and dinension of $v_{n}^{\prime}$ is equal to that of $V_{n}{ }^{*}$. Thus $V$ can be rewritten as $V=V_{n} \oplus V_{n}{ }^{*}$.

If $e_{q_{1}}, \ldots, e_{q_{n}}$ is any basis for $V_{n}$, take the dual basis for $V_{n}{ }^{*}$ and let $e_{p_{1}}, \ldots, e_{p_{n}}$ be the isomorphic basis of $v_{n}^{\prime}$, then $e_{q_{1}}, \ldots, e_{q_{n}}, e_{p_{1}}, \ldots, e_{p_{n}}$ is a symplectic bases for $(V, \omega)$. In this basis $\omega$ can be represented by the matrix


Now proceed by induction. Find a basis for $V_{1}$, say $e_{\hat{p}}^{1}, \ldots,{ }_{\hat{p}}^{1}{\underset{r}{1}}_{1}^{1}$, which are of course, linear combinations of the $e_{p_{i}} s_{\text {. Orthonormalise }}$ using the Gram-Schmidt process using the Euclidean metric, and then denote the resultant basis for $V_{1}$ by $e_{p} \frac{1}{1}, \ldots,{ }^{e} p_{r_{1}}^{1}$.

Now find a vector $e_{\hat{p}}^{2} 2$ such that $e_{\hat{\mathrm{p}}}^{1} 2 \in V_{2}$ but $\mathrm{e}_{\hat{\mathrm{p}}_{1}^{2}}^{2} \& \mathrm{~V}_{1}$. Again use the Gram-Schmidt process to ensure $e_{p}^{1}, \ldots, e_{p_{r}}^{1}, e_{p}^{2}$ is Euclidean orthonormal.

Repeat this for a vector $e_{\dot{p}_{2}}^{2} \in \operatorname{span}\left\{V_{l}, e_{p_{1}^{2}}^{2}\right\}$. Continue this until an orthonormal basis for $V_{2}$ is found. The same procedure is repeated for $V_{3}, \ldots, V_{n}$.

Now let

$$
e_{q_{k}^{j}}=J e_{p_{k}} \quad 1 \leqslant j \leqslant n, \quad 1 \leqslant k \leqslant r_{i}, 1 \leqslant i \leqslant n
$$

Clearly,

$$
e_{p_{j}}^{\prime \prime} \quad J e_{p_{i}^{k}}=0
$$

since $e_{p_{j}^{i}}$ 's are linear combinations of the $e_{p_{i}}$ 's which have this property. Similarly,

$$
\begin{aligned}
& e_{q_{j}^{i}}^{\prime} J e_{q_{\ell}^{k}}=0, \text { since } \\
& e_{q} q_{j}^{i}=J e_{p_{j}^{i}}^{i} \quad \text { and } J^{\prime} J J=J
\end{aligned}
$$

Now $e_{q_{j}}{ }^{\prime} \quad J e_{p_{l}}^{k}=e_{p_{j}^{i}} \quad e_{p_{\ell}}$, these are zero unless $i=k$ and $j=\ell$ by construction.

If the Lie subalgebra $\mathcal{L}$ of $g 1(V)$ over an algebraically closed field Fis solvable then the classical theorem in the literature is Lie's theorem. Recall that a Lie algebra $\mathcal{L}$ is solvable if it satisfies $\mathcal{L}^{(1)}=\mathcal{L}$, $\mathcal{L}^{(k)}=\left[\mathcal{L}^{(k-1)}, \mathcal{L}^{(k-1)}\right], k>1$ and there exists some integer $n$ such that $\mathcal{L}^{(n)}=0$.

Proposition 3.1.12 Let $\mathcal{L}$ be a Lie algebra.
(i) If $\mathcal{L}$ is solvable then so are all sub and quotient algebras.
(ii) If $I$ is an ideal of $\mathcal{L}$ and $\mathcal{L} / I$ and $I$ are solvable then so is $\mathcal{L}$.
(iii) If $I$ and $J$ are solvable ideals of $L$ then $I+J$ is a solvable ideal.

Lemma 3.1.13 Let $V$ be a non-zero, finite dimensional vector space and a solvable Lie subalgebra over an algebraically closed field $F$ of
gl(V). Then $V$ contains a common eigenvector, $v$, for each element of $\mathcal{l}$. i.e. $X . v=\lambda(X) v, \forall X \in \mathcal{L} \quad$ where $\lambda: L \longmapsto F$.

Theorem 3.1.14 (Lie) If $\mathcal{L}$ is a solvable subalgebra of $g l(V),\{0\} \neq \mathrm{V}$ a finite dimensional vector space, then $\mathcal{L}$ stabilizes some flag in $V$.

Replace the finite dimensional vector space by a finite dimensional symplectic vector space to obtain.

Theorem 3.1.15 Let $(V, \omega)$ be a symplectic vector space of dimension $2 n$. If $\mathcal{L}$ is a solvable Lie subalgebra of Hamiltonian endomorphisms contained in $g l(V)$ over an algebraically closed field $F$, then there exists a sequence of isotropic subspaces $\mathrm{V}_{\mathrm{i}}$ of dimension i

$$
\{0\}=v_{0} \subset v_{1} \subset \ldots \in v_{n-1} \subset v_{n}
$$

with $V_{n}$ a Lagrangian subspace, and

$$
\mathcal{L}: v_{i} \longmapsto v_{i}, \quad 1 \leqslant i \leqslant n .
$$

Moreover, $L: \mathrm{v}_{\mathrm{i}}^{\perp} \longmapsto \mathrm{V}_{\mathrm{i}}^{\perp}$, and so there exists a flag of the form

$$
\{0\}=v_{0} \subset v_{1} \subset \ldots \in v_{n-1} \subset v_{n}=v_{n}^{\perp} \subset v_{n-1}^{\perp} \subset \ldots v_{1}^{\perp} \subset v_{0}^{\perp}=v
$$

and $\mathcal{L}$ stabilizes the flag.

Proof. Identical to the proof of theorem 3.1.9. Let $\alpha_{1} \neq 0$ be a common eigenvector of all elements of $\mathcal{L}$, which exists by lemma 3.1.13., and consider the subspace $V_{1}=F a_{1} . V_{1}$ is isotropic since it is onedimensional.

Again proceed by induction. Assume $V_{i}$ is an i-dimensional isotropic subspace of $V$ with $\mathcal{L}: V_{i} \longrightarrow V_{i}$ then by lemma 3.1.8. $\mathcal{L}: V_{i}^{\perp} \longmapsto V_{i}^{\perp}$ Consider the reduced space $W_{i}=V_{i}^{\perp} / V_{i}$ and the induced Lie subalgebra of endomorphisms, $\mathcal{L}^{+}$, on $W_{i}$. By proposition 3.1.12. $\mathcal{L}^{+}$is solvable and therefore there exists a common eigenvector $0 \neq \beta_{i+1} \subset W_{i}$ such that $X^{\dagger} \beta_{i+1}=\lambda\left(X^{\dagger}\right) \beta_{i+1}, \Psi X^{\dagger} \subset \sum^{\dagger}+$ Let $\alpha_{i+1}$ be a representative for $\beta_{i+1}$ in $V_{i}^{+} \in V$ then define

$$
v_{i+1}=v_{i}+F \alpha_{i+1}
$$

and then $\mathcal{L}: V_{i+1} \longmapsto V_{i+1}$ by construction.
That $V_{i+1}$ is isotropic follows exactly as in theorem 3.1.9. Again, stop the procedure when $V_{n}$ is reached. $V_{n}$ is isotropic and dim $V_{n}=\frac{1}{2}$ dim $V$ $=n$ and therefore $V_{n}$ is a Lagrangian subspace of $(V, \omega)$.

Use of proposition $3.1 .6(i)$ and lema 3.1 .8 completes the flag which is stabilized under $\mathcal{L}$.

> Q.E.D.

Essentially, this theorem says that there exists an ordered basis for ( $V, w$ ) such that each element of $\mathcal{L}$ can be expressed by a matrix of the form relative to the basis,

$$
X=\left(\begin{array}{cc}
S & 0 \\
-Q & -S^{\prime}
\end{array}\right)
$$

where $S, Q \in E^{n \times n}$
$Q=Q^{\prime}$ and $S$ is lower triangular

### 3.2 Canonical coordinates

In a globalization of theorem 3.1.11 one would assume the existence of a sequence of involutive distributions of the form

$$
\{0\}=\Delta_{0} \in \Delta_{1} \subset \ldots \in \Delta_{\mathrm{n}-1} \in \Delta_{\mathrm{n}}
$$

on a symplectic manifold ( $M, \omega$ ) with

$$
\Delta_{n}(x)=\left(\Delta_{n}(x)\right)^{\perp}, \forall x \in \quad M_{0}
$$

then find conditions for the existence of canonical coordinates $\left(q_{1}^{1}, \ldots, q_{r_{1}}^{1}, q_{1}^{2}, \ldots, q_{r_{n}}^{n}, p_{1}^{1}, \ldots, p_{r_{1}}^{1}, p_{1}^{2}, \ldots, p_{r_{n}}^{n}\right)$ on $M$ such that locally
$\Delta_{1}$ is spanned by ${ }^{\partial} / \partial p_{1} 1, \ldots,{ }^{\partial} / \partial p_{r_{1}}^{1}$
:
$\Delta_{i}$ is spanned by ${ }^{\partial} / \partial p_{1} 1, \ldots,{ }^{\partial} / \partial p_{r_{1}}^{1}, \ldots,{ }^{\partial} / \partial p_{1}^{i}, \ldots,{ }^{\partial} / \partial p_{r_{i}}^{i}$
:
$\Delta_{n}$ is spanned by ${ }^{\partial} / \partial p_{1}, \ldots,{ }^{\partial} / \partial p_{r_{1}}^{1}, \ldots,{ }^{\partial} / \partial p_{1}^{n}, \ldots,{ }^{\partial} / \partial p_{r_{n}}^{n}$

This is the goal of this section.

Definition 3.2.1 Let $(M, \omega)$ be a symplectic manifold and $i=L \longmapsto M$ be an immersion. We say $L$ is an isotropic (co-isotropic) immersed submanifold of ( $M, \omega$ ) if

$$
\left(T_{x} i\right)\left(T_{x} L\right) \subset T_{i(x)}{ }^{M} \text { is an isotropic (co-isotropic) subspace }
$$

for each $x \in L$.

The same terminology is used for submanifolds of $M$ and for subbundles of TM over submanifolds of $M$.

A submanifold $L \subset M$ is called Lagrangian if it is isotropic and there is an isotropic subbundle $\left.E \subset T M\right|_{L}$ such that $\left.T M\right|_{L}=T L E E$.

Note that i : $\mathrm{L} \longrightarrow \mathrm{M}$ is isotropic if and only if $\mathrm{i}^{*}{ }_{\omega}=0$, and also if $L=M$ is Lagrangian then $\operatorname{dim} L=\frac{1}{2} \operatorname{dim} M$ and $\left(T_{x}\right)^{-1}=T_{x}, \forall x \in L$.

Proposition 3.2.2 Let $(M, \omega)$ be a symplectic manifold and $L \subset M$ a submanifold. Then $L$ is Lagrangian if and only if $L$ is isotropic and $\operatorname{dim} L=\frac{1}{2} \operatorname{dim} M$.

A number of preparatory results are required before the proof of the main theorem in this section can be given. The first says that symplectomorphisms preserve Hamiltonian vector fields.

Theorem 3.2.3 (Jacobi) Let ( $M_{1}, \omega_{1}$ ) and ( $M_{2}, \omega_{2}$ ) be symplectic manifolds and $\psi: M_{1} \longmapsto M_{2}$ a diffeomorphism. Then $\psi$ is symplectic if and only if for all $H \in C\left(M_{2}\right), \psi_{\star}^{-1} X_{H}=X_{H_{0}} \psi^{\circ}$

The following interpretation of this theorem will be used repeatedly. Suppose $X_{f}$ and $X_{g}$ are Hamiltonian vector fields with Hamiltonian functions $f$ and $g$ respectively, on a symplectic manifold ( $M, \omega$ ) Let $(t, x) \longmapsto Y(t) x$ be the flow of $X_{f}(x)$, then since the flows of Hamiltonian vector fields are symplectic Jacobi's theorem says,

$$
\gamma(-t)_{*} X_{g}(\gamma(t) x)=X_{g(\gamma(t))}(x)
$$

The following results are standard to Symplectic Geometry and can be found for example in Abraham and Marsden [1] and Guillemin and Sternberg [22].

Theorem 3.2.4 (Darboux-Weinstein) Let N be a submanifold of the symplectic manifold $M$, and let $\omega_{0}$ and $\omega_{1}$ be two non-degenerate closed two forms on $M$ such that $\left.\omega_{0}\right|_{N}=\left.\omega_{1}\right|_{N}$. Then there exists a neighbourhood, $U$, of $N$ and a diffeomorphism $f: U \longmapsto M$ such that $\left.f\right|_{N}=$ id and $f^{*} \omega_{1}=\omega_{0}$.

This extends the classical Darboux Theorem which is the above theorem when $N$ is a point. It then states that if two symplectic forms agree on the tangent space at a point, then, up to a symplectomorphism, they agree in a neighbourhood of a point.

The case where N is a Lagrangian submanifold is of particular interest.

Corollary 3.2.5. (Kostant) Let $L$ be a Lagrangian submanifold of a symplectic manifold $(M, \omega)$. Let $L$ also be regarded as the zero section of $T^{*} L$ and let $\omega_{0}$ be the natural symplectic form on $T^{*} L$. Then there exists a neighbourhood, $U$, of $L$ in $M$ and a diffeomorplism $h$ of $U$ into $T^{*} L$ such that $\left.h\right|_{L}=$ id and $h^{*} \omega_{0}=\omega_{1}$.

Definition 3.2 .6 Let $(M, \omega)$ be a symplectic manifold. A polarization of $M$ is an involutive distribution $\Delta$ of $M$ such that $\Delta(x) \subset T_{x} M$ is a Lagrangian subspace for all $x \in M$.

The above type of polarization is usually referred to as a real polarization in the literature, it will be the kind most used in this thesis.

The fundamental example of a polarization is the foliation of a cotangent bundle by its fibres. In fact,

Proposition 3.2.7 (Kostant-Weinstein) Suppose $\Delta$ is a polarization of a symplectic manifold ( $M, \omega$ ) which is transverse to a Lagrangian submanifold L. Then there is a symplectomorphism, $f$, of some neighbourhood, $U$, of $L$ in $M$ onto some neighbourhood, $V$, of the zero section of $T^{*} L$ carrying the leaves of $\Delta$ onto the fibres of $T^{*} L$.

On $T^{*} L$ there is a naturally defined one form, $\theta_{L}$, such that $\omega_{L}=-d \theta_{L}$ is a symplectic form on $T^{*} L$. This form $\theta_{L}$ can be characterised as (i) $-d \theta_{L}=\omega_{L}$, (ii) $i(X) \theta_{L}=0$ if $X$ is tangent to the cotangent foliation and (iii) $\left.\theta_{L}\right|_{L}=0$ i.e. $\theta_{L}$ restricted to the zero section vanishes. By proposition 3.2.7. if $\Delta$ is any polarization of a symplectic manifold ( $M, \omega$ ) which is transversal to a Lagrangian submanifold $L$ then locally, about $L$, there is a uniquely defined one form $\beta$ satisfying (i) $\mathrm{d} \beta=\omega$ (ii) $i(X) \beta=0$ for $X$ lying in $\Delta$ and (iii) $\left.\beta\right|_{L}=0$.

Let $(M, \omega)$ be a $2 m$ - dimensional symplectic manifold and suppose there exists a sequence of involutive distributions on $M$ of the form
$\{0\}=\Delta_{0} \subset \Delta_{1} \subset \ldots \subset \Delta_{n-1} \subset \Delta_{n}$
such that they satisfy:
(1) $\Delta_{n}$ is a polarization of $M$ which is transversal to a Lagrangian submanifold $L$ in $M$.

This imediately implies that each $\Delta_{i}, 1 \leqslant i \leqslant n$ is isotropic in the sense that

$$
\Delta_{i}(x) \subset \Delta_{n}(x)=\left(\Delta_{n}(x)\right)^{\perp} \subset\left(\Delta_{i}(x)\right)^{\perp}, \forall x \in M
$$

(2) there exists Hamiltonian vector fields on $M$,

$$
x_{1}^{1}, \ldots, x_{r_{1}}^{1}, \ldots, x_{1}^{n}, \ldots, x_{r_{n}}^{n} \text { such that } x_{1}^{1}, \ldots, x_{r_{1}}^{1}, \ldots, x_{1}^{i}, \ldots, x_{r_{i}}^{i}
$$

$\operatorname{span} \Delta_{i}$ for $1 \leqslant i \leqslant n$.
(3) $\left[x_{j}^{i}, x_{\ell}^{k}\right] \in \Delta_{k}$ for all $k \geqslant i$
(4) there exists m Hamiltonian vector fields $X_{1}, \ldots, X_{m}$ on $M$ such that

$$
X_{1}(x), \ldots, X_{m}(x) \operatorname{spans} T_{x} L \forall x \in L
$$

$$
\text { and }\left[x_{i}, x_{l}^{k}\right] \in \Delta_{k}, \quad 1 \leqslant i \leqslant m, x_{l}^{k} \in \Delta_{k}
$$

$$
1 \leq k \leq n
$$

Theorem 3.2.8 Let $\left(M^{2} m, \omega\right)$ be a symplectic manifold in which there exists a sequence of involutive distributions satisfying (1) - (4) above. Then there exists canonical coordinates
$\left(q_{1}^{1}, \ldots, q_{r_{1}}^{1}, \ldots, q_{1}^{n}, \ldots, q_{r_{n}}^{n}, p_{1}^{1}, \ldots, p_{1}^{n}, \ldots, p_{r_{n}}^{n}\right)$ on a neighbourhood, $U^{\prime}$, of $L$ in $M$ such that ${ }^{\partial} / \partial p_{1}^{1}, \ldots,{ }^{\partial} / \partial P_{r_{1}}^{1}, \ldots,{ }^{\partial} / \partial p_{1}^{i}, \ldots,{ }^{\partial} / \partial p_{r_{i}}^{i}$ spans $\Delta_{i}$.

Proof. As the analysis performed here will be local, M can be identified with $T^{*} L$ on application of corollary 3.2.5. Fix $x \in L \subset M=T^{*} L$ and let $(U, \phi)$ be a chart in $T^{*} L, x \in U$ on which there exist canonical coordinates $\left(\AA_{1}, \ldots, q_{m}, \beta_{1}, \ldots, p_{m}\right)$. Thus

$$
\phi^{*} \omega=\sum_{l=1}^{m} d q_{i} \cdot d \beta_{i}
$$

The necessary $p$ coordinates are first constructed. Introduce the map

$$
\psi: R^{m} \longmapsto I\left(\Delta_{n}, x\right) \quad\left(I \quad\left(\Delta_{n}, x\right)\right. \text { the integral submanifold of }
$$ $\Delta_{n}$ through $x$ ) defined by

$$
\begin{aligned}
\psi:\left(p_{1}^{1}, \ldots p_{r_{1}}^{1}, \ldots, p_{1}^{n}, \ldots, p_{r_{n}}^{n}\right) & \mapsto r_{1}^{1}\left(p_{1}^{1}\right)_{0 \ldots 0} r_{r_{1}}^{1}\left(p_{r_{1}}^{1}\right)_{0 \ldots 0} \\
& r_{1}^{n}\left(p_{1}^{n}\right)_{0 \ldots 0} \gamma_{r_{n}^{n}}^{n}\left(p_{r_{n}^{n}}^{n}\right)_{x}
\end{aligned}
$$

where $\gamma_{j}^{i}$ is the flow of $X_{j}^{i}$.
Let $p_{i}=\left(p_{i}^{1}, \ldots, p_{r_{1}}^{1}, \ldots, p_{i}^{i}, \ldots, p_{r_{i}}^{i}, 0 \ldots, 0\right)$
Now using exactly the same argument as in the proof of Krener's theorem 2.1.7 along with condition (2) a neighbourhood $\bar{U} \subset U$ with $x \in \bar{U}$ is obtained such that for $1 \leqslant i \leqslant n$,
$p_{i} \longmapsto \psi\left(p_{i}\right)$ is a diffeomorphism onto $\bar{U} \cap I\left(\Delta_{i}, x\right)$, where $I\left(\Delta_{i}, x\right)$ is the integral submanifold of $\Delta_{i}$ through $x$. These are the required p coordinates.

Let $\gamma_{i}\left(q_{i}\right)$ be the flow of $X_{i}, 1 \leqslant i \leqslant m$ and define the map
$\Psi: R^{2 m} \longrightarrow \bar{U} \subset M$ by $\Psi(q, p)=\gamma(q)_{0} \psi(p)$ where
$Y(\phi)=\gamma_{1}\left(\phi_{1}\right)_{0 \ldots 0} \gamma_{m}\left(\theta_{m}\right)$
Claim. (a) $\overline{\underline{I}}$ is a diffeomorphisn of some neighbourhood of the origin in ${\underset{\sim}{2}}^{2 m}$ onto some neighbourhood $U^{\prime} \leftharpoonup \vec{U} \subset M$.
(b) $\Psi_{*}{ }^{\partial} / \partial p_{j}^{i}$ is a Hamiltonian vector on $U^{\prime}$
(c) $\Psi_{*}{ }^{\partial} / \partial p_{j}^{i}$ spans $\Delta_{i}$ on $U^{\prime}$ for $1 \leqslant j \leqslant r_{i}$
$1 \leqslant i \leqslant n$

Suppose these claims are valid for the moment and proceed to prove the theorem.

Define the coordinates $q_{j}^{i}, 1 \leqslant i \leqslant n, 1 \leqslant j \leqslant r_{i}$ by
$d q_{j}^{i}=-i\left({ }_{j}^{\partial} / \partial p_{j}^{i}\right)\left(i^{*} \omega\right)$
for then if this were true
$\Psi^{*} \omega=\sum d q_{j}^{i}, d p_{j}^{i}$.

For (3.2.1) to be valid the following must be checked.

> (i) $d q_{j}^{i}\left({ }^{\partial} / \partial p_{m}^{l}\right)=0$. That is $\left(\Psi^{*} \omega\right)\left({ }_{\partial}^{\partial} p_{m}^{l},{ }^{\partial} / \partial p_{j}^{i}\right)=0$. But $\omega\left(\Psi_{*}^{\partial} / \partial p_{m}^{l}, \Psi_{*}^{\partial} / \partial p_{j}^{i}\right)=0$ since by (c) $\Psi_{*}^{\partial} / \partial p_{j}^{i} \in \Delta_{i}$ which by assumption is isotropic for $1 \leq i \leq n$.

$$
\text { (ii) d }\left(i\left({ }^{\partial} / \partial p_{j}^{i}\right)\right)\left(\Psi^{*} W\right)=0
$$

Since then $i\left({ }^{\partial} / \partial p_{j}\right)\left(\underline{\Psi}^{*}{ }^{W}\right)$ can be expressed locally, by the derived derivative of some function.

For this Carton's identity is required ie.

$$
L_{X} \omega=\operatorname{di}(X) \omega+i(X) d \omega, X \in V(N), \omega \in \Omega^{k}(M)
$$

$s \circ d\left(i\left({ }^{\partial} / \partial_{j} \mathrm{P}_{\mathrm{j}}\right)\left(\dot{\Psi}^{*}{ }_{\omega}\right)\right)$
$=L_{a} /{ }_{\partial} p_{j}^{i} \quad \Psi^{*} \omega-i\left({ }_{\partial} \partial p_{j}^{i}\right) \quad d\left(\left(\Psi^{*} \omega\right)\right)$
but $\mathrm{d} \Psi^{*} \omega=\bar{\Psi}^{\star} \mathrm{d} \omega=0$ since $\omega$ is closed. Also,

$$
{ }^{L_{\partial}} /{ }_{\partial} p_{j}^{i} \Psi^{*} \omega=\Psi^{*} \Psi_{*}{ }^{*} /{ }_{\partial} p_{j}^{i}
$$

But from (b) $\Psi_{*}{ }^{\partial} / \partial P_{j}^{i}=X$ is a Hamiltonian vector field on $U^{\prime}$ thus $L_{X} \omega=0$ and $d\left(i\left(\partial / \partial p_{j}^{i}\right)\left(\Psi^{*} \omega\right)\right)=0$ as required.
(iii) $\bar{U}^{*} \omega=\bar{w}$, clearly $d \bar{w}=0$ and non-degeneracy follows from (a). Therefore $(3.2 .1)$ is a valid choice of $q_{j}^{i}$. This is the required coordinate system satisfying the conditions of the theorem, since they are now seen to be equivalent to (c).

Now it just remains to prove claims (a), (b) and (c).
(a) Proof. Differentiating the expression for $\mathbb{I}$ and evaluating at $\mathbb{Q}=0$, $p=0$ gives

$$
\begin{aligned}
& { }^{\partial} /\left.\partial q_{i} \Psi(\uparrow, p)\right|_{(0,0)}=x_{i}(x) \\
& { }^{\partial} /\left.\partial_{\partial} p_{j}^{i} \Psi(\uparrow, p)\right|_{(0,0)}=x_{j}^{i}(x)
\end{aligned}
$$

These vectors span $T_{X} M$ by construction, now on application of the inverse function theorem the result is obtained.
(b) Proof. Since

$$
\begin{equation*}
\Psi_{\star}{ }^{\partial} /\left.\partial p_{j}^{i}\right|_{q, p}=\gamma(q) * \frac{\partial \psi}{\partial p_{j}^{i}}(p) \tag{3.2.2}
\end{equation*}
$$

it is necessary to evaluate $\frac{\partial \psi}{\partial p_{j}^{i}}$
This is equal to

$$
\begin{equation*}
\gamma_{1}^{1}\left(p_{1}^{1}\right)_{* \ldots *} \gamma_{j-1}^{i}\left(p_{j-1}^{i}\right) * x_{j}^{i}\left(\gamma_{j}^{i}\left(p_{j}^{i}\right){ }_{o} \ldots o^{\gamma_{r_{n}}^{n}}\left(p_{r_{n}}^{n}\right), x\right) \tag{3.2,3}
\end{equation*}
$$

where $X_{j}^{i}$ is the Hamiltonian vector field on $M$ with Hamiltonian function $H_{j}^{i}$. Then by Jacobi's theorem

$$
\begin{aligned}
& \gamma_{j-1}^{i}\left(p_{j-1}^{i}\right) X_{H}^{i}\left(\gamma_{j-i}^{i}\left(-p_{j-1}^{i}\right) \cdot\right) \text { can be expressed by } \\
& X_{H_{j}^{i}}^{i} \quad \gamma_{j-1}^{i}\left(-p_{j}^{i}\right)(\cdot)
\end{aligned}
$$

Thus (3.2.3) can be written as,

$$
\mathrm{H}_{\mathrm{j}}^{\mathrm{i}} \circ \gamma_{j-1}^{i}\left(-p_{j-1}^{i}\right) \circ \ldots \gamma_{1}^{1}\left(-p_{1}^{1}\right)^{(y)}
$$

where $y=\gamma_{1}^{1}\left(p_{1}^{1}\right){ }_{0} \ldots o^{\gamma_{r_{n}}^{n}}\left(p_{r_{n}}^{n}\right) x=\psi(p)$ and (3.2.2) becomes, since $\gamma_{i}\left(\mathcal{q}_{i}\right) \quad 1 \leqslant i \leqslant m$ are flows of Hamiltonian vector fields and thus symplectic,

$$
X_{H_{j}^{i}} o^{\gamma_{j-1}^{i}}\left(-p_{j-1}^{i}\right)_{0 \ldots o^{\gamma}}^{1}\left(-p_{1}^{1}\right)_{o} \gamma_{m}\left(-\hat{q}_{m}\right)_{o \ldots o^{\gamma_{1}}\left(-\hat{q}_{1}\right)}(\Psi(\AA, p))
$$

and this can be re-expressed by considering

$$
\begin{aligned}
& H_{j_{o}}^{i} \gamma_{j-1}^{i}\left(-p_{j-1}^{i}\right) \\
&= H_{j o}^{i} \gamma_{j}^{i}\left(p_{j}^{i}\right)_{0} \ldots o_{1}\left(-q_{1}\right)(\Psi(q, p)) \\
& r_{n}^{n}\left(p_{r_{n}}^{n}\right)(x)
\end{aligned}
$$

let $k_{j}^{i}: \mathbb{R}^{2 m} \longmapsto{\underset{\sim}{R}}^{2 m}$ be defined by

$$
k_{j}^{i}\left(q_{1}, \ldots, q_{m}, p_{1}^{1}, \ldots, p_{r_{n}}^{1}\right)=\left(0, \ldots, 0, p_{j}^{i}, \ldots, p_{r_{n}}^{n}\right)
$$

therefore $(3.2 .4)$ can be written as

$$
H_{j o \Psi_{0}}^{i} k_{j o \Psi^{-1}}^{i}(\Psi(\mathbb{T}, p))
$$

which is a Hamiltonian vector field on $U^{\prime}$.
(c) Proof. Again proved by differentiation and repeated use of the Campbe11-Baker-Hausdorff formula.

$$
/_{\partial} p_{j}^{i} \Psi(Q, p)=\gamma(Q)_{*} /_{\partial} p_{j}^{i}(\psi(p))
$$

But

$$
\begin{aligned}
& { }^{\partial} / \partial_{j} p_{j}^{i} \psi(p)={ }^{\partial} / \partial_{\partial} p_{j}^{i}\left(r_{1}^{1}\left(p_{1}^{l}\right) \ldots \ldots o^{\gamma_{r}^{n}}\left(p_{r_{n}}^{n}\right) \cdot x\right) \\
& =\gamma_{1}^{1}\left(p_{1}^{1}\right)_{* \ldots *} \gamma_{j-1}^{i}\left(p_{j-1}^{i}\right) x_{k}^{i}\left(\gamma_{j-1}^{i}\left(-p_{j-1}^{i}\right)_{0 \ldots o^{r}}^{1}\left(-p_{1}^{1}\right) y\right) \\
& \text { where } y=r_{1}^{1}\left(p_{1}^{1}\right) \ldots \ldots 0^{r_{r}^{n}}\left(p_{r_{n}}^{n}\right) . x
\end{aligned}
$$

But by the Campbell-Baker-Hausdorf formula and condition (3) this is seen to be contained in $\Delta_{i}$. However, by condition (4) it is clear that on further application of the Canpbell-Baker-Hausdorff formula that $Y\left(q_{*} \Delta_{i} \subset \Delta_{i}\right.$. Further, since $\Psi$ is a diffeomorphism on $U^{\prime}$ it
follows that

$$
\partial / \partial p_{j}^{i} \Psi(\mathbb{q}, p)=\Psi_{*}^{\partial} /\left.\partial p_{j}^{i}\right|_{(q, p)}
$$

spans $\Delta_{i}\left(\mathbb{I}(q, p) \quad 1 \leqslant j \leqslant r_{i}, 1 \leqslant i \leqslant n\right.$
in a neighbourhood $U^{\prime}$ of $x \in L$ in $M$.

> Q.E.D.

A generalization of theorem 3.1.9 to symplectic manifolds would involve searching for a sequence of isotropic distributions and the existence of a Lie algebra consisting of Hamiltonian vector fields on a symplectic manifold ( $M, \omega$ ) which in some sense stabilizes the sequence.

Let $L$ be a Lie Algebra of Hamiltonian vector fields on a symplectic $\operatorname{manifold}(M, \omega)$ i.e. $\mathcal{L}=\left\{X_{H_{1}}, \ldots, X_{H_{m}}{ }^{\prime}{ }_{L}\right.$.A. Suppose there exists a sequence of involutive distributions satisfying the conditions of theorem 3.2.8 and let

$$
\begin{aligned}
& q_{1}=\left(q_{1}^{1}, \ldots, q_{r_{1}}^{1}\right), \ldots, q_{n}=\left(q_{1}^{n}, \ldots, q_{r_{n}}^{n}\right) \\
& p_{1}=\left(p_{1}^{1}, \ldots, p_{r_{1}}^{1}\right), \ldots, p_{n}=\left(p_{1}^{n}, \ldots, p_{r_{n}}^{n}\right)
\end{aligned}
$$

and also let $f\left(q_{1}^{1}, \ldots, p_{r_{n}}^{n}\right)^{\partial} /{ }_{\partial} p_{j}$ represent

$$
\begin{gathered}
f^{1}\left(q_{1}^{l}, \ldots, p_{r_{n}}^{n}\right){ }^{\partial} /{ }_{\partial} p_{1}^{j}+\ldots+f^{r_{j}}\left(q_{1}^{1}, \ldots, p_{r_{n}}^{n}\right)^{\partial} /{ }_{\partial} p_{r_{j}}^{j}, \text { and } \\
\text { e.g. } \frac{\partial H}{\partial p_{j}}=\left(\begin{array}{c}
\frac{\partial H}{\partial p_{1}} \\
\vdots \\
\frac{\partial H_{j}}{\partial p_{r_{j}}}
\end{array}\right), \text { for } 1 \leqslant j \leqslant n .
\end{gathered}
$$

```
    Then, locally
\Delta
\Delta
\vdots
i
```

This can be used to give the following generalization of theorem 3.1.9.

Theorem 3.2.9 Let $(M, \omega)$ be a symplectic manifold. Let $L$ be a Lie algebra consisting of Hamiltonian vector fields on $M$ and suppose

$$
\{0\}=\Delta_{0} \subset \Delta_{1} \subset \ldots \subset \Delta_{n-1} \subset \Delta_{n}
$$

is a sequence of involtive distribution on $M$ satisfying the conditions of theorem 3.2.8. with local coordinates as described above.

If each $X_{H} \in \mathcal{L}$ with Hamiltonian function $H$ satisfies

$$
\begin{equation*}
\left[X_{H}, \Delta_{i}\right] \subset \Delta_{i} \quad 1 \leqslant i \leqslant n \tag{3.2.5}
\end{equation*}
$$

then locally each $H$ has the form

$$
H=p_{1}^{\prime} f_{1}\left(q_{1}\right)+\ldots+p_{n}^{\prime} f_{n}\left(q_{1}, \ldots, q_{n}\right)+Q\left(q_{1}, \ldots, q_{n}\right)
$$

Proof. Locally ( $q_{1}, \ldots, q_{n}, P_{1}, \ldots, p_{n}$ ) are canonical coordinates on $M$ and each Hamiltonian vector field $X_{H}$ on $M$ can be expressed by

$$
x_{H}=\sum_{i=1}^{n} \frac{\partial H}{\partial p_{i}} \partial / \partial q_{i}-\frac{\partial H}{\partial q_{i}} \quad \partial / \partial p_{i}
$$

Consider the case $i=n$ in (3.2.5) and, without loss of generality, calculate

$$
\left[X_{H}, f \frac{\partial}{\partial}{ }^{P_{j}}\right], f:{\underset{\sim}{R}}^{2 N} \longrightarrow{\underset{\sim}{R}}^{R_{j}}, j \in\{1, \ldots, n\}
$$

an analytic map, $N=\sum_{k=1}^{n} r_{k}$.

$$
\begin{aligned}
{\left[X_{H}, f\right.} & \left.\frac{\partial}{\partial}{ }_{\partial} p_{j}\right]
\end{aligned}=\left[\sum_{i=1}^{n} \frac{\partial H}{\partial p_{i}} \frac{\partial}{\partial} q_{i}-\frac{\partial H}{\partial q_{i}} \partial / \partial p_{i}, f\left(\partial / \partial p_{j}\right]\right] .
$$

$$
\left.-f^{\prime} \frac{\partial^{2} H}{\partial p_{i} \partial p_{j}} \quad \partial / \partial q_{i}+f^{\prime} \frac{\partial^{2} H}{\partial q_{i} \partial p_{j}} \quad \partial / \partial p_{i}\right) \text { (3.2.6) }
$$

$$
(3.2 .5) \Rightarrow \quad \frac{\partial^{2} H}{\partial p_{i} \partial p_{j}}=0, \quad 1 \leqslant i, j \leqslant n
$$

$$
\Rightarrow \quad H \text { is linear in } p
$$

$$
\text { thus } H=\sum_{i=1}^{n} p_{i}^{\prime} f_{i}\left(q_{1}, \ldots, q_{n}\right)+Q\left(q_{1}, \ldots, q_{n}\right) \text {. }
$$

For the case $i=n-1$ it is obvious from (3.2.6) that (3.2.5) implies

$$
\frac{\partial^{2} H}{\partial q_{n} \partial p_{j}}=0 \quad 1 \leqslant j \leqslant n-1
$$

thus $H=\sum_{i=1}^{n-1} p_{i}^{\prime} f_{i}\left(q_{1}, \ldots, q_{n-1}\right)+p_{n}^{\prime} f_{n}\left(q_{1}, \ldots, q_{n}\right)+Q\left(q_{1}, \ldots, q_{n}\right)$

Proceed by induction. Assume $\left[X_{H}, \Delta_{k+1}\right] \subset \Delta_{k+1}$ implies
$H=\sum_{i=1}^{k+1} p_{i}^{\prime} f_{i}\left(q_{1}, \ldots, q_{k+1}\right)+\sum_{i=k+2}^{n} p_{i}^{\prime} f_{i}\left(q_{1}, \ldots, q_{i}\right)$
$+Q\left(q_{1}, \ldots, q_{n}\right)$
(3.2.7)
then consider $\left[X_{H}, \Delta_{k}\right] \subset \Delta_{k}$. By calculation (3.2.6)

$$
\frac{\partial^{2} H}{\partial q_{i} \partial p_{j}}=0, \quad 1 \leqslant j \leqslant k<i \leqslant n
$$

which by (3.2.7) implies

$$
H=\sum_{i=1}^{k} p_{i}^{\prime} f_{i}\left(q_{1}, \ldots, q_{k}\right)+\sum_{i=k+1}^{n} p_{i}^{\prime} f_{i}\left(q_{1}, \ldots, q_{i}\right)+Q\left(q_{1}, \ldots, q_{n}\right)
$$

as required.
Continue this procedure until $i=1$ in (3.2.5) which then gives the required result.
Q.E.D.

Essentially this theorem says that systems of the form

$$
\begin{array}{r}
\dot{x}=X_{H_{1}}(x)+\sum_{i=2}^{m} u_{i} X_{H_{i}}(x), \quad x \in(M, \omega) \\
x(0)=x_{0}
\end{array}
$$

with $\mathcal{L}=\left\{\mathrm{X}_{\mathrm{H}_{1}}, \ldots, \mathrm{X}_{\mathrm{H}_{\mathrm{m}}}\right\}_{\text {L.A. }}$.
and satisfying the above conditions, can be integrated by quadratures.

Example 3.2.10 Consider the following system on the symplectic manifold $\left(R^{2} \backslash\{0\}, d q, d p\right)$

$$
\dot{x}=X_{H}(x)+u X_{H}(x) \quad x(0)=x_{0}, \quad x \in{\underset{\sim}{R}}^{2}-\{0\}
$$

where $x=(q, p)^{\prime}$
where $H=\frac{1}{2} p^{2}+\frac{1}{2} q^{2}$

$$
\mathrm{H}_{\mathrm{u}}=\mathrm{q} /\left(\mathrm{p}^{2}+\mathrm{q}^{2}\right)^{\frac{1}{2}}
$$

Therefore,

$$
x_{H}=p^{\partial} / \partial q-q^{\partial} / \partial p
$$

$$
X_{H}=-\frac{p}{\left(p^{2}+q^{2}\right)^{3 / 2}}\left(q^{\partial} / \partial q+p^{\partial} / \partial{ }_{p}\right)
$$

Let $L=\left\{x \in{\underset{\sim}{R}}^{2},\{0\}: \frac{1}{2} p^{2}+\frac{1}{2} q^{2}=\right.$ constant $\}$

Obviously L is a Lagrangian submanifold with tangent space spanned by $p^{\partial} / \partial q-q^{\partial} / \partial p$.

Let $\Delta_{1}$ be spanned by $\frac{q}{p^{2}+q^{2}} \quad \frac{\partial}{\partial q}+\frac{p}{p^{2}+q^{2}} \quad \partial / \partial p \quad$ a Hamiltonian vector field on $\mathbb{R}^{2}\{0\}$ with Hamiltonian function $\tan ^{-1} p / q$. It is clear that $\Delta_{1}$ is transversal to $L$.

Changing to polar coordinates gives
$\mathrm{p}=\mathrm{r} \cos \theta$
$q=r \sin \theta$
$d q \cdot d p=(\sin \theta d r+r \cos \theta d \theta) \cdot(\cos \theta d r-r \sin \theta d \theta)$
$=r \cos ^{2} \theta d \theta, d r-r \sin ^{2} \theta d r \wedge d \theta$
$=d \theta, d\left(r^{2} / 2\right)$
Let $\phi=r^{2} / 2, \quad I=\theta$ then
$H(I, \phi)=\phi$
$H_{u}(I, \phi)=\sin I$
a symplectic change of coordinates. Thus
$\dot{I}=1$
$I(0)=I_{0}$
$\dot{\phi}=-u \cos I$
$\phi(0)=\phi_{0}$

Obvious ly integrable by quadratures and satisfies
$\left[x, \Delta_{1}\right] \Leftarrow \Delta_{1} \quad \forall x \in \mathcal{L}$
where $\mathcal{L}=\left\{X_{H}, X_{H_{u}}\right\}_{\text {L.A. }}$
Notice $\mathcal{L}$ is solvable and $S$ the ideal in $\mathcal{L}$ generated by $X_{H_{u}}$ is nilpotent i.e.
$S=\left\{\cos I^{\partial} / \partial \phi^{\prime}, \sin I^{\partial} / \partial \phi{ }_{L}, A\right.$.

In fact $S$ is abelian.

This leads into the investigation of the existence of such distributions as described in this section.

### 3.3 Existence of Sequences of Involutive Distributions

In the previous section the existence of a sequence of involutive distributions was assumed and conditions found for a particularly interesting set of canonical coordinates. Here, it is shown by the use of concepts developed by Kostant [19], Kirillov [21] and Souriau [20], how by the introduction of a nilpotent structure that the above objects arise naturally.

The main references used in this section will be Abraham and Marsden [1], Guillemin and Sternberg [22] and Wallach [25].

Definition 3.3.1 Let $(M, \omega)$ be a symplectic manifold. It is called a symplectic G - space if $G$ is a Lie group which acts on $M$ by diffeomorphisms such that if

$$
\begin{gather*}
\phi: G \times M \longrightarrow M \\
\phi(g, x)=\phi_{g}(x)=g \cdot x \\
\text { then } \phi_{g}^{*} \omega=\omega, \forall g \in G \tag{3.3.1}
\end{gather*}
$$

Furchermore ( $M, \omega$ ) is called a homogeneous symplectic G - Space if $G$ acts transitively on $M$.

Throughout the rest of this section $G$ will be assumed to be a connected Lie group unless otherwise stated.

Define
$d \phi(X)(f)(x)=\left.\frac{d}{d t}\right|_{t=0}(f(\exp -t X, x)), f \in C(M)$
which gives a map from $\mathcal{L}$ to $V(M)$, where $X \in \dot{L}$, the Lie algebra associated to G.

It can be shown that $d \phi$ is a Lie algebra homomorphism i.e.
$[\mathrm{d} \phi(\mathrm{X}), \mathrm{d} \phi(\mathrm{Y})]=\mathrm{d} \phi([\mathrm{X}, \mathrm{Y}]), \quad \forall \mathrm{X}, \mathrm{Y} \in \mathcal{L}$

By (3.3.1) $d \phi(X)$ is a locally Hamiltonian vector field for all $x \in \mathcal{L}$. However, if

$$
\mathrm{d} \phi(X) \subset \operatorname{Ham}(M, \omega)
$$

where Ham ( $M, \omega$ ) is the set of all Hamiltonian vector fields on ( $M, \omega$ ), then with this further condition $(M, \omega)$ is called a strongly symplectic G - Space.

Definition 3.3.2 Let $(M, \omega)$ be a strongly symplectic $G$ - space. A lift, $\lambda$, of $d \phi: L \longmapsto \operatorname{Ham}(M, \omega)$ is a Lie algebra homomorphism

$$
\lambda: \mathcal{L} \longmapsto \mathrm{C}(\mathrm{M})
$$

such that


Homogeneous simply connected symplectic manifolds can be completely classified, as in Wallach [25], with the aid of co-adjoint actions.

The adjoint representation is a map from the Lie group $G$ to the isomorphisms of $L$, the Lie algebra of $G$, i.e.

$$
\text { Ad }: G \longmapsto \text { Aut }(L) \text { defined by }
$$

$$
{ }_{g}^{A d} X=\left.\frac{d}{d t}\right|_{t=0} g \exp t X g^{-1}
$$

Ad is called the adjaint operator. More importantly, the coadjoint is defined by

$$
\begin{aligned}
& \qquad A d_{g}^{*}: L^{*} \longmapsto \ell^{*} \quad \forall g \in G \\
& A d_{g}^{*} f(X)=\left\langle f, A_{g^{-1}} X>, f \in \ell^{*}, x \in \mathcal{L}\right. \\
& \text { where }<,>\text { denotes the duality pairing. }
\end{aligned}
$$

Also, note that $A_{g}$ is a Lie algebra homomorphism for all $g \in G$. Let 0 be an orbit of $G$ acting on $\mathcal{L}^{*}$ i.e. $0=A d_{G}^{*}$, f $\in \mathcal{L}^{*}$. 0 has the manifold structure of $G / G_{f}$ where $G_{f}=\left\{g \in G: A d_{g}^{*} f=f\right\}$ Define an action on $O$ by $G$ as
$\sigma: G \times 0 \longrightarrow 0$
$\sigma(g, f)=A d_{g}^{*} f, \quad \forall f \in 0, \forall g \in G$.

Also note that $T_{f} 0 \approx \mathcal{L}_{f} \mathcal{L}_{f}$ where $\mathcal{L}_{f}=\{X \in \mathcal{L}:$ do $(X) \quad(f)=0\}$

Lemma 3.3.3 Define a skew form $\mathrm{B}_{\mathrm{f}}(\mathrm{X}, \mathrm{Y})=-\mathrm{f}([\mathrm{X}, \mathrm{Y}]), \mathrm{X}, \mathrm{Y} \in \mathcal{L}$, $f \in \mathcal{L}^{*}$, then $B_{f}$ is a bilinear form with kernel $\mathcal{L}_{f}$.

Thus $B_{f}$ induces a non-degenerate 2 -form $\omega(f)$ on $\mathcal{L} \mathcal{L}_{f}$ and hence gives a non-degenerate 2 - form on $\mathrm{T}_{\mathrm{f}} 0$ defined by

$$
w(f)(d \sigma(X)(f), d \sigma(Y)(f))=-f([X, Y])
$$

Hence, define $\omega_{0} \in \Omega^{2}(0)$ by $\omega_{0}\left(\left(\operatorname{Ad}^{*} \mathrm{f}\right)\right)\left(\mathrm{d} \sigma(\mathrm{X})\left(\mathrm{Ad}_{\mathrm{g}} \mathrm{f}\right)\right.$, $\mathrm{d} \sigma(\mathrm{Y})\left(\mathrm{Ad}_{\mathrm{g}}{ }_{\mathrm{F}} \mathrm{f}\right)$ $=-\operatorname{Ad}_{\mathrm{g}}^{\mathrm{f}} \mathrm{f}([\mathrm{X}, \mathrm{Y}])$, for $\mathrm{g} \in \mathrm{G}, \mathrm{f} \in 0$.

Suppose a lift exists on a symplectic G-space ( $M, W$ ) then the action $\phi: G \times M \longmapsto M$ is said to be Poisson. ( $M, \omega$ ) may then be referred to as a Poisson G-space.

Theorem 3.3.4 $\left(0, \omega_{0}\right)$ is a symplectic manifold and the action $\sigma: \mathrm{GxO} \longrightarrow 0$ is Poisson with lift given by
$\lambda(X) f \quad=f(x), \quad \forall f<0$.
Let $(M, W)$ be a Poisson $G$-space. Fix $x \in M$, then there exists a linear mapping

$$
\begin{aligned}
& \tau(x): X \longrightarrow \lambda(X)(x) \\
& \text { so } \tau(x): \mathcal{R} \longrightarrow \mathbb{R} \\
& \text { and thus } \tau: M \longrightarrow \mathcal{L}^{*}
\end{aligned}
$$

This is usually referred to as the moment map. It acts as a map which intertwines the action on $M$ with the coadjoint action on $\mathcal{L}^{*}$.

In the situation of theorem 3.3.4 $\tau: 0 \longmapsto \mathcal{L}^{*}$ is the inclusion map.

Lemma 3.3.5 If $\tau: M \longmapsto \mathcal{L}^{*}$ is the moment map of a Poisson action $\phi$ then
$\tau_{o} \phi_{g}=A d_{g}^{*} o^{\tau}$.
This is equivalent to saying that $\tau$ is $A d^{*}$ - equivariant.

Definition 3.3.6 If $\left(M_{1}, \omega_{1}\right)$ and $\left(M_{2}, \omega_{2}\right)$ are Poisson G-spaces with moment maps $\tau_{1}$ and $\tau_{2}$ respectively then a morphism is a smooth map $\mathrm{f}: \mathrm{M}_{1} \longmapsto \mathrm{M}_{2}$ such that $\mathrm{f}^{*} \omega_{2}=\omega_{1}$ and $\tau_{2}{ }^{\mathrm{o}} \mathrm{f}=\tau_{1}$.

Theorem 3.3.7 (Kostant-Souriau) Let ( $M, \omega$ ) be a homogeneous Poisson G-space with moment map $\tau$. Then $\tau(M)=0 \subset L^{*}$ is an orbit, $\tau$ is a morphism from ( $M, \omega$ ) to ( $0, \omega_{0}$ ) and a covering map.

Corollary 3.3.8 There is a bijection between simply connected homogeneous Poisson G-spaces and co-adjoint orbits.

It is not always the case that a lift exists for a given Lie Algebra L. Normally Lie algebra cohomology provides necessary conditions for existence and uniqueness. Here, however, an alternative method is considered for the specific case of interest. Further, it it always possible to extend the Lie algebra $L$ to gain a lift, as will be shown.

Suppose $T: M \longmapsto \mathcal{L}^{*}$ is a moment map on a symplectic G-space $(M, \omega)$. Let $O$ be an orbit of $G$ acting on $l^{*}$ and set

$$
N_{0}=\tau^{-1}(0)
$$

Let $x \in N_{0}$ and $f=\tau(x)$. Define

$$
v_{x}=(d \tau)^{-1}(x)\left(T_{f}\right)
$$

Assume $N_{0}$ is a smooth manifold at $x$, so $V_{x}$ is the tangent space at $x$. By definition

$$
\tau(x)(X)=\lambda(X)(x), \quad \forall x \in \mathcal{L}
$$

Let $E \subset V_{x}$ then the value of

$$
d \tau(\xi) \in \mathcal{L}^{*} \text { on } X \in \mathcal{L} \text { is } d \lambda(X)(x)(\xi)=\omega(x)\left(X_{\lambda(X)}, \xi\right)
$$

Let $\sigma: G \times 0 \longmapsto 0$ be the action on 0 defined by

$$
\sigma(g, f)=A d_{g}^{\star} f
$$

$d \sigma(X)(f)$ is the tangent vector to 0 at $f, X \in \mathcal{L}, i . e$.

$$
\mathrm{d} \sigma(\mathrm{X})(f) \in \mathrm{T}_{\mathrm{f}} 0
$$

Thus $\xi \in V_{x}$ if there exists a $Y \in \mathcal{L}$ such that

$$
\begin{equation*}
\omega(x)\left(X_{\lambda(X)}, \xi\right)=\operatorname{d\sigma }(Y)(f)(X), \quad \psi X \in \mathcal{L} . \tag{3.3.2}
\end{equation*}
$$

Now

$$
V_{x}^{\perp}=\left\{X_{\lambda(X)}(x) \in T_{x} M: \omega(x)\left(X_{\lambda(X)}(x), \xi(x)\right)=0 \quad \forall \xi \in V_{x}\right\}
$$

$\operatorname{From}(3.3 .2$.$) if X_{\lambda(X)} \in V_{x}^{\perp} \quad$ then $\operatorname{d\sigma }(Y)(f)(X)=0, \quad \forall X \in \mathcal{L}$. By definition of $\mathcal{L}_{f}=\{X \in \mathcal{L}: \operatorname{do}(X)(f)=0\}$ thus $V_{X}^{\perp}=\mathcal{L}_{f}(x)$ where $\mathcal{L}_{f}(x)$ is the corresponding subspace of the tangent space of $M$ at $x$.

Theorem 3.3.9 If $\phi: G \times M \longmapsto M$ is a nilpotent transitive action on a symplectic G-space $(M, w)$ then this action is not Poisson.

Proof. Assume the action is Poisson, so there exists a moment map $\tau: M \longmapsto \mathbb{i}^{*}$. Since the action is transitive $V_{x}=T_{x} M \forall x \in M$ which implies that $V_{x}^{\perp}=\{0\}$.

Let $Z(\mathcal{L})$ denote the centre of $\mathcal{L i . e .}$

$$
\begin{aligned}
Z(L)=\{X & \in L:[X, Y]=0 \quad v Y \in L\} \\
\text { since } d \sigma(X)(f) & =\left.\frac{d}{d t}{A d^{*}}_{\text {exp -tX }} f\right|_{t=0} \\
& =-a d_{X}^{\star} f
\end{aligned}
$$

thus $d \sigma(X)(f)(Y)=-\mathrm{ad}_{X}^{*} f(Y)=\langle f,[X, Y]\rangle$ and $\mathcal{L}_{\mathrm{f}}=\{\mathrm{X} \in \mathcal{L}: \mathrm{d} \sigma(\mathrm{X})(\mathrm{f})=0\}$. If $X \in Z(\mathcal{L})$ obviously $X \in \mathcal{L}_{f}$ as $\operatorname{do}(X)(f)(Y)=\langle f,[X, Y]\rangle=0, \forall Y \in \mathcal{L}$ Thus $\{0\}=V_{X}^{\perp}=L_{f}(x)=2(L)(x), \forall x \in M$

So $Z(\mathcal{l})=\{0\}$ which contradicts proposition 3.1.1.

Now

$$
V_{x}^{\perp}=\left\{X_{\lambda(X)}(x) \in T_{X}^{M: \omega(x)}\left(X_{\lambda(X)}(x), \xi(x)\right)=0 \quad \forall \xi \in V_{x}\right\}
$$

From (3.3.2.) if $X_{\lambda(X)} \in V_{x}^{\perp}$ then $\operatorname{do}(Y)(f)(X)=0, \forall X \in \mathcal{L}$. By definition of $\mathcal{L}_{f}=\{X \in \mathcal{L}: \operatorname{d\sigma }(X)(f)=0\}$ thus $V_{x}^{\perp}=\mathcal{L}_{f}(x)$ where $\mathcal{L}_{f}(x)$ is the corresponding subspace of the tangent space of $M$ at $x$.

Theorem 3.3.9 Lf $\phi: G \times M \longrightarrow M$ is a nilpotent transitive action on a symplectic $G$-space $(M, w)$ then this action is not Poisson.

Proof. Assume the action is Poisson, so there exists a moment map $\tau: M \longmapsto \mathcal{L}^{*}$. Since the action is transitive $V_{x}=T_{x} M \quad \boldsymbol{\psi} \times M$ which implies that $V_{x}^{\perp}=\{0\}$.

Let $Z(\mathcal{L})$ denote the centre of $\mathcal{L} i, e$.
$Z(\mathcal{L})=\{X \in L:[X, Y]=0 \quad \forall Y \in L\}$
since $d \sigma(X)(f)=\frac{d}{d t} A d^{*} \exp -\left.t X f\right|_{t=0}$
$=-a d_{X}^{*} f$
thus $\mathrm{do}(\mathrm{X})(\mathrm{f})(\mathrm{Y})=-\mathrm{ad}_{\mathrm{X}}^{*} \mathrm{f}(\mathrm{Y})=\langle\mathrm{f},[\mathrm{X}, \mathrm{Y}]\rangle$
and $\mathcal{L}_{f}=\{X \in \mathcal{L}: \operatorname{do}(X)(f)=0\}$.
If $X \in Z(\mathcal{L})$ obviously $X \in \mathcal{L}_{f}$ as $\operatorname{d\sigma }(X)(f)(Y)=\langle f,[X, Y]\rangle=0, \Downarrow Y \in \mathcal{L}$
Thus $\{0\}=V_{x}^{\perp}=L_{f}(x) \Rightarrow Z(\mathcal{L})(x), \forall x \in M$

So $Z(\mathcal{P})=\{0\}$ which contradicts proposition 3.1.1.
Q.E.D.

It is possible to avoid this problem by extending the lie algebra as follows

If $X_{1} \ldots, X_{m}$ is a basis for $\mathcal{L}$ and $d \phi\left(X_{i}\right)=X_{H_{i}}$. Let $H_{o}\left(X_{i}\right)=H_{i}$ and define $\mu_{0}\left(\sum_{i=1}^{\mathrm{m}} a_{i}, H_{i}\right)=\sum_{i=1}^{\mathrm{m}} a_{i .} \mu_{0}\left(X_{i}\right), \forall a_{i} \in \underset{\sim}{R}$.

Then $\mu_{0}: \mathcal{L} \longmapsto C(M)$ and $X_{\mu_{0}}(X)=d \phi(X)$. If $\mu_{0}$ were a Lie algebra homomorphism it would be a lift. Assume it is not a lift and let

$$
\tilde{L}=\mathcal{L} \times \mathrm{R}
$$

and define

$$
[(X, t),(Y, s)]=([X, Y], B(X, Y))
$$

where $B(X, Y)=\left\{\mu_{0}(X), \mu_{0}(Y)\right\}-\mu_{0}([X, Y])$

Then $\mathcal{L}$ is a Lie algebra. Define

$$
\lambda(x, t)(x)=\mu_{0}(x)(x)+t
$$

then $\{\lambda(X, t), \lambda(Y, s)\}=\left\{\mu_{0}(X), \mu_{0}(Y)\right\}$
$=\mu_{0}([\mathrm{X}, \mathrm{Y}])+\beta(\mathrm{X}, \mathrm{Y})$
$=\lambda([(X, t),(Y, s)])$
thus $\lambda$ is a Lie algebra homomorphism.

Let $\tilde{G}$ be the connected simply connected Lie group with Lie algebra $\tilde{\mathcal{L}}$. Let $H$ be the connected subgroup with Lie algebra $\{(0, t): t \in \underset{\sim}{R}\}$. then $\tilde{G} / \mathrm{H}$ is the simply connected Lie group with Lie algebral. Let

$$
v: \tilde{G} / H \longmapsto G
$$

be the covering map.
Let $\mu: \tilde{G} \longmapsto G$ be defined by $\mu(\tilde{g})=v(\tilde{g} H)$. Then $\tilde{G}$ acts on $M$ by $\tilde{g} . x=\mu(\tilde{g}), x, x \in M$ and $\lambda: \tilde{L} \longmapsto C(M)$ is a lift. In this case, see Wallach [25] for example, there exists the following.

Theorem 3.3.10 If ( $M, \omega$ ) is a simply connected homogeneous G-space, then there exists a $\beta$ of the form (3.3.3) and $f \in \tilde{l}^{*}$ so that the moment map $\tau: M \mapsto \tilde{L}^{*}$ is a morphism from $(M, \omega)$ to $\left(0, \omega_{0}\right)$ and a covering map, where $0=\operatorname{Ad}_{\tilde{G}}^{*} \tilde{f}$.

Essentially theorems 3.3 .7 and 3.3 .10 allow all work to be done on the coadjoint orbit and then by use of the moment map restated on the symplectic manifold ( $M, \omega$ ).

The nilpotent case is now considered in depth using the work of Kirillov [21].

Definition 3.3.11 Let $\mathcal{L}$ be a Lie algebra and let $f \in \mathcal{L}^{*}$. Then a polarizing subalgebra of $\mathcal{L}$ for $f$, sometimes called a polarization of $f$, is a subalgebra $h$ of $\mathcal{L}$ satisfying,
(i) $f([h f])=0$
(ii) If $V \subset \mathcal{L}$ is a subspace so that $f([X, Y])=0$ for $X, Y \in V$ and $v \Rightarrow h$ then $v=h$.

In general there does not exist a polarization for $f$ for all Lie algebras. For nilpotent Lie algebras, however, there exists polarizations for arbitrary $f \in \mathcal{L}^{*}$.

Let $V$ be a finite dimensional real vector space. Let $B: V \times V \longmapsto \underset{\sim}{R}$ be bilinear and skew-symmetric, although not necessarily non-degenerate. If $W$ is a subspace of $V$, set

$$
\mathcal{W}^{\perp}=\{v \in V: B(v, w)=0, \forall w \in W\}
$$

Then
$\operatorname{dim} W+\operatorname{dim} W^{\perp}=\operatorname{dim} V+\operatorname{dim}\left(V^{\perp} \cap W\right)$ since for $v c V$, set $\tilde{v}(u)=B(v, u)$, let $\zeta(v)=\left.\tilde{v}\right|_{W}$ then $\zeta: v \longmapsto\left(W / W \cap v^{\perp}\right)^{*}$, $\operatorname{ker} \zeta=W^{\perp} . \quad$ Thus, $\operatorname{dim} V=\operatorname{dim} W-\operatorname{dim}\left(W \cap V^{\perp}\right)+\operatorname{dim} W^{\perp}$

A subspace $W$ of $V$ is said to be isotropic if

$$
\left.B\right|_{W \times W}=0 \quad \text { i.e. } W \subset W^{\perp}
$$

If $W$ is maximal as an isotropic subspace then $V^{\perp} \subset W$, hence $2 \operatorname{dim} W \leq \operatorname{dim} V+\operatorname{dim} V^{\perp}$. More precisely, see Wallach [25],

Lemma 3.3.12 If $W \subset V$ is an isotropic subspace, then $W$ is maximal isotropic if and only if

$$
\operatorname{dim} W=\frac{1}{2}\left(\operatorname{dim} V+\operatorname{dim} V^{\perp}\right)
$$

Lemma 3.3.13 Let $V^{\prime} \subset \quad \mathrm{V}$ be a subspace of codimension one. Let $N^{\prime}=\left\{v \in V^{\prime}: B\left(v, v^{\prime}\right)=0 \quad \forall v^{\prime} \in V^{\prime}\right\}$. If $V^{\perp} \in V^{\prime}$ then $V^{\perp}$ is of codimension one in $N^{\prime}$. If $W \subset V^{\prime}$ is maximal isotropic in $V^{\prime}$, then $W$ is maximal isotropic in $V$.

Proof. If $V^{\perp} \subset V^{\prime}$ then since

$$
v^{\perp}=\{v \in V: B(v, u)=0, \forall u \in V\}
$$

it can rewritten as $V^{\perp}=\left\{v \subset V^{\prime}: B(v, u)=0, \psi u \in V\right\} \subset N^{\prime}(3.3 .4)$ Let $x \in V$ be so that $x \notin V^{\prime}$. Then it is claimed $V^{\perp}=\mathbb{R}$ where

$$
R=\left\{v \in N^{\prime}: B(x, v)=0\right\}
$$

Let $v_{1} \in v^{1}$ then $v_{1} \in N^{\prime}$ so $B\left(v_{1}, v^{\prime}\right)=0$ from (3.3.4) and $B\left(v_{1}, x\right)=0$ so $V^{\perp} \subset \mathcal{R}$. If $v \in Q, B(x, v)=0, V=V^{\prime}+R x$. Also $v \in N^{\prime}$ so $B\left(v, v^{\prime}\right)=0 \Rightarrow B\left(v, R x+v^{\prime}\right)=0 \Rightarrow B(v, v)=0 \Rightarrow v \in v^{\perp}$.

So claim is proved.
So dim $\left(N^{\prime} / V^{\perp}\right) \leqslant 1$. But since $N^{\prime}=V^{\perp}$ with respect to the vector space $V^{\prime}$ there exists a non-degenerate 2-form on $V^{\prime} / V^{\prime}=V^{\prime} / N^{\prime}$ and so the $\operatorname{dim}\left(V^{\prime} / N^{\prime}\right)$ is even. Also, $\operatorname{dim}\left(V / V^{+}\right)$is even. Thus
$\operatorname{dim} V / V^{\perp}-\operatorname{dim} V^{\prime} / N^{\prime}$
$=\operatorname{dim} V-\operatorname{dim} V^{\perp}-\operatorname{dim} V^{\prime}+\operatorname{dim} N^{\prime}$
$=\operatorname{dim} V-\operatorname{dim} V^{\prime}+\operatorname{dim} N^{\prime}-\operatorname{dim} V^{\perp}$
$=1+\operatorname{dim} N^{\prime} / v^{\perp}$
is even, and since $\operatorname{dim}\left(N^{\prime} / V^{\perp}\right) \leqslant 1$ it follows that $\operatorname{dim}\left(N^{\top} / V^{\perp}\right)=1$.

By lemma 3.3.12, $W$ is maximal if and only if $\operatorname{dim} W=\frac{1}{2}\left(\operatorname{dim} V+\operatorname{dim} V^{\perp}\right)$ thus, $\frac{1}{2}\left(\operatorname{dim} V+\operatorname{dim} V^{\perp}\right)=\frac{1}{2}\left((\operatorname{dim} V-1)+\left(\operatorname{dim} V^{\perp}+1\right)\right)=\frac{1}{2}\left(\operatorname{dim} V^{\prime}+\operatorname{dim} N^{\prime}\right)$ $=\frac{1}{2}\left(\operatorname{dim} V^{\prime}+\operatorname{dim} V^{\prime}\right)$

Hence the result.

Lemma 3.3.14 Let $\mathcal{L}$ be a nilpotent Lie algebra. Suppose that the centre of $\mathcal{L}, Z(\mathcal{L})$, is one dimensional. That is,

$$
z(\mathcal{L})=\{x \in L:[x, L]=0\}
$$

and $\operatorname{dim}(Z(L))=1$.
Suppose that $\mathrm{R} Z=Z(\mathcal{O})$. Then there exists $X, Y \in \mathcal{L}$ such that $[\mathrm{X}, \mathrm{Y}]=\mathrm{Z}$ and if

$$
s=\{W \in \mathcal{L}:[Y, W]=0\}
$$

then $S$ is a codimension one subalgebra of $\mathcal{L}$.

Proof. By Engel's theorem 3.1.2 and by considering the adjoint representation of $\mathcal{L}$ i.e. $a_{\mathcal{L}}: \mathcal{L} \longmapsto \mathcal{L}, \mathcal{L}$ can be decomposed as,

$$
\mathcal{L}=L_{n} \supset \mathcal{L}_{n-1} \supset \ldots \supset \mathcal{L}_{1} \supset\{0\}
$$

a sequence of ideals with $\left[\mathcal{L}, \mathcal{L}_{\mathrm{i}}\right] \subset \mathcal{L}_{\mathrm{i}-1}$ and $\operatorname{dim} \mathcal{L}_{\mathrm{i}}=\mathrm{i}$.

$$
\text { Since } \operatorname{dim} Z(L)=1 \Rightarrow z(\mathcal{L})=\mathcal{L}_{1} \text {. }
$$

Let $Y \in \mathcal{L}_{2}$ be such that

$$
\mathcal{L}_{2}=\underset{\sim}{R} Y+\underset{\sim}{R} Z
$$

If we $\mathcal{L}$, then
$\operatorname{ad}_{W} Y=[W, Y]=f(W) Z, f: \mathcal{L} \rightarrow \underset{\sim}{R}$ linear, Claim $f \neq 0$, since
if $f=0[W, \underset{R}{R} Y+\underset{R}{R} Z]=0 \quad \forall \in \mathcal{R} Y+\underset{\sim}{R} Z \in Z(\mathcal{L})$ which
is a contradiction, since dim $Z(\mathcal{L})=1$.
Let $X \in \mathcal{L}$ be such that $f(X)=1$. Then
$S=\{W \subset \mathcal{L}:[W, Y]=0\}$

Then $S=\operatorname{ker} \mathrm{f}$
Hence $\operatorname{dim} \mathcal{L}=\operatorname{dim} S+1$
Q.E.D.

Theorem 3.3.15 [Kirillov [21]] If $\mathcal{L}$ is a nilpotent Lie algebra and $f \in \ell^{*}$, then there exists a polarization of $f$.

Proof. The result is proved by induction on the dimension of $\mathcal{L}$. If $\operatorname{dim} \mathcal{L}=1$, then $\mathcal{L}$ is a polarization of $f$ is immediate from the definition.

Suppose the theorem has been proved for all nilpotent Lie algebras of dimension $\leqslant n-1$ and $\operatorname{dim} \mathcal{L}=n$. Let $Z(\mathcal{L}$ be the centre of $\mathcal{L}$, choose $f$, a real linear functional on $\mathcal{L}$ such that $Z(\mathcal{L}) \not \neq$ er $f$, and denote $z_{0}(\mathcal{L})$ the subspace of $Z\left(\mathcal{L}\right.$ consisting of all $z_{o}$ such that $f\left(z_{0}\right)=0$. Clearly $Z_{0}(\mathcal{L}) \subset$ kerf $f$ so,
$\operatorname{dim} z_{o}(\mathcal{L}) \leqslant \operatorname{dim}(\operatorname{ker} f)$.

Consider $\left.f\right|_{Z(\mathbb{L})}$ then

$$
\operatorname{dim} Z(L)=1+\operatorname{dim}\left(\left.\operatorname{ker} f\right|_{Z(U)}\right)
$$

$\geqslant 1+\operatorname{dim} Z_{o}(\mathcal{l})$
$\Rightarrow \operatorname{dim} z_{0}(l) \leq \operatorname{dim} z(\mathcal{L})-1$.

Case 1: $\operatorname{dim} Z_{0}(\mathbb{D}) \geqslant 1 \Rightarrow \operatorname{dim} Z(\mathcal{L}) \geqslant 2$. So there exists Y $\in Z(\mathcal{L})$ such that $f(Y)=0, Y \notin 0$. Then $f$ defines an element of $\left(\mathcal{L}_{\text {RY }}\right)$, $\bar{f}$. Let $\bar{h}$ be a polarization of $\bar{f}$. Obviously it is possible to choose RY $Y \notin, h \in h$ and if $h \mapsto \not \subset$ under the canonical map, then it can be shown that $h$ is a polarization of $f$ as follows.

$$
\text { Set } B_{f}\left(x_{1}, x_{2}\right)=-f\left(\left[x_{1}, x_{2}\right]\right), \forall x_{1}, x_{2} \in \mathcal{L} \text {. }
$$

Clearly, $h$ is isotropic for $B_{f}$, since $Y<2(D)$. If $h \subset V$ with $V$ an isotropic subspace of $\mathcal{L}$ relative to $B_{f}$, then $V \mapsto \bar{V} \supset \bar{h}$ under the canonical map and $\overline{\mathrm{V}}$ is isotropic for $\mathrm{B}_{\mathrm{F}}$. Hence by definition $\overline{\mathrm{V}}=\overline{\mathrm{h}}$ But $\underset{\sim}{R} Y \subset R$ hence $V=h$.

Case 2 : $\operatorname{dim} Z_{0}(\mathcal{L})=0 \Rightarrow \operatorname{dim} Z(\mathcal{L}) \geqslant 1$. The only case not yet covered is $\operatorname{dim} Z(\mathcal{L})=1$.

Let $Z(\mathcal{\ell})=\underset{\sim}{R} Z, f(Z)=1$. Let $X, Y \in \mathcal{L}$ be such that $[X, Y]=Z$ and
$S=\{W \in \mathcal{L}:[Y, W]=0\}$
as in lemma 3.3.14.

It has already been seen that $\operatorname{dim} S=n-1$. Let $K \subset S$ be a polarization for $\left.f\right|_{S}$. Again, set $B_{f}\left(X_{1}, X_{2}\right)=-f\left(\left[X_{1}, X_{2}\right]\right)$. than it is claimed, $\mathcal{L}^{1}=\left\{W \in \mathcal{L}: B_{f}(W, \mathcal{L})=0\right\} \subset S$. Since $\mathcal{L}=\underset{\sim}{R} X \oplus S$ and if $W=c X+s, c \in R, s \in S$, and $W \in \mathcal{L}^{\mathcal{L}}$ then $B_{f}(c X+s, Y)=c B_{f}(X, Y)+B_{f}(s, Y)$
$=-c f([X, Y])+f([Y, s])$
$=-c \mathbf{f}(Z)+\mathbf{f}(0)$
$=-c \quad, Y \in \mathcal{L}$
Hence $c=0$. Thus $W \in S$ which implies $\mathcal{L}^{\perp} \subset S$. Now lemma 3.3.13 implies that $h$ is a polarization for $f$ on $\mathcal{L}$.

Corollary 3.3 .15 If $h$ is a polarization with respect to f then $A d \underset{g}{ }$ is a polarization with respect to $\mathrm{Ad}_{\mathrm{g}}^{\mathrm{*}}$.
 homomorphism and $h$ is a polarization of $f$. So it is sufficient to show that $A d{ }_{g} h$ is maximal isotropic.
 f which implies that

$$
\begin{align*}
& {A d_{g}-1}^{W} \in K \text { or } W \in A_{g} h \\
\Rightarrow \quad & A d_{g} h \text { is maximal isotropic. }
\end{align*}
$$

Obviously $\mathcal{L}_{f}=\{X \in \mathscr{L}:\langle f,[X, Y]\rangle=0$ Y $X \in \mathcal{L}\}$ is a subalgebra contained in $L$ since $\langle f,[X, Y]\rangle=0 \quad \forall X, Y \in \mathcal{L}_{\mathrm{E}}$.

$$
\begin{aligned}
& \text { Let } \\
& \sigma: G \times O \longmapsto 0
\end{aligned}
$$

as before with $0=A d_{G}^{*} f$, for sone fixed $f \in \mathcal{L}^{*}$.
The following lemma can be found in Kirillov [21].

Lemma 3.3.16 If $h$ is a polarization at $f \in \mathcal{L}^{*}$ in the nilpotent Lie algebra $L$ and $O$ is the orbit in $L^{*}$ containing $f$, then

$$
\begin{equation*}
\operatorname{dim} h=\operatorname{dim} l-\frac{1}{2} \operatorname{dim} 0 \tag{3.3.5}
\end{equation*}
$$

This leman along with theorem 3.3.14 says there exists a Lagrangian subspace of $\mathrm{T}_{\mathrm{f}} \mathrm{O}, \mathrm{f} \in \mathrm{O}$ since (3.3.5) implies

$$
\begin{aligned}
\operatorname{dim} h / \mathcal{L}_{f} & =\operatorname{dim} \mathcal{L}-\frac{1}{2}\left(\operatorname{dim} \mathcal{L}-\operatorname{dim} L_{\mathrm{f}}\right)-\operatorname{dim} \mathcal{L}_{\mathrm{f}} \\
& =\frac{1}{2} \operatorname{dim} \mathcal{L}-\frac{1}{2} \operatorname{dim} L_{f}=\frac{1}{2} \operatorname{dim} L / L_{\mathrm{f}}
\end{aligned}
$$

So this subspace is given by

$$
\mathrm{d} \sigma(\mathcal{W})(f) \subset \mathrm{T}_{\mathrm{f}} \mathrm{o}^{0} .
$$

Note $h$ is a nilpotent Lie subalgebra of $\mathcal{Q}$. Also it is known that

$$
v / L_{f} \cap v=\left(L_{f}+v\right) / L_{f} \quad v \in L_{1} \quad \text { therefore it is possible }
$$

to construct the following sequences of quotient algebras,

$$
\begin{equation*}
\{0\}=l_{f} / L_{f} \subset L_{f}+h^{m} / L_{f} \subset \ldots \subset L_{f}+h^{j} / L_{E} \tag{3.3.6}
\end{equation*}
$$

where $h^{1}=h, h^{i}=\left[h, h^{i-1}\right], h^{n+1}=0$.
$\{0\}=l_{f} / \mathcal{L}_{f} \subset h \cap l^{n}+L_{f} / L_{f}=\ldots \subset h \cap l^{1}+L_{f} / l_{f}$
where $L^{1}=\mathcal{L}, \mathcal{L}^{i}=\left[\mathcal{L}, \mathcal{L}^{i-1}\right], \mathcal{L}^{n+1}=0$.

It is worthwhile observing that the proceeding theory is valid whether sequence (3.3.6) or (3.3.7) is used. However, the calculations are performed for (3.3.7) as they will appear in chapter 4.2 .

It is now possible to construct a sequence of isotropic distributions as in theorem 3.2.8. But first some notation.

Let $\Delta$ be a distribution on ( $0, \omega_{0}$ ) then it is said to be G-invariant if

$$
\mathrm{Ad}_{\mathrm{g}}^{*} \quad \Delta \subset \Delta, \quad \forall \mathrm{~g} \in \mathrm{G}
$$



Lemma 3.3.17 On the Poisson $G-s p a c e\left(0, \omega_{0}\right)$ there exists a sequence of G-invariant isotropic involutive distributions, if $G$ is nilpocent, of the form

$$
\{0\}=\Delta_{0} \subset \Delta_{1} \subset \ldots \subset \Delta_{n-1} \subset \Delta_{n} \text { with } \Delta_{n} \text { Lagrangian. }
$$

Proof. Define $\Delta_{i}, 0 \leqslant i \leqslant n$ by

$$
\Delta_{i}\left(\operatorname{Ad}_{g}^{*} f\right)=\operatorname{da}\left(\operatorname{Ad}_{g}\left(L_{f} \cap L^{n-i+1}+L_{f}\right)\right)\left(\operatorname{Ad}_{g}^{*} f\right)
$$

note er $d \sigma\left(\operatorname{Ad}_{g}^{*} f\right)=\operatorname{Ad}_{g} \mathcal{L}_{\mathrm{E}}, \Psi g \in G, f \in 0$, thus $\Delta_{i}\left(\operatorname{Ad}_{g}^{*} f\right)$ is isomorphic to $\mathrm{Ad}_{\mathrm{g}}\left(\mathcal{L}_{\mathrm{f}} \cap \ell^{n-i+1}+\mathcal{L}_{f}\right) /_{\mathrm{Ad}_{g}} \mathcal{L}_{\mathrm{f}}$ then from (3.3.7)

$$
\Delta_{i}\left(A d_{g}^{\star} f\right) \subset \Delta_{i+1}\left(A d_{g}^{\star} f\right), \forall g \in G, f \in 0
$$

Which defines the sequence of distributions

$$
\{0\}=\Delta_{0} \subset \Delta_{1} \subset \ldots \subset \Delta_{\mathrm{n}-1} \subset \Delta_{\mathrm{n}}
$$

Since $L$ is nilpotent $\left[L^{i}, L^{j}\right] \subset \mathcal{L}^{i+j}$ by induction i.e. $\left[\mathcal{L}, L^{j}\right]=L^{j+1}$ by definition for all $j$, suppose $\left[L^{i-1}, L^{j}\right] \subset L^{i+j-1}$ then $\left[\mathcal{L}^{i}, L^{j}\right]$ can be rewritten using the Jacobi identity as

$$
\begin{aligned}
{\left[\left[L, L^{i-1}\right], \mathcal{L}^{j}\right] } & =\left[\left[L, L^{j}\right], \mathcal{L}^{i-1}\right]+\left[L,\left[\mathcal{L}^{i-1}, L^{j}\right]\right] \\
& \subset L^{i+j}+\left[L, L^{i+j-1}\right] \\
& =L^{i+j}
\end{aligned}
$$

Also do and $A{ }_{g}$ are Lie algebra homomorphisms so $\Delta_{i}$ is involutive, for let $X, Y \in \mathcal{L}_{f} \cap \mathcal{L}^{\mathrm{n}-\mathrm{i}+1}+\mathcal{L}_{\mathrm{f}}$ then

$$
\begin{aligned}
& =\Delta_{i}\left(\operatorname{Ad}_{g}{ }^{\mathrm{f}}\right) .
\end{aligned}
$$

For G-invariance it is necessary to prove that

$$
\operatorname{Ad}_{\mathrm{g}}^{\star} \Delta_{i} \subset \Delta_{i}, g \in G, 1 \leqslant i \leqslant n
$$

i.e. $\operatorname{Ad}_{g_{1} *}^{*} \Delta_{i}\left(\operatorname{Ad}_{g_{2}}^{*} f\right) \subset \Delta_{i}\left(\operatorname{Ad}_{g_{1} g_{2}}{ }^{\mathrm{f}}\right), \quad 1 \leqslant i \leqslant n$ Let $y \in L, X \in L_{f} \cap L^{n-i+1}+\mathcal{L}_{f}, g_{1}, g_{2} \in G$
$\operatorname{Ad}_{\mathrm{g}_{1} *}^{*} \operatorname{do}\left(\mathrm{Ad}_{\mathrm{g}_{2}} \mathrm{X}\right)\left(\mathrm{Ad}_{\mathrm{g}_{2}}^{\mathrm{f}}\right)(\mathrm{Y})$
$=\operatorname{Ad}_{\mathrm{g}_{1}}^{*} \quad \mathrm{do}\left(\mathrm{Ad}_{\mathrm{g}_{2}}^{\mathrm{X}}\right) \quad\left(\mathrm{Ad}_{\mathrm{g}_{2}}^{\mathrm{f}} \mathrm{f}\right)(\mathrm{Y}) \quad\left(\mathrm{Ad}_{\mathrm{g}_{1}}^{*}\right.$ is linear)

$=\left\langle\mathrm{Ad}_{\mathrm{Ad}_{\mathrm{g}_{2}}} \mathrm{X} \quad \mathrm{Ad}_{\mathrm{g}_{2}}^{\mathrm{F}}, \quad \mathrm{Ad}_{\mathrm{B}_{1}{ }^{1}} \mathrm{Y}\right\rangle$
$=-\left\langle A d_{\mathrm{g}_{2}}^{*} \mathrm{f},\left[\operatorname{Ad}_{\mathrm{g}_{2}} \mathrm{X}, \mathrm{Ad}_{\mathrm{g}_{1} 1} \mathrm{Y}\right]\right\rangle$
$=-\left\langle\mathrm{f}, \mathrm{Ad}_{\mathrm{g}_{2}^{-1}}\left[\operatorname{Ad}_{\mathrm{g}_{2}} \mathrm{X}, \mathrm{Ad}_{\mathrm{ra}_{1}^{-1}} \mathrm{Y}\right]\right\rangle$
$=-\left\langle f, A d_{g_{2}^{I}}{ }^{A d} g_{1}^{-1}\left[A d_{g_{1}}{ }_{g_{2}} X, Y\right]\right\rangle$
$=-\left\langle A d_{g_{1} g_{2}}{ }^{\mathrm{f}}, \quad\left[\operatorname{Ad}_{\mathrm{g}_{1} g_{2}} \mathrm{X}, \mathrm{Y}\right]>\right.$
$=\mathrm{d} \sigma\left(\operatorname{Ad}_{\mathrm{g}_{1} \mathrm{~g}_{2}} \mathrm{X}\right)\left(\mathrm{Ad}_{\mathrm{g}_{1} \mathrm{~g}_{2}}^{\star} \mathrm{f}\right)(\mathrm{Y})$
therefore since $Y$ was arbitrary

$$
\operatorname{Ad}_{g_{1}}^{*} * \Delta_{i}\left(\operatorname{Ad}_{g_{2}}^{*} \mathrm{f}\right) \subset \Delta_{i}\left(\operatorname{Ad}_{g_{1} g_{2}}^{*} \mathrm{f}\right)
$$

as required.
The fact that each $\Delta_{i}, 0 \leq i \leq n-1$ is isotropic and $\Delta_{n}$ is Lagrangian follows from theorem 3.3.14, corollary 3.3.15 and lema 3.3.16. Q.E.D.

Observe that $\Delta_{n}$ is a real polarization of $\left(0, \omega_{0}\right)$ in the sense of definition 3.2.6.

For the case in question, an orbit 0 in $\hat{\ell}^{*}$ given by $0=A d_{G}^{*} f$, for some $f \in \mathcal{L}^{\star}$, the invariant two form, $d f$, is given by
$d f(X, Y)=-\langle f,[X, Y]\rangle$,
see chapter 5.3 of Guillemin and Sternberg [22], and the invariant polarization on the orbit through $f$ corresponds to the subalgebra constructed above.

> As the two-form $\omega_{0}$ on 0 is invariant under $\sigma_{g}, \forall g \in G$, i.e. $\sigma_{g}^{*} \omega_{0}=\omega_{0}$, df is equal to $\omega_{0}$.

The next proposition due to Kostant, see Guillemin and Sternberg [22], is of particular importance.

Proposition 3.3.18 Let $f$ be a one-form on $M$, a symplectic manifold such that its symplectic form is given by $\mathrm{df}=\omega$ and such that the set $\left\{x:\left.f\right|_{x}=0\right\}$ is a $n$-dimensional submanifold, L. Then $L$ is Lagrangian and there is a unique polarization defined on $M$ in a neighbourhood of $L$ which is transversal to $L$ and whose associated one - form is $f$.

Lemma 3.3.19 The sequence of distributions defined on a homogeneous strongly symplectic $G$-space ( $0, \omega_{0}$ ) in lemma 3.3 .17 satisfy the conditions (1) - (3) of theorem 3.2.8.

Proof. Take the Lagrangian submanifold $L$ to be that constructed in proposition 3.3.18 and the sequence of distributions are then that of lemma 3.3.17. By assumption $d \sigma(X) \subset$ Ham ( $0, \omega_{0}$ ) - Hamiltonian vector
fields on 0 - and thus by (3.3.7) and transitivity there exists Hamiltonian vector fields.

$$
\begin{aligned}
& d \sigma\left(A d_{g} Y_{1}^{i}\right)\left(A d{ }_{g}^{*} f\right), \ldots, d \sigma\left(A d{ }_{g} Y_{r_{i}}^{i}\right)\left(A d_{g}^{*} f\right)
\end{aligned}
$$

which spans $\Delta_{i}\left(\underset{g}{\left(d^{\star} f\right.}\right), \forall g \in G, 1 \leqslant i \leqslant n$
and so condition (2) is satisfied. Condition (3) follows immediately from the facts that do and $\mathrm{Ad}_{\mathrm{g}}$ are Lie algebra homomorphisms and $\left[L^{i}, L^{j}\right] \subset L^{i+j} \subset R^{j}, j \geqslant i$
Q.E.D.

Before investigating condition (4) some results from Kostant [19] are required.

Definition 3.3.20 If $M$ and $N$ are manifolds and $\psi: M \mapsto N$ is an analytic map, $X \in V(M), Y \in V(N) X$ is said to be $\psi$ - related to $Y$ if

$$
\psi_{\star} X(x)=Y(\psi(x)), \forall x \in M .
$$

Lemma 3.3.21 If $\psi: M \mapsto N$ and $X_{i}$ is $\psi$ - related to $Y_{i}, i=1,2$ then $\left[X_{1}, X_{2}\right]$ is $\psi$ - related to $\left[Y_{1}, Y_{2}\right]$.

Proposition 3.3.22 (Kostant [19]). Let ( $M, \omega$ ) be a simply connected homogeneous symplectic G - space with a Poisson action

$$
\phi: G \times M \mapsto M
$$

Fix $x_{0} \in M$, let $\tau: M \longrightarrow 0 \subset \ell^{*}$ be the moment map with $0=A d_{G}^{*} \tau\left(x_{0}\right)$ and Poisson action on 0 given by
$\sigma: G \times 0 \longmapsto 0$
$\sigma(g, f)=A d_{g}^{*} f$.
Then $d \phi(X)$ is $\tau$ - related to do (X).

Proof. Need to show that $\tau_{k} \mathrm{~d} \phi(\mathrm{X})(\mathrm{x})=\mathrm{d} \sigma(\mathrm{X}) \tau(\mathrm{x})$. By definition da (X) (f) $(\psi)=\left.\frac{d}{d t}\right|_{t=0} ^{\psi}(\exp -t X . f) \psi \in \quad C\left(\mathcal{L}^{*}\right)$
Let $Y \in \mathcal{L}$ and define $\psi^{Y}: L^{*} \longmapsto R$ by $\psi^{Y}(f)=\langle E, Y\rangle$
then do $(X)(f)\left(\psi^{Y}\right)=\psi^{[X, Y]}$
since
$d o(X)(f)\left(\psi^{Y}\right)=\left.\frac{d}{d t}\right|_{t=0} ^{Y}(\exp -t X . f)$
$=\left.\frac{d}{d t}\right|_{t=0}\langle\exp -t X . f, Y\rangle$
$=-\langle X . f, Y\rangle$
$=-\langle f,[X, Y]>$
$=\psi^{[X, Y]}$

By the definition of the moment map i.e.

$$
\begin{aligned}
\tau(x)(X) & =\lambda(X)(x) \\
\text { then } \psi^{X}(\tau(x)) & =\lambda(X)(x)
\end{aligned}
$$

Let $\quad v=\tau_{*} d \phi(Y)(x)$
then $v \psi^{X}=\tau_{*} d \phi(Y)(x)\left(\psi^{X}\right)$
$=\psi^{X}{ }_{\star} \tau_{\star} d \phi(Y)(x)$
$=\lambda(\mathrm{X}){ }_{\star} \mathrm{d} \Phi(\mathrm{Y})(\mathrm{x})$
$=\lambda(X){ }_{\star} X_{\lambda(Y)}{ }^{(x)}$
$=\{\lambda(Y), \lambda(X)\}(x)$
$=\lambda([Y, X])(x)$

So $d \phi(Y) \lambda(X)=\{\lambda(Y), \lambda(X)\}=\lambda([Y, X])$.

But $\lambda([Y, X])(x)=\tau(x)([Y, X])$
$=\psi^{[Y, X]} \tau(x)$

Thus $v \psi^{X}=\psi^{[Y, X]} \tau(x) \quad \forall X \in \mathcal{E}$.
$\Rightarrow \tau_{*} d \phi(Y)(x) \psi^{X}=d \sigma(Y)(\tau(x)) \psi^{X} \quad$ (by (3.3.8))
$\Rightarrow \tau_{k} d \phi(Y)=d \sigma(Y)$

Hence result
Q.E.D.

So by lemma 3.3.21 it is imnediate that
$\left[\tau_{*} d \phi(X), \tau_{*} d \phi(Y)\right]=\tau_{\star}[d \phi(X), d \phi(Y)]$.
From theorem 3.3.9 it is necessary for nilpotent group actions to go to the extension in order for a Poisson action to exist. Notice that the extended Lie algebra $\tilde{\mathcal{L}}$ described after theorem 3.3 .9 is nilpotent if $\mathcal{L}$ is nilpotent.

Theorem 3.3.23 If $(M, \omega)$ is a simply connected symplectic manifold with a transitive Poisson action by the connected nilpotent lie group $\tilde{G}$, with nilpotent Lie algebra $\tilde{L}$, given by

$$
\begin{align*}
& \tilde{\phi}: \tilde{G} \times M \mapsto M \\
& \tilde{\phi}(\tilde{g}, x)=\phi(\mu(\tilde{g}), x) \tag{3.3.9}
\end{align*}
$$

where $\mu: \widetilde{G} \mapsto G$ and $\phi: G \times M \longrightarrow M$ the action given by

$$
\phi\left(\exp t_{1} X_{1_{0}} \ldots 0 \exp ^{t_{k}} X_{k}, x\right)=\exp t_{1} x_{1_{0}} \ldots 0 \exp _{t_{k}} X_{k} \cdot x
$$

where $X_{1}, \ldots, X_{k} \in \mathcal{L}$. Then there exists a sequence of distributions on $M$ satisfying the conditions (1) - (4) of theorem 3.2 .8 and theorem 3.2.9.

Proof. Take the sequence of distributions constructed in lemma 3.3.17 and then take $\tau_{*}{ }^{-1} \Delta_{i}$ where $\tau$ is the moment map, $\tau: M \mapsto O \subset \tilde{L}^{*}$. By proposition 3.3.22 $\tau_{*}{ }^{-1}$ can be taken through the brackets and the validity of conditions (1) - (3) of theorem 3.2.8 follow trivially from lemma 3.3.19.
$\tilde{G}$ - invariance i.e. $\left.\tilde{\phi}_{\tilde{g}}{ }^{\left(\tau_{*}\right.}{ }^{-1} \Delta_{i}\right) \subset \tau_{*}^{-1} \Delta_{i}$ follows from lemma 3.3.5. Since,

$$
\begin{aligned}
& \tau_{0} \tilde{\phi}_{\tilde{g}}=\operatorname{Ad}_{\tilde{g} o}^{*} \tau, \forall \tilde{g} \in \tilde{G} \\
\Rightarrow & \tau_{*} \tilde{\phi}_{\tilde{g} *} \tau_{*}^{-1} \Delta_{i}=\operatorname{Ad}_{\tilde{g}}^{*} * \Delta_{i} \subset \Delta_{i} \quad \psi \tilde{g} \in \tilde{G} \text { by lemraa } 3.3 .17 \\
\Rightarrow & \tilde{\phi}_{\tilde{g} *}\left(\tau_{*}^{-1} \Delta_{i}\right) \subset \tau_{*}^{-1} \Delta_{i}, \forall \tilde{g} \in \tilde{G} .
\end{aligned}
$$

Then by application of the Campbell-Baker-Haus dorff formula i.e.

$$
\exp -t X{ }_{\star} Y(\exp t X . x)=\sum_{i=0}^{\infty} \frac{t^{i}}{i!} a d_{X}^{i} Y(x) \text { and (3.3.9) it }
$$ follows that $\widetilde{G}$ - invariance is equivalent to

$$
\left[X_{H}, \tau_{*}^{-1} \Delta_{i}\right] \subset \tau_{*}^{-1} \Delta_{i} \quad 0 \leqslant i \leqslant n \quad \psi X_{H} \in \mathcal{L}
$$

Condition (4) of theorem 3.2.8 now follows trivially by $\tilde{G}$ - invariance and transitivity since each $X_{k} 1 \leqslant k \leqslant n$ of theorem 3.2 .8 is a Hamiltonian vector field and $X_{\ell}^{i} \subset \tau_{*}^{-1} \Delta_{i}, 1 \leqslant i \leqslant n$.
Q.E.D.

Similarly the existence of polarization for solvable Lie algebras have been studied by Auslander and Kostant [26]. For this it is necessary to introduce complex polarizations. The questions arising from solvable Lie algebras will not be gone into in full detail but an interesting lemma found in the above paper will be stated for the purpose of application in the next chapter.

Let $\mathcal{L}$ be an arbitrary Lie algebra over $\underset{\sim}{R}$ and let $G$ be the corresponding connected Lie group. As before define the alternating bilinear from $B_{f}$ on $\mathcal{L}$ by

$$
\begin{aligned}
& B_{f}(X, Y)=-\left\langle f,[X, Y]>, f \in L^{*}, X, Y \in L \quad\right. \text { and } \\
& \mathcal{L}_{f}=\left\{X \in L: B_{f}(X, Y)=0 \quad \Psi Y \in \mathcal{L}\right\}
\end{aligned}
$$

 valued linear functional on $\mathcal{L}_{\mathrm{C}_{\mathrm{C}}}$. A polarization at $f$ is a complex subalgebra $h \in \mathcal{L}_{\underline{C}}$ such that
(i) $\mathcal{L}_{f} \subset K$ and $\mathcal{L}_{f}$ is invariant under $\mathrm{Ad}_{G_{f}}$
(ii) $\operatorname{dim}_{\underline{C}} L_{\underline{C} / h}=\frac{1}{2} \operatorname{dim}_{\underline{R}} L / L_{\mathrm{f}}$
(iii) <f, $[h, h]>=0$
(iv) $h+\bar{h}$ is a Lie algebra of $L_{C}$ $h \quad$ can be reinterpreted as a maximal isotropic subspace of $\mathcal{L}_{\underset{\sim}{C}}$ relative to $B_{\bar{I}}$.

Let $G$ be a connected, simply connected solvable Lie group with Lie algebra $\mathcal{L}$ and let $s$ denote the nil-radical of $\mathcal{P}$ i.e. maximal nilpotent ideal of $\mathcal{L}$, with Lie group $N$ then the following lemma is proved in Auslander and Kostant [26].

Lemma 3.3.25 Let $N$ be a simply connected nilpotent lie group and let $S$ be it's Lie algebra. Let Aut $S$ be the group of all lie algebra automorphisms of $S$ so that $A d$ is a subgroup of Aut S. Regard Aut $S$ as operating by contragradience on the dual $S^{*}$. Let $f \in S^{*}$. Assume $F$ is a group and a homomorphism $F \mapsto$ Aut $S$ is given (so that $F$ operates on $S$ and $S^{*}$ ) such that (i) the commutator subgroup $F^{l}$ maps into $A d_{N}$ and (ii) F. $f=f$. Then there exists a polarization $h$ at $f$ which is invariant under $F$.

Essentially this lemma can be used along with all the preceding theory on the construction of sequences of involutive distributions to extend to a special case of a solvable group action on a symplectic manifold which will be of importance in the next chapter.

## Chapter 4

## HAMILTONIAN REALIZATIONS OF EINITE VOLTERRA SERIES

Chapter 4 uses the theory developed in chapter 3 to find necessary and sufficient conditions for the existence of Hamiltonian realizations of stationary finite Volterra series. Then using the structure of chapter 3.3 , section two contains a canonical coordinate realization for minimal linear analytic Hamiltonian realizations as well as an algorithm for computing the Volterra kernels if such a system is given. In the last section it is shown that such realizations are closely related to the concept of interconnection as found in electrical network theory.

### 4.1 On finite Volterra Series which admit Hamiltonian realizations

This section begins with a brief review of the work of Brockett [7] and Crouch [5] on nonlinear systems whose input-output behaviour is represented by a Volterra series - in particular linear analytic systems and finite Volterra Series.

So the type of system investigated here is of the form (2.1.1) with the relevant conditions as described in chapter 2.1. The formal volterra series will be written as

$$
y(t)=W_{0}(t)+\int_{0}^{t} W_{1}\left(t, \sigma_{1}\right) u\left(\sigma_{1}\right) d \sigma_{1}+\int_{0}^{t} \int_{0}^{\sigma_{1}} W_{2}\left(t, \sigma_{1}, \sigma_{2}\right)
$$

$$
\mathrm{u}\left(\sigma_{1}\right) u\left(\sigma_{2}\right) \mathrm{d} \sigma_{2} \mathrm{~d} \sigma_{1}
$$

Theorem 4.1.1. (Brockett [7]). Given $T>0, \exists \varepsilon>0$ such that for all locally integrable $u$ satisfying $\int_{0}^{T}|u(s)| d s<\varepsilon$, the following Volterra series converges uniformly and absolutely on $[0, T]$ to the input-output map of system (2.1.1)

$$
y(t)=W_{0}(t)+\sum_{i=1}^{\infty} \int_{0}^{t} \ldots \int_{0}^{\sigma_{i-1}} W_{i}\left(t, \sigma_{1}, \ldots, \sigma_{i}\right) u\left(\sigma_{i}\right) \ldots u\left(\sigma_{1}\right) d \sigma_{i} \ldots d \sigma_{1}
$$

for unique analytic kernel functions

$$
t, \sigma_{1}, \ldots, \sigma_{i} \longrightarrow w_{i}\left(t, \sigma_{1}, \ldots, \sigma_{i}\right)
$$

A series teminating with the term involving the $r^{\text {th }}$ kernel is called a finite Volterra series of length $r$.

In Krener and Lesiak [27], the kernels are shown to be given inductively by the equations

$$
\begin{align*}
& W_{0}^{i}(t, x)=h_{i}\left(\gamma_{f}(t) x\right) \\
& W_{n} j_{0} j_{1} \ldots j_{n}\left(t, \sigma_{1}, \ldots, \sigma_{n}, x\right)=\gamma_{f}\left(-\sigma_{n}\right) \approx g_{j_{n}}\left(\gamma_{f}\left(\sigma_{n}\right) x\right) \\
&  \tag{4,1.2}\\
& \quad\left(W_{n-1} j_{0 .} \ldots j_{n-1}\left(t, \sigma_{1}, \ldots, \sigma_{n-1}, .\right)\right)
\end{align*}
$$

where $\gamma_{f}: \underset{\sim}{R} \times M \longmapsto M$ is the flow of $f \in V(M)$.
The kemels include $x<M$ since each initial state defines an inputoutput map and so the kernels depend on these initial states.

In the case of a finite Volterra series Brockett [7] gives necessary and sufficient conditions in order that the finite Volterra series has a linear analytic realization.

The kernel $W_{r}\left(t, \sigma_{1}, \ldots, \sigma_{r}\right)$ is said to be separable if it can be expressed as a finite sum

$$
W_{r}\left(t, \sigma_{1}, \sigma_{2}, \ldots, \sigma_{r}\right)=\sum_{i=1}^{N} \gamma_{0}^{i}(t) \gamma_{1}^{i}\left(\sigma_{1}\right) \ldots r_{r}^{i}\left(\sigma_{r}\right)
$$

If the component functions $\gamma_{j}^{i}$ are differentiable then it is called differentiably separable.

Further, it is said to be stationary if

$$
W_{r}\left(t, \sigma_{1}, \ldots, \sigma_{r}\right)=W_{r}\left(0, \sigma_{1}-t, \ldots, \sigma_{r}-t\right)
$$

that is the kernels satisfy

$$
\partial / \partial t+\sum_{i=1}^{r}{ }^{\partial} / \partial \sigma_{i} \quad W_{r}\left(t, \sigma_{1}, \ldots, \sigma_{r}\right) \equiv 0
$$

A linear analytic realization is said to be stationary if $f, 8$ and $h$ do not depend explicitly on time and $f(x(0))=0$.

Theorem 4.1.2 (Brockett [7]). A finite Volterra series has a (stationary) linear analytic realization if and only if the kernels are (stationary and) differentiably separable if and only if it has a (stationary) bilinear realization.

```
Further in Brockett [7] for the bilinear realization
\dot{x}=Ax+uBx }x(0)=\mp@subsup{x}{0}{
y=C'x
```

it may be assumed that the subalgebra generated by the matrices $a d_{A}^{k}(B)$ is nilpotent. In fact the bilinear realizations constructed there are such that the matrices $\operatorname{ad}_{A}^{k}(B)$ are strictly lower triangular.

As has been already mentioned, many physical systems evolve on a symplectic manifold, Abraham and Mars den [1] and Arnold [2] being two of the more recent texts which treat this subject from a Symplectic geometric view point. In both these texts only free systems are considered. This then motivates the investigation of linear analytic systems of the form

$$
\dot{x}=X_{f}(x)+u X_{g}(x) \quad x(0)=x_{0} \quad x \in(M, \omega)
$$

$y=g(x)$
where $X_{f}$ and $X_{g}$ are complete analytic Hamiltonian vector fields on a 2 n dimensional analytic connected symplectic manifold ( $M, \omega$ ).

The question to be answered here is under what conditions on the kernels of a stationary finite Volterra series does there exist a realization of the form (4.1.4) and conversely? See Crouch and Irving [28].

Recall that the Poisson bracket of functions on a symplectic manifold $(M, \omega)$ is given by

$$
\begin{equation*}
\{a, b\}=u\left(X_{a}, X_{b}\right)=X_{b}(a), a, b \in C(M) \tag{4.1.5}
\end{equation*}
$$

Theorem 4.1.3 (Crouch [29]) For the system (4.1.4)

$$
\begin{aligned}
& \left.W_{n}\left(t, \sigma_{1}, \ldots, \sigma_{n}, x\right)=\left\{\ldots\left\{h(t), h\left(\sigma_{1}\right)\right\}, h\left(\sigma_{2}\right) \ldots\right\} h\left(\sigma_{n}\right)\right\}(x) \quad(4.1 .6) \\
& W_{1}\left(t, \sigma_{1}, x\right)=\left\{h(t), h\left(\sigma_{1}\right)\right\}(x)
\end{aligned}
$$

where $h(\sigma)(x)=g_{0} \gamma_{f}(\sigma)(x), \gamma_{f}(\sigma)$ the flow of $X_{f}$

Proof. From (4.1.2) and Jacobi's theorem 3.2.3

$$
W_{n}\left(t, \sigma_{1}, \ldots, \sigma_{n}, x\right)=X_{g_{0} \gamma_{f}}\left(\sigma_{n}\right)(x)\left(W_{n-1}\left(t, \sigma_{1}, \ldots, \sigma_{n-1}, x\right)\right)
$$

which by (4.1.5) becomes

$$
\omega(x)\left(X_{W_{n-1}}(x), X_{h\left(\sigma_{n}\right)}(x)\right)=\left\{W_{n-1}\left(t, \ldots, \sigma_{n-1}, x\right), h\left(\sigma_{n}\right)(x)\right\}
$$

So for $n=1$ applying Jacobi's theorem to (4.1.2) it becomes $\left\{h(t), h\left(\sigma_{1}\right)\right\}(x)$. The result now follows easily from induction using the above.
Q.E.D.

In Crouch [5] the following bracket operation applied to Volterra kernels was given, defined inductively by

$$
\begin{aligned}
& W_{n}\left(\sigma_{0} \ldots \sigma_{r}\left[\ldots\left[\sigma_{r+1}, \sigma_{r+2}\right] \ldots \sigma_{q}\right] \sigma_{q+1} \ldots \sigma_{n}\right) \\
= & W_{n}\left(\sigma_{0} \ldots \sigma_{r}\left[\ldots\left[\sigma_{r+1}, \sigma_{r+2}\right] \ldots \sigma_{q-1}\right] \sigma_{q} \sigma_{q+1} \ldots \sigma_{n}\right) \\
- & W_{n}\left(\sigma_{0} \ldots \sigma_{r} \sigma_{q}\left[\ldots\left[\sigma_{r+1}, \sigma_{r+2}\right] \ldots \sigma_{q-1}\right] \sigma_{q+1} \ldots \sigma_{n}\right)
\end{aligned}
$$

These relations are used to give the following necessary condition.

Lemma 4.1.4 The Volterra kernels of a Hamiltonian system (4.1.3) satisfy the following for $n \geqslant 0$

$$
\begin{equation*}
W_{n}\left(\left[\ldots\left[t, \sigma_{1}\right] \ldots \sigma_{q}\right] \sigma_{q+1} \ldots \sigma_{n}, x\right) \tag{4.1.7}
\end{equation*}
$$

$=(q+1) W_{n}\left(t, \sigma_{I} \ldots \sigma_{q}, \sigma_{q+1} \ldots \sigma_{n}, x\right)$ for $1 \leqslant q \leqslant n$

Proof. The proof is by induction on $q$. By (4.1.6) for $q=1$

$$
W_{n}\left(\left[t, \sigma_{1}\right], \sigma_{2} \ldots \sigma_{n}, x\right)=W_{n}\left(t, \sigma_{1}, \sigma_{2} \ldots \sigma_{n}, x\right)
$$

$$
-W_{n}\left(\sigma_{1}, t, \sigma_{2} \ldots \sigma_{n}, x\right)
$$

$\left.=\left\{\ldots\left\{h(t), h\left(\sigma_{1}\right)\right\} h\left(\sigma_{2}\right) \ldots\right\} h\left(\sigma_{n}\right)\right\}(x)$
$\left.-\left\{\ldots\left\{h\left(\sigma_{1}\right), h(t)\right\} h\left(\sigma_{2}\right) \ldots\right\} h\left(\sigma_{n}\right)\right\}(x)$
$\left.=2\left\{\ldots\left\{h(t), h\left(\sigma_{1}\right)\right\} h\left(\sigma_{2}\right) \ldots\right\} h\left(\sigma_{n}\right)\right\}(x)$
by the anti-symmetry of Poisson brackets, $=2 \mathrm{~W}_{\mathrm{n}}\left(\mathrm{t}, \sigma_{1}, \sigma_{2} \ldots \sigma_{\mathrm{rr}} \mathrm{x}\right)$
Assuming (4.1.7) is true for $q$ then by definition and induction

$$
w_{n}\left(\left[\ldots\left[t, \sigma_{1}\right] \ldots \sigma_{q+1}\right] \sigma_{q+2} \ldots \sigma_{n}, x\right)
$$

$=(q+1) W_{n}\left(t, \sigma_{1} \ldots \sigma_{q} \sigma_{q+1} \ldots \sigma_{n}, x\right)$

$$
-W_{n}\left(\sigma_{q+1},\left[\ldots\left[t, \sigma_{1}\right] \ldots \sigma_{q}\right] \sigma_{q+2} \ldots \sigma_{n}, x\right)
$$

By (4.1.6) and the Jacobi identity, applied to Poisson brackets, it is clear that

$$
W_{n}\left(t, \sigma_{1} \ldots \sigma_{r}\left[\sigma_{r+1}, \sigma_{r+2}\right] \sigma_{r+3} \ldots \sigma_{n}, x\right)
$$

$\left.=\left\{\ldots\left\{h(t), h\left(\sigma_{1}\right)\right\} \ldots\right\},\left(\left\{h\left(\sigma_{r+1}, h\left(\sigma_{r+2}\right)\right\}\right)\right\} \ldots h\left(\sigma_{n}\right)\right\}(x)$

Then by repeated application of the Jacobi identity

$$
W_{n}\left(\sigma_{q+1},\left[\ldots\left[t, \sigma_{1}\right] \ldots \sigma_{q}\right] \sigma_{q+2} \ldots \sigma_{n}, x\right)
$$

$\left.\left.=-\left\{\ldots\left(\left\{h(t), h\left(\sigma_{1}\right)\right\} \ldots h\left(\sigma_{q}\right)\right\}\right), h\left(\sigma_{q+i}\right)\right\} \ldots h\left(\sigma_{n}\right)\right\}(x)$
$=-W_{n}\left(t, \sigma_{1}, \ldots, \sigma_{q}, \sigma_{q+1}, \ldots, \sigma_{n}, x\right)$

This verifies (4.1.7) for $q+1$, and hence result.
Q.E.D.

To prove the converse some results are required from the theory of permutation groups and their associated group algebras. These results were proved in Crouch and Irving [28] by Crouch, where they
can be found in detail. For the purpose of this thesis these results will only be stated in order to give a proof of the main theorem, the consequences of which will be used in a later section.

$$
\text { Let } s^{n+1} \text { denote the symmetric group on }(n+1) \text { letters }\{0,1, \ldots, n\}
$$ and ${\underset{\sim}{n}}^{n+1}$ the space of $(n+1)$ - tuples $\left(\sigma_{0}, \sigma_{1}, \ldots, \sigma_{n}\right)$, then the action of $\mathrm{S}^{\mathrm{n+1}}$ on $\mathbb{R}^{\mathrm{n}+1}$ is given by

$$
\rho\left(\sigma_{0}, \sigma_{1}, \ldots, \sigma_{n}\right)=\left(\sigma_{\rho(0)}, \sigma_{\rho(1)}, \ldots, \sigma_{\rho(n)}\right)
$$

Let $A^{n+1}$ be the group algebra generated by $S^{n+1}$, and $F^{n+1}$ the space of real valued functions on $\mathbb{R}^{n+1}$

$$
\left(\sigma_{0}, \sigma_{1}, \ldots, \sigma_{n}\right) \longmapsto f\left(\sigma_{0}, \sigma_{1}, \ldots, \sigma_{n}\right) .
$$

$A^{n+1}$ is represented as an algebra of linear operators on $F^{n+1}$ by setting $(\rho f)\left(\sigma^{n}\right)=f\left(\rho\left(\sigma^{n}\right)\right)$ for $\rho \in s^{\mathrm{n}+1}$ and $\sigma^{n}=\left(\sigma_{0}, \sigma_{1}, \ldots, \sigma_{n}\right)$

Further, $A^{k+1}$ may be considered as a subalgebra of $A^{n+1}$ for $k<\pi$, by defining the action of $A^{k+1}$ on ${\underset{\sim}{n}}^{n+1}$ to be that obtained by fixing $\sigma_{k+1}, \ldots, \sigma_{n}$.

Let $p_{k} \in s^{k+1}$ be the permutation defined by the cycle ( $k, k-1, \ldots, 1,0$ ), for $k \geqslant 1$ and set $\rho_{o}$ to be the identity in $s^{k+1}$. If $\rho_{k} \in A^{k+1}$ considered as subalgebra of $A^{n+1}$ then

$$
o_{k}\left(\sigma_{o} \sigma_{1} \ldots \sigma_{k}, \sigma_{k+1}, \ldots, \sigma_{n}\right)=\left(\sigma_{k} \sigma_{k-1} \ldots \sigma_{0} \sigma_{k+1} \ldots \sigma_{n}\right)
$$

Lemma 4.1.5 If $\beta_{n}=\left(\rho_{o}-\rho_{1}\right)\left(\rho_{o}-\rho_{2}\right) \cdots\left(\rho_{o}-\rho_{n}\right) \subset A^{n+1}$
then $\beta_{n}^{2}=(n+1) \beta_{n}$.

$$
\text { Now let } B_{k}^{a,} b \ldots c \text { - for } a, b, \ldots, c \text { belonging to some index set - }
$$ be the function

$$
\sigma^{k} \longmapsto \gamma_{a}\left(\sigma_{o}\right) \gamma_{b}\left(\sigma_{1}\right) \ldots \gamma_{c}\left(\sigma_{k}\right)
$$

and set $W_{k}^{a, b} \ldots c=\beta_{k} B_{k}^{a, b \ldots c}$, then

Theorem 4.1.6 If $W_{n}\left(\sigma_{0}, \sigma_{1}, \ldots, \sigma_{n}, x\right)$ is a differentiably separable kernel satisfying the conditions of lemma 4.1.4 then

$$
W_{n}\left(\sigma_{0}, \sigma_{1} \ldots \sigma_{n}, x\right)=\sum_{i=1}^{N} W_{n} a_{i} b_{i} \cdots c_{i}\left(\sigma_{0}, \ldots, \sigma_{n}, x\right)
$$

for differentiable functions $\gamma_{a_{i}} \ldots \gamma_{c_{i}} \quad 1 \leqslant i \leqslant N$. Moreover, this is miquely defined.

Now with the aid of lemma 4.1.5 and theorem 4.1.6 the main result may be proven.

Theorem 4.1.7 A finite Volterra series in which the kernels $W_{n} l \leq n \leqslant r$, satisfying the stationarity condition (4.1.3) has a symplectic realization in the form of system (4.1.4) if and only if
(i) The kernels are differentiably separable
(ii) The kernels satisfy condition (4.1.7) of lemma 4.1.4, that is

$$
W_{n}\left(\left[\ldots\left[t, \sigma_{1}\right] \ldots \sigma_{q}\right] \sigma_{q+1} \ldots \sigma_{n} x\right)=(q+1) W_{n}\left(t, \sigma_{1}, \ldots, \sigma_{n}, x\right) ; 1 \leq q \leq n
$$

Proof. Implication follows from theorem 4.1.2 and lemma 4.1.4
Conversely, if the kernels are differentiably separable then by Brockett [7] the Volterra series has a bilinear realization. It therefore follows that the kernels may be written in the form

$$
\begin{equation*}
W_{n}\left(t, \sigma_{1}, \ldots, \sigma_{n}\right)=\sum_{i=1}^{N} \gamma_{n}^{i}\left(t-\sigma_{1}\right) \gamma_{n-1}^{i}\left(\sigma_{1}-\sigma_{2}\right) \ldots \gamma_{1}^{i}\left(\sigma_{n-1}-\sigma_{n}\right) \tag{4.1.8}
\end{equation*}
$$

where each function $\gamma_{k}^{i}$ has the form

$$
\gamma_{k}^{i}(\sigma)=c_{i k}^{\prime} e^{A_{i k}{ }^{\sigma}} b_{i k}
$$

for suitable $c_{i k}, b_{i k} \in \mathbb{R}^{s}$ and $A_{i k} \in \mathbb{R}^{s \times s}$ for some $s>0$. This sum may be split into minimal groups still having the form given by the right hand side of equation (4.1.8), and also satisfying condition (4.1.7). Clearly each group is also differentiably separable. Select one such group and apply theorem 4.1 .6 to see that it must be expressible in the form,

$$
W_{n}^{\star}\left(t, \sigma_{1}, \ldots, \sigma_{n}\right)=\beta_{n} \gamma_{n}\left(t-\sigma_{1}\right) \gamma_{n-1}\left(\sigma_{1}-\sigma_{2}\right) \ldots \gamma_{1}\left(\sigma_{n-1}-\sigma_{n}\right)
$$

for suitable functions $\gamma_{1} \ldots \gamma_{n}$. This is the sum of $2^{n}$ terms, with each term constructed from the product of the functions $\gamma_{1} \ldots \gamma_{n}$.

$$
\text { Let } \gamma_{i}(t)=c_{i}^{\prime} e^{A_{i} t} b_{i}, \quad 1 \leqslant i \leqslant n \text {. Now a realization of the }
$$ input-output map

$$
\begin{equation*}
g^{*}(t)=\int_{0}^{t} \cdots \int_{0}^{\sigma_{n-1}} W_{n}^{*}\left(t, \sigma_{1} \ldots, \sigma_{n}\right) u\left(\sigma_{n}\right) \ldots u\left(\sigma_{1}\right) d \sigma_{n} \ldots d \sigma_{1} \tag{4.1.9}
\end{equation*}
$$

is constructed.

With $W_{n}^{*}$ replaced by

$$
\begin{equation*}
\gamma_{n}\left(t-\sigma_{1}\right) \gamma_{n-1}\left(\sigma_{1}-\sigma_{2}\right) \ldots \gamma_{1}\left(\sigma_{n-1}-\sigma_{n}\right) \tag{4.1,10}
\end{equation*}
$$

the above input-output map realizes the following system, see Brockett [7].

$$
\begin{array}{ll}
\dot{q}_{1}=A_{1} q_{1}+b_{1} u & q_{1}(0)=0 \\
\dot{q}_{2}=A_{2} q_{2}+b_{2} c_{1}^{\prime} q_{1} u & q_{2}(0)=0 \\
\vdots & q_{n}(0)=0
\end{array}
$$

$y_{1}=c_{n}^{\prime} q_{n}$
Now, consider the composite system obtained by adjoining these equations to the equations

$$
\begin{array}{ll}
-\dot{p}_{1}=A_{1}^{\prime} p_{1}+c_{1} b_{2}^{\prime} p_{2} u & p_{1}(0)=0 \\
-\dot{p}_{2}=A_{2}^{\prime} p_{2}+c_{2} b_{3}^{\prime} p_{3}^{u} & p_{2}(0)=0 \\
\vdots & \\
-\dot{p}_{n}=A_{n}^{\prime} p_{n}+c_{n}^{u} & p_{n}(0)=0
\end{array}
$$

This composite system may be expressed as a Hamilonian system by defining the Hamiltonian functions.

$$
\begin{aligned}
& H(q, p)=p_{1}^{\prime} A_{1} q_{1}+\ldots+p_{n}^{\prime} A_{n} q_{n} \\
& H_{u}(q, p)=p_{1}^{\prime} b_{1}+c_{n}^{\prime} q_{n}+c_{1}^{\prime} q_{1} b_{2}^{\prime} p_{2}+\ldots+c_{n-1}^{\prime} q_{n-1} b_{n}^{\prime} p_{n}
\end{aligned}
$$

$$
\text { If } q_{i} \in \mathbb{R}^{n_{i}} \text { and } \sum_{i=1}^{n} n_{i}=N \text {, the following system on }{\underset{\sim}{2}}^{2 N}
$$

coincides with the composite system above,

$$
\dot{q}=A q+u N q+b u \quad q(0)=0
$$

$$
\begin{equation*}
-\dot{p}=A^{\prime} p+u^{\prime} q+c u \quad p(0)=0 \tag{4.1.11}
\end{equation*}
$$

$$
y_{2}=b^{\prime} p+c^{\prime} q+p^{\prime} N q
$$

where

$$
\begin{aligned}
& b=\left(\begin{array}{c}
b_{1} \\
0 \\
0 \\
0
\end{array}\right) \quad, \quad c=\left(\begin{array}{l}
0 \\
0 \\
0 \\
c_{n}
\end{array}\right) \quad, \quad\left(\begin{array}{c}
q_{1} \\
\vdots \\
\vdots \\
q_{n}
\end{array}\right) \quad, \quad p=\left(\begin{array}{c}
p_{1} \\
0 \\
\vdots \\
p_{n}
\end{array}\right)
\end{aligned}
$$

With respect to the canonical symplectic structure $J$ on $\mathbb{R}^{2 N}$ this is a Hamiltonian system, since

$$
\begin{aligned}
\dot{q} & =\frac{\partial H}{\partial p}(q, p)^{\prime}+\frac{\partial H_{u}}{\partial p}(q, p)^{\prime} \\
-\dot{p} & =\frac{\partial H}{\partial q}(q, p)^{\prime}+\frac{\partial H_{u}}{\partial q}(q, p)^{\prime} \\
y_{2} & =H_{u}(q, p) .
\end{aligned}
$$

It now remains to show that the input-output map of system (4.1.11) coincides with that in equation (4.1.9). Note that the output function in system (4.1.11) contains the term $c^{\prime} q=c_{n}^{\prime} q_{n}=y_{1}$, and so the input-output map of system (4.1.11) contains a term with Volterrakernel given by equation ( $4,1,10$ ) . It is an easy calculation to check that if $b_{r+1}^{\prime} p_{r+1}$ and $c_{r}^{\prime} q_{r}$ are viewed as outputs of system (4.1.11) then

$$
\begin{array}{r}
c_{r}^{\prime} q_{r}=\int_{0}^{t} \ldots \int_{0}^{\sigma_{r-1}} \gamma_{r}\left(t-\sigma_{1}\right) \gamma_{r-1}\left(\sigma_{1}-\sigma_{2}\right) \ldots \gamma_{1}\left(\sigma_{r-1}-\sigma_{r}\right) \\
u\left(\sigma_{r}\right) \ldots u\left(\sigma_{1}\right) d \sigma_{r} \ldots d \sigma_{1}
\end{array}
$$

$$
b_{r+1}^{\prime} p_{r+1}=\int_{0}^{t} \cdots \int_{0}^{\sigma_{n-r-1}} \gamma_{r+1}^{*}\left(t-\sigma_{1}\right) \ldots \gamma_{n}^{*}\left(\sigma_{n-r-1}-\sigma_{n-r}\right)
$$

$$
u\left(\sigma_{n-r}\right) \ldots u\left(\sigma_{1}\right) d \sigma_{n-r} \ldots d \sigma_{1}
$$

where $\gamma_{k}^{*}(\sigma)=-\gamma_{k}(-\sigma)$. It follows that $y_{\hat{z}}(t)$ is a Volterra series with one term involving $n$ iterated integrals of $u$, and associated Volterra kernel constructed from a sum of terms, each of which is in turn constructed from a product of all the functions $\gamma_{1} \ldots \gamma_{n}$. It follows from Crouch [30], that the contribution of the term
$c_{r}^{\prime} q_{r} b_{r+1}^{\prime} p_{r+1}, 1 \leqslant r \leqslant n-1$ is exactly $n!/ r!(n-r)!$ such terms. Thus the Volterra defining $y_{2}(t)$ as Volterra series consists of

$$
1+\sum_{r=1}^{n-1} n!/ r!(n-r)!+1=2^{n}
$$

such terms. As system (4.1.11) is Hamiltonian, lemma 4.1.4 and the uniqueness of theorem 4.1.6 imply that the input-output map for system (4.1.11) coincides with the input-output map in equation (4.1.9) as claimed.

The proof is completed by noting that given a finite Volterra series satisfying the hypothesis of the theorem, it may be split into a finite number of terms, each of which is realizable by a Hamiltonian system of the form given by system (4.1.11). Obviously, the composite system is still Hamiltonian, and realizes the given finite Volterra series. Q.E.D.

System (4.1.11) may be rewritten with the introduction of two new variables $q_{0}$ and $p_{0}$ as

$$
\begin{align*}
& y=p^{* \prime} N^{*} q^{*}+\frac{1}{2} q^{*^{\prime}} Q^{*} q \\
& q_{0}(0)=1 \quad p_{0}(0)=p^{0} \tag{4.1.12}
\end{align*}
$$

where $N^{*}=\left(\begin{array}{ll}0 & 0 \\ b & N\end{array}\right), Q^{*}=\left(\begin{array}{ll}0 & c^{\prime} \\ c & 0\end{array}\right), \quad q^{*}=\binom{q_{0}}{q}, \quad p^{*}=\binom{p_{0}}{p}$

System (4.1.12) has the same input-output map as system 4.1.11, since

$$
=y_{2}
$$

So system (4.1.12) is exactly the form found in chapter 3.1 if the Lie algebra $\left\{\operatorname{ad}_{A}^{k} \bar{B}\right\}_{\text {L. A. }}$. is assumed to be nilpotent. Here, of course, the matrices are block diagonal or block lower/upper triangular whereas before each was an individual element. Where

$$
\bar{A}=\left(\begin{array}{cccc}
0 & 0 & 0 \\
0 & A & & \\
& 0 & 0 & 0 \\
& & 0 & -A^{\prime}
\end{array}\right), \bar{B}=\left(\begin{array}{cc}
N^{*} & 0 \\
-Q^{*} & -N^{*^{\prime}}
\end{array}\right)
$$

$$
\begin{aligned}
& p^{*^{\prime}} N^{*} q^{*}+\frac{1}{2} q^{*^{\prime}} Q^{*} q=\left(p_{0}^{\prime} p^{\prime}\right)\left(\begin{array}{ll}
0 & 0 \\
b & N
\end{array}\right)\left(\begin{array}{l}
q \\
q_{0} \\
q
\end{array}\right) \\
& +\frac{1}{2}\left(q_{0}^{\prime} q^{\prime}\right)\left(\begin{array}{cc}
0 & 0 \ldots c_{n}^{\prime} \\
0 & \\
\vdots & 0
\end{array}\right)\binom{q_{0}}{c_{n}} \\
& =\left(p_{1}^{\prime} b_{1} p^{\prime} N\right)\binom{q}{q}+\frac{1}{2}\left(q_{n}^{\prime} c_{n} 0 . .00 q_{0}^{\prime} c_{n}^{\prime}\right)\left(\begin{array}{l}
q \\
0 \\
q
\end{array}\right) \\
& =p_{1}^{\prime} b_{1} q_{0}+p^{\prime} N q+\frac{1}{2} q_{n}^{\prime} c_{n} q_{0}+\frac{1}{2} q_{0}^{\prime} c_{n}^{\prime} q_{n} \\
& =p_{1}^{\prime} b_{1}+p^{\prime} N q+c_{n}^{\prime} q_{n} \\
& \text { Since } q_{0}=1
\end{aligned}
$$

Example 4.1.8 Consider the system

$$
\begin{array}{ll}
\dot{q}_{1}=q_{1}+u & q_{1}(0)=0 \\
\dot{q}_{2}=q_{2}+u q_{1} & q_{2}(0)=0 \\
-\dot{p}_{1}=p_{1}+u p_{2} & p_{2}(0)=0 \\
-\dot{p}_{2}=p_{2}+u & p_{2}(0)=0 \\
y=p_{2} q_{1}+p_{1}+q_{2} &
\end{array}
$$

on the symplectic manifold $\left(I^{*} R^{2}, \sum_{i=1}^{2} d q_{i}, d p_{i}\right)$

This is clearly a Hamiltonian system i.e. it is of the form of equation (4.1.4), with Hamiltonian functions given by

$$
\begin{aligned}
& f=p_{1} q_{1}+p_{2} q_{2} \\
& g=p_{2} q_{1}+p_{1}+q_{2} \\
& \text { and } X_{f}=q_{1}{ }^{\partial} / \partial q_{1}+q_{2} \partial / \partial q_{2}-p_{1}^{\partial / \partial} p_{1}-p_{2} \partial / \partial p_{2} \\
& x_{g}=\partial / \partial q_{1}+q_{1} \partial / \partial q_{2}-p_{2} \partial / \partial p_{1}-\partial / \partial p_{2} \\
& S=\left\{\operatorname{ad}_{X_{f}}, X_{g}\right\}_{L . A} \text {. is the ideal of } \dot{L} \text { generated by }
\end{aligned}
$$

$$
\begin{aligned}
& x_{g}=\partial / \partial q_{1}+q_{1} \partial / \partial q_{2}-p_{2} \partial / \partial p_{1}-\partial / \partial p_{2}, \partial / \partial q_{1}+\partial / \partial p_{2}, \\
& \partial / \partial q_{2}+\partial / \partial p_{1}, \partial / \partial q_{1}-\partial / \partial p_{2}, \partial / \partial q_{2}-\partial / \partial p_{1}
\end{aligned}
$$

which is obviously strongly accessible since $\operatorname{dim} S(x)=T T^{*} R^{2} \forall x \in T^{*} R^{2}$ and thus by Van der Schaft [31] weakly observable. Note also that $S$ is a nilpotent Lie algebra.

Since,

$$
\begin{aligned}
& q_{1}(t)=\int_{0}^{t} e^{t-\sigma_{1} u\left(\sigma_{1}\right) d \sigma_{1}} \\
& q_{2}(t)=\int_{0}^{t} \int_{0}^{\sigma_{1}} e^{t-\sigma_{2} u\left(\sigma_{2}\right) u\left(\sigma_{1}\right) d \sigma_{2} d \sigma_{1}} \\
& P_{1}(t)=\int_{0}^{t} \int_{0}^{\sigma_{1}} e^{-\left(t-\sigma_{2}\right)} u\left(\sigma_{2}\right) u\left(\sigma_{1}\right) d \sigma_{2} d \sigma_{1} \\
& P_{2}(t)=-\int_{0}^{t} e^{-\left(t-\sigma_{1}\right)} u\left(\sigma_{1}\right) d \sigma_{1}
\end{aligned}
$$

the Volterra series has the form

So $W_{2}\left(t, \sigma_{1}, \sigma_{2}, x\right)=e^{\sigma_{2}-t}+e^{t-\sigma_{2}}-e^{\sigma_{2}-\sigma_{1}}-e^{\sigma_{1}^{-\sigma} 2}$
which is seen to satisfy condition (4.1.7) of lemma 4.1.4 as follows,

$$
\begin{aligned}
W_{2}\left(\left[t, \sigma_{1}\right], \sigma_{2}\right) & =W_{2}\left(t, \sigma_{1}, \sigma_{2}\right)-W_{2}\left(\sigma_{1}, t, \sigma_{2}\right) \\
& =e^{\sigma_{2}-t}+e^{t-\sigma_{2}}-e^{\sigma_{2}-\sigma_{1}}-e^{\sigma_{1}-\sigma_{2}} \\
& -e^{\sigma^{-\sigma_{1}}-e^{\sigma_{1}-\sigma_{2}}+e^{\sigma_{2}-t}+e^{t-\sigma_{2}}} \\
& =2\left(e^{\sigma_{2}-t}+e^{t-\sigma_{2}}-e^{\sigma_{2}-\sigma_{1}}-e^{\sigma_{1}-\sigma_{2}}\right) \\
\text { and } W_{2}\left(\left[\left[t, \sigma_{1}\right]\right.\right. & \left.\left.\left.\sigma_{2}\right]\right]\right)=W_{2}\left(\left[t, \sigma_{1}\right], \sigma_{2}\right)-W_{2}\left(\sigma_{2}\left[t, \sigma_{1}\right]\right) \\
& =2 W_{2}\left(t, \sigma_{1}, \sigma_{2}\right)-W_{2}\left(\sigma_{2}, t, \sigma_{1}\right) \\
& +W_{2}\left(\sigma_{2}, \sigma_{1}, t\right) \\
& =2 W_{2}\left(t, \sigma_{1}, \sigma_{2}\right)-e^{\sigma_{1}-\sigma_{2}}-e^{\sigma_{2}-\sigma_{1}+e^{\sigma_{1}-t}} \\
& +e^{t-\sigma_{1}}+e^{t-\sigma_{2}}+e^{\sigma_{2}-t}-e^{t-\sigma_{1}}-e^{\sigma_{1}-t} \\
& =3 W_{2}\left(t, \sigma_{1}, \sigma_{2}\right) \text { as required. }
\end{aligned}
$$

### 4.2 A Canonical Realization

In Crouch [5] dynamical realizations of finite Volterra series are investigated. In particular canonical coordinate charts are chosen to express strongly accessible, observable linear analytic realizations of finite Volterra series. Here the work of chapter 3 is applied to find realizations of the above type which include a symplectic structure.

Sussmann [9] guarantees the existence and uniqueness of minimal realizations i.e. for analytic systems accessible and observable realizations. However, for Hamiltonian systems the following observations from Goncalves [16] must be made.

Let $G$ be a connected Lie group and suppose there exists a Poisson action of $G$ on a strongly symplectic $G$-space, $(M, \omega)$. Note that ( $M, \omega$ ) is strongly symplectic if for example $M$ is simply connected. The fact that the action is Poisson means that there exists a moment map $\tau$.

Let ( $\left[, x_{0}\right.$ ) represent an initialized analytic system of complete vector fields on the symplectic manifold ( $M, \omega$ ), and let $O$ be the orbit in $\mathcal{L}^{*}$, the Lie algebra of $G$, of $\tau\left(x_{0}\right)$ by the coadjoint action of $G$. Let ( $\sum_{0}, \tau\left(x_{0}\right)$ ) be the corresponding system on che co-adjoint orbit. Then with strong equivalence as defined in chapter 2.1 it is proved in Goncalves [16],

Theorem 4.2.1 $\left(\sum_{0}, \tau\left(x_{0}\right)\right)$ is strongly equivalent to ( $\sum, x_{0}$ ), and moreover it is an accessible Hamiltonian system.

It is noted in [16] that $\sum_{0}$ is accessible but may not be observable or weakly observable. However, the following theorem may be found in either Goncalves [16] or Van der Schaft [31].

Theorem 4.2.2 If $\sum_{0}$ is strongly accessible it is weakly observable.
The following results are to be found in Crouch [5], and will prove useful throughout this section.

Lema 4.2.3 Given a strongly accessible realization of a finite Volterra series of length $n$ then
(i) The Volterra series has length n when evaluated at any point in the state space $M$.
(ii) The kernel $W_{n}$ depends only on the time parameters $t, q_{1}, \ldots, \sigma_{n}$, not on the state $x \in M$.

Theorem 4.2.4 Given a strongly accessible weakly observable realization of a finite Volterra series of length $n$, $S$ is a nilpotent lie algebra, with a descending central series of length less than or equal to $n$, and the Lie algebra $\mathcal{L}$ is solvable and finite dimensional.

Recall that $S$ is the ideal in $\hat{\mathcal{L}}$ generated by the vector fields $g_{1}, \ldots, g_{r}$ of (2.1.1).

Since the systems under consideration are strongly accessible, the connected Lie subgroup $N$ of $G$, corresponding to $S$, acts transitively on $M$ and therefore $M$ is analytically diffeomorphic to the homogeneous space $N / N_{x_{0}}$ where $N_{x_{0}}$ is the isotropy subgroup
$N_{x_{0}}=\left\{g \in N: g . x_{0}=x_{0}\right\}$
As $N$ is nilpotent this is a nilmanifold. For further details of
this point see Helgas on [15].

One of the most itportant theorems in Crouch [5] states,

Theorem 4.2.5 A strongly accessible observable realization of a finite Volterra series has a state space which is diffeomorphic to a Euclidean space.

This shows that the underlying manifold is simply connected.
Finally from [5],

Corollary 4.2.6 A strongly accessible analytic realization of a finite Volterra series is observable if and only if it is weakly observable.

Now the results of chapter 3 are utilized to find a canonical realization of a linear analytic Hariltonian strongly accessible, observable realization.

As N acts as a nilpotent transitive action on M by strong accessibility and theorem 4.2.4, then theorem 3.3 .9 says it must be extended in order to make the action Poisson i.e. let $\tilde{N}$ be the nilpotent connected Lie group of the Lie subalgebra $\tilde{S}=S \times R$. Further note that ( $M, w$ ) is always strongly symplectic as $M$ in this case is simply connected by theorem 4.2.5.

Recall that the extended action was defined as follows, let

$$
\mu: \tilde{N} \longmapsto N
$$

and if $\tilde{g} \subset \tilde{N}, x \in M$ then the action is

$$
\tilde{g} \cdot x=\mu(\tilde{g}) \cdot x
$$

So obviously $\tilde{N}$ acts transitively on $M$.

A particularly interesting case of when accessibility and strong accessibility are equivalent is found when considering the concept of stationarity of linear analytic systems as described in Brockett [7].

Consider a linear analytic Hamiltonian system of the form

$$
\begin{aligned}
\dot{x}=X_{H}(x)+\sum_{i=1}^{T} u_{i} X_{H_{u_{i}}}(x) \quad & x \in(M, \omega) \\
& x(0)=x_{0}
\end{aligned}
$$

$$
\begin{equation*}
y^{i}=H_{u_{i}}(x) \tag{4.2.1}
\end{equation*}
$$

where ( $M, \omega$ ) is a $2 m$ - dimensional symplectic manifold.
Then it is stationary if $X_{H}\left(x_{0}\right)=0$, which is equivalent to expt $X_{H} \cdot X_{0}=X_{o}$ where expt $X_{H}$ is the flow of $X_{H}$. By Sussmann [32] it then follows that $\operatorname{dim} S(x)=\operatorname{dim} L(x), \psi x \in M$ and thus accessibility and strong accessibility are equivalent.

Furthermore, if this system is lifted to the corresponding system on ( $0, \omega_{0}$ ) it too has this property, since if $X_{H} \in S$, let $\tilde{X}_{H} \in \widetilde{S}$ then by definition of the action it follows that

$$
\begin{aligned}
& \qquad \operatorname{expt} \tilde{X}_{H} \cdot x_{0}=\mu\left(\operatorname{expt} \tilde{X}_{H}\right) \cdot x_{0}=\exp t X_{H} \cdot x_{0}=x_{0} \\
& \text { so by } A d^{\star}-\text { equivalence (see lema 3.3.5) }
\end{aligned}
$$

$$
\begin{aligned}
& A d^{\operatorname{expt}} \tilde{X}_{H}^{*}{ }_{\circ} \tau\left(x_{0}\right)=\tau_{0} \operatorname{expt} \tilde{X}_{H} \cdot x_{o} \\
\Rightarrow \quad & A d_{\operatorname{expt}}^{*} \tilde{X}_{H} f=\tau\left(x_{o}\right)=f \\
\Rightarrow & \sigma_{\operatorname{expt}} \tilde{X}_{H}(f)=f
\end{aligned}
$$

where $\sigma: \tilde{N} \times 0 \longmapsto 0$ as in chapter 3.3. Differentiating gives do $\left(\tilde{X}_{H}\right)(f)=0$ as required.

Thus from theorem 4.2.1, theoren 4.2.2 and the fact that a stationary finite Volterra series which admits a linear analytic Hamiltonian realization is such that the linear analytic system is stationary, see Brockett [7], then the following corollary holds.

Corollary 4.2.7 A stationary linear analytic Hamiltonian realization of a stationary finite Volterraseries on a symplectic manifold ( $M, \omega$ ) is strongly equivalent to a minimal i.e. accessible and observable stationary linear analytic Hamiltonian realization on ( $0, \omega_{0}$ ).

The remainder of this section is mostly concerned with systems of the form (4.2.1) which have nilpotent Lie groups associated to them, and to which the results of chapter 3.3 apply. There are basically two cases to be considered. The first corresponds to the case where the system is strongly accessible and $\mathcal{L}=\mathrm{S}$. The second case corresponds to the system being strongly accessible but $\mathcal{L} \neq \mathrm{S}$.

First an in depth consideration of the case $s=\mathcal{L}$ on a strongly accessible, observable linear analytic Hamiltonian realization of the form 4.2.1 of a finite Volterra series. Theorem 4.2 .5 gives $S$ nilpotent
and theorems 3.3.23 and 3.2.9 give rise to the following coordinate system.

Let $\left(q_{1}, \ldots, q_{n}, p_{1}, \ldots, p_{n}\right)$ be canonical coordinates on the symplectic manifold ( $M^{2 \mathrm{~m}}, \omega$ ) with

$$
\begin{array}{ll}
q_{i}=q_{1}^{i}, \ldots, q_{r_{i}}^{i} & 1 \leqslant i \leqslant n \\
p_{i}=p_{1}^{i}, \ldots, p_{r_{i}}^{i} & 1 \leqslant i \leqslant n
\end{array}
$$

then the realization can be given by

$$
\begin{aligned}
& \dot{q}_{1}=f_{1}\left(q_{1}\right) \quad g_{1}^{i}\left(q_{1}\right) \\
& \begin{array}{ll}
\dot{q}_{2}=f_{2}\left(q_{1}, q_{2}\right) & g_{2}^{i}\left(q_{1}, q_{2}\right) \\
\vdots & \\
\vdots &
\end{array} \\
& \dot{q}_{n}=f_{n}\left(q_{1}, \ldots, q_{n}\right) \quad g_{n}^{i}\left(q_{1}, \ldots, q_{n}\right) \\
& -\dot{p}_{1}=\frac{\partial f^{\prime}}{\partial q_{1}} p_{1}+\ldots+\frac{\partial f^{\prime}}{\partial q_{1}^{\prime}} p_{n}+\frac{\partial f^{\prime}}{\partial q_{1}} \\
& +\sum_{i=1}^{q} u_{i} \frac{\partial g_{l}^{i,}}{\partial q_{1}} \quad p_{1}+\ldots+\frac{\partial g_{n}^{i \prime}}{\partial q_{1}} p_{n}+\frac{\partial g}{\partial q_{1}}{ }_{1}^{i_{1}} \\
& -\dot{p}_{n}=\frac{\partial f^{\prime}}{\partial q_{n}^{\prime}} p_{n}+\frac{\partial f^{\prime}}{\partial q_{n}^{\prime}}+\sum_{i=1}^{r} u_{i} \frac{\partial g_{n}^{i}}{\partial q_{n}^{\prime}} p_{n}+\frac{\partial g^{i}}{\partial q_{n}^{\prime}}
\end{aligned}
$$

where $x=\left(q_{1}, \ldots, q_{n}, p_{1}, \ldots, p_{n}\right)^{\prime} \quad x(0)=x_{0}$
and $H=p_{1}^{\prime} f_{1}\left(q_{1}\right)+\ldots+p_{n}^{\prime} f_{n}\left(q_{1}, \ldots, q_{n}\right)+f_{o}\left(q_{1}, \ldots, q_{n}\right)$
$y^{i}=H_{u_{i}}=p_{1}^{\prime} g_{1}^{i}\left(q_{1}\right)+\ldots+p_{n}^{\prime} g_{n}^{i}\left(q_{1}, \ldots, q_{n}\right)+g_{o}^{i}\left(q_{1}, \ldots, q_{n}\right)$
where $f_{j}$ and $g_{j}$ are analytic functions in $q_{1}, \ldots, q_{j}, l \leqslant j \leqslant n$ and $f_{o}$ and $g_{o}^{i}$ are analytic functions in $q_{1}, \ldots, q_{n}$.

In fact under further scrutiny the following theorem is obtained.

Theorem 4.2.8 Any strongly accessible observable linear analytic Hamiltonian realization of a finite Volterra series which satisfies condition 4.1 .7 of lemma 4.1 .4 and such that $S=L$ has a canonical realization as follows, let $\left(q_{1}, \ldots, q_{n}, p_{1}, \ldots, p_{n}\right)$ be canonical coordinates on the state space $\left(\mathrm{m}^{2 \mathrm{~m}}, \omega\right)$ where

$$
\begin{array}{ll}
q_{i}=q_{1}^{i}, \ldots, q_{r_{i}}^{i}, & 1 \leqslant i \leqslant n \\
p_{i}=p_{1}^{i}, \ldots, p_{r_{i}}^{i}, & 1 \leqslant i \leqslant n
\end{array}
$$

then the realization is expressed by

$$
\begin{array}{ll}
\dot{q}_{1}=f_{1} & g_{1}^{i} \\
\dot{q}_{2}=f_{2}\left(q_{1}\right) & \sum_{i=1}^{r} \\
\vdots & g_{i}^{i}\left(q_{1}\right) \\
\dot{q}_{n}=f_{n}^{i}\left(q_{1}, \ldots, q_{n-1}\right) & g_{n}^{i}\left(q_{1}, \ldots, q_{j-1}\right) \\
&
\end{array}
$$

$-\dot{p}_{1}=\frac{\partial f^{\prime}}{\partial q_{1}^{\prime}} p_{1}+\ldots+\frac{\partial f^{\prime}}{\partial q_{1}^{n}} \quad p_{n}+\frac{\partial f^{\prime}}{\partial q_{1}^{\prime}}$

$$
+\sum_{i=1}^{r} u_{i} \frac{\partial g_{2}^{i}}{\partial q_{1}^{\prime}} \quad p_{2}+\ldots+\frac{\partial g_{1}^{i}}{\partial q_{1}^{\prime}} p_{n}+\frac{\partial g_{0}^{i}}{\partial q_{1}^{\prime}}
$$

$-\dot{p}_{n}=\frac{\partial f^{1}}{\partial q_{n}^{o}}+\sum_{i=1}^{r} u_{i} \frac{\partial g_{o}^{i}}{\partial q_{n}^{o}}$
set $x=\left(q_{1}, \ldots, q_{n}, p_{1}, \ldots, p_{n}\right)^{\prime}$ then $x(0)=x_{0}$
$H=p_{1}^{\prime} f_{1}+p_{2}^{\prime} f_{2}\left(q_{1}\right)+\ldots+p_{n}^{\prime} f_{n}\left(q_{1}, \ldots, q_{n-1}\right)+f_{0}\left(q_{1}, \ldots, q_{n}\right)$
$y^{i}=H_{u_{i}}=p_{1}^{\prime} g_{1}^{i}+p_{2}^{\prime} g_{2}^{i}\left(q_{1}\right)+\ldots+p_{n}^{\prime} g_{n}^{i}\left(q_{1}, \ldots, q_{n-1}\right)+g_{0}^{i}\left(q_{1}, \ldots, q_{n}\right)$
where $f_{0}$ and $g_{j}^{i}$ are analytic functions in $q_{1}, \ldots, q_{j-1}, 1 \leqslant j \leqslant n$ and $f_{0}$ and $g_{o}^{i}$ are analytic functions in $q_{1}, \ldots, q_{n}$. Also
$2 \sum_{i=1}^{n} r_{i}=2 m$ the dimension of $M$.
Proof. Construct the coadjoint orbit on $\widehat{S}^{*}=\tilde{l}^{*}$ as described previously.
From lemma 3.3.17 the involutive distributions on ( $0, \omega_{0}$ ) were constructed from a sequence of quotient algebras of the form


The above sequence (4.2.3) can now be used to rewrite the sequence of distributions as follows

$$
\Delta_{i}\left(\operatorname{Ad}_{\tilde{g}}^{*} f\right)=\Delta_{n}\left(\operatorname{Ad}_{\tilde{g}}^{*} f\right) \cap \bar{\Delta}_{i}\left(\operatorname{Ad}_{\tilde{g}}^{*} f\right) \quad \forall \tilde{g} \in \tilde{N}
$$

where $\bar{\Delta}_{i}\left(\operatorname{Ad}_{\hat{g}}^{*} f\right)$ is defined by $\left.d \sigma\left(\operatorname{Ad}_{\tilde{G}}\left(\tilde{S}^{n-i+1}\right)+\tilde{S}_{f}\right)\right)\left(\operatorname{Ad}_{\tilde{g}}^{*} f\right)$
Let $\tilde{X} \in \tilde{S}, \quad \tilde{Y} \in \hat{S}^{n-i+1}+\tilde{S}_{f}$
so $\tilde{Y}=\tilde{Y}_{1}+\tilde{Y}_{2}$ where $\tilde{Y}_{1} \in \tilde{S}^{n-i+1}, \tilde{Y}_{2}=\tilde{S}_{f}$

Now $\operatorname{do}\left(\operatorname{Ad}_{\tilde{g}} \tilde{\mathrm{Y}}_{2}\right)\left(\operatorname{Ad}_{\tilde{\mathrm{g}}}^{*} \mathrm{f}\right)=0$ by definition of
$\tilde{S}_{f}=\{\tilde{X} \subset \tilde{S}: d \sigma(\tilde{X})(f)(\tilde{Z})=0, \forall \tilde{Z} \in \tilde{S}\}$ alternatively $\tilde{\mathrm{Y}}_{2} \subset \tilde{\mathrm{~S}}_{\mathrm{E}}$ if
$\left\langle f,\left[\tilde{Y}_{2}, \tilde{Z}\right]\right\rangle=0, \dot{\tilde{Z}} \in \tilde{\mathrm{~S}}$

Thus $\left\langle\operatorname{Ad}_{\tilde{g}}^{*} \mathrm{f},\left[A d_{\tilde{g}} \tilde{\mathrm{Y}}_{2}, \tilde{z}\right]\right\rangle$
$=\left\langle\mathrm{f}, \quad\left[\tilde{\mathrm{Y}}_{2}, \mathrm{Ad}_{\tilde{\mathrm{g}}} \tilde{z}\right]\right\rangle=0$ by $(4.2 .4)$ so $\mathrm{Ad}_{\tilde{\mathrm{g}}} \tilde{\mathrm{Y}}_{2} \in \tilde{\mathrm{~S}}_{\mathrm{Ad}_{\tilde{\mathrm{g}}}^{\mathrm{t}}}$

Hence $\left[\operatorname{do}(\operatorname{Ad} \tilde{\mathbb{g}} \tilde{X}), \operatorname{do}\left(\operatorname{Ad}_{\tilde{\mathbf{g}}} \tilde{\mathrm{Y}}\right)\right] \quad\left(\operatorname{Ad}_{\tilde{\mathbf{E}}}^{*} \mathrm{f}\right)$
$=\left[d \sigma(\operatorname{Ad} \underset{\mathcal{G}}{ } \tilde{X}), \operatorname{d\sigma }\left(\operatorname{Ad}_{\tilde{g}} \tilde{\mathrm{Y}}_{1}\right)\right] \quad\left(\operatorname{Ad}_{\tilde{g}}^{*} \mathrm{f}\right)$
$=\operatorname{d\sigma }\left(\operatorname{Ad}_{g}\left[\tilde{X}, \tilde{Y}_{l}\right]\right)\left(\operatorname{Ad}_{g}^{*} f\right) \in \operatorname{do}\left(\operatorname{Ad}_{\bar{G}} \bar{S}^{n-i+2}\right)\left(\operatorname{Ad}_{g}^{*} f\right)$

$$
=\bar{\Delta}_{i-1}\left(\operatorname{Ad}_{\tilde{\mathbf{g}}}^{*} \mathrm{f}\right) \quad \forall \tilde{\mathrm{g}} \in \tilde{\mathbb{N}}
$$

Thus if $\tilde{Y} \subset \hat{h} \cap \tilde{\mathrm{~s}}^{\mathrm{n}-\mathrm{i}+1}+\tilde{\mathrm{S}}_{\mathrm{f}}$ then by $\hat{N}$ - invariance on $\Delta_{\mathrm{n}}$ and the above calculations

$$
\begin{align*}
& {\left[\tau_{*}^{-1} \mathrm{~d} \mathrm{\sigma}\left(\operatorname{Ad}_{\dot{\mathrm{g}}} \overline{\mathrm{X}}\right), \tau_{*}^{-1} \mathrm{~d} \mathrm{\sigma}\left(\operatorname{Ad}_{\tilde{\mathrm{g}}} \tilde{\mathrm{Y}}\right)\right]\left(\operatorname{Ad}_{\tilde{\mathrm{g}}}^{\mathrm{g}} \mathrm{f}\right)} \\
& =\tau_{*}^{-1}\left[\operatorname{do}\left(\operatorname{Ad}_{\bar{g}} \widehat{X}\right), \operatorname{d\sigma }\left(\operatorname{Ad}_{\bar{g}} \hat{Y}\right)\right]\left(\operatorname{Ad}_{\tilde{G}}^{*} f\right) \\
& \left.\varepsilon \quad{ }_{*}^{-1}\left(\Delta_{n} \cap \bar{\Delta}_{i-1}\right)\left(\operatorname{Ad}_{\hat{\delta}}^{*} f\right)\right) \quad \forall \tilde{g} \in \tilde{N} \tag{4.2.5}
\end{align*}
$$

So by construction there exists a basis of Hamiltonian vector fields

$$
\begin{aligned}
& ., \operatorname{do}\left(\operatorname{Ad}_{\hat{g}} \hat{\mathrm{Y}}_{r_{i}}^{\mathrm{i}}\right)\left(\operatorname{Ad}_{\hat{g}}^{*} \mathrm{f}\right), \quad 1 \leqslant \mathrm{i} \leqslant \mathrm{n}
\end{aligned}
$$

which $\operatorname{span} \Delta_{i}\left(\operatorname{Ad}_{\bar{g}}^{\star} f\right), \psi \tilde{g} \in \tilde{N}$ on $\left(0, \omega_{o}\right)$. Let $d \sigma\left(\bar{Y}_{1}\right), \ldots, d o\left(\bar{Y}_{m}\right)$ be the Hamiltonian vector fields, which exist by transitivity, that span the involutive distribution transversal to $\Delta_{n}$.

Notice that it is possible to construct canonical coordinates on $\left(0, \omega_{o}\right)$ similar to that on $(M, \omega)$ in theorem 3.3 .23 by simple application of the moment map $\tau$, i.e. the condition

$$
\left[x_{k}, x_{j}^{i}\right] \in \tau_{*}^{-1} \Delta_{i}
$$

was required, let $\tau_{*}^{-1}$ d $\sigma\left(\tilde{Y}_{k}\right)=X_{k}$ and $\tau_{\star}^{-1}$ do $\left(\tilde{Y}_{j}^{i}\right)=X_{j}^{i}$ then
$\left[d \sigma\left(\bar{Y}_{\mathrm{K}}\right), \operatorname{d\sigma }\left(\bar{Y}_{\mathrm{j}}^{\mathrm{i}}\right)\right] \epsilon \Delta_{i}$
and the re exists canonical coordinates

$$
\begin{aligned}
& \left(\tilde{q}_{1}, \ldots, \tilde{q}_{m}, \tilde{p}_{1}, \ldots, \tilde{p}_{m}\right) \text { on }\left(0, \omega_{0}\right) \text { such that } \\
& { }^{\partial} / \partial \tilde{p}_{1}, \ldots,{ }^{\partial} / \partial \tilde{p}_{i} \text { spans } \Delta_{i} \quad 1 \leqslant i \leqslant n
\end{aligned}
$$

where

$$
\begin{array}{ll}
q_{i}=\left(q_{1}^{i}, \ldots, q_{r_{i}}^{i}\right) & 1 \leqslant i \leqslant n \\
p_{i}=\left(p_{1}^{i}, \ldots, p_{r_{i}}^{i}\right) & 1 \leqslant i \leqslant n
\end{array}
$$

Observe that $\operatorname{do}(\operatorname{Ad} \underset{\mathbf{g}}{\hat{X}})\left(\operatorname{Ad}_{\tilde{g}}^{*} f\right)=\operatorname{Ad}_{\tilde{g}_{x}^{*}} \quad \operatorname{do}(\bar{X})(f)$
since

$$
\begin{aligned}
& \operatorname{do}\left(\operatorname{Ad}_{\tilde{g}} \tilde{X}\right)\left(\operatorname{Ad}_{\mathrm{g}}^{*} f\right)(\tilde{Z})
\end{aligned}
$$

$$
\begin{aligned}
& =-\left\langle\operatorname{Ad}_{\overline{\mathrm{g}}}^{\mathrm{E}}, \quad\left[\operatorname{Ad}_{\mathrm{g}} \tilde{X}, \tilde{Z}\right]\right\rangle \\
& =-\left\langle f, \operatorname{Ad}_{\bar{g}}^{-1} \quad\left[\operatorname{Ad}_{\bar{g}} \bar{X}, \bar{Z}\right]\right\rangle \\
& =-\left\langle f, \quad\left[\bar{X}, \operatorname{Ad}_{\dot{g}}^{-1} \tilde{z}\right]\right\rangle \\
& =\left\langle\operatorname{ad} \hat{X}^{\star} f, \text { Ad }_{\hat{g}-1} \bar{Z}\right\rangle \\
& =\operatorname{Ad}_{\bar{g}}^{*} \operatorname{d\sigma }(\bar{X})(f)(\tilde{Z})
\end{aligned}
$$

thus

$$
\begin{align*}
& {\left[\operatorname{d\sigma }\left(\operatorname{Ad}_{\bar{g}} \bar{X}\right)\left(\operatorname{Ad}_{\dot{g}}^{\star} f\right), \operatorname{d\sigma }\left(\operatorname{Ad}_{\bar{g}} \bar{Y}_{j}^{i}\right)\left(\operatorname{Ad}_{\dot{g}}^{\star} f\right)\right] } \\
= & \operatorname{Ad}_{\bar{g}}^{\star} *\left[\operatorname{d\sigma }(\bar{X}), \operatorname{do}\left(\bar{Y}_{j}^{i}\right)\right](f) \tag{4.2.6}
\end{align*}
$$

Also from the construction of the proof of theorem 3.2.8

$$
d \sigma\left(\bar{Y}_{j}^{i}\right)(f)=\partial / \partial \bar{P}_{j}^{i}
$$

and by invariance

$$
\begin{aligned}
& \operatorname{Ad}_{\bar{g}}^{\star} * \mathrm{da}\left(\mathrm{Y}_{\mathrm{j}}^{\mathrm{i}}\right)(\mathrm{f})={ }^{\partial} /{ }_{\partial \overline{\mathrm{p}}_{\mathrm{j}}}
\end{aligned}
$$

$$
\begin{aligned}
& =\partial / \partial \tilde{p}_{j}^{i}
\end{aligned}
$$

where $X_{q_{j}}$ is the Hamiltonian vector field of ${ }^{\partial} / \partial \hat{p}_{j}^{i}$.

So as da $(\bar{X}) \in \operatorname{Ham}\left(0, \omega_{0}\right)$, locally
$\left[\sum_{\ell=1}^{r_{k}} \sum_{k=1}^{n} \frac{\partial \hat{H}^{n}}{\partial \tilde{p}_{\ell}^{k}} \quad \partial / \partial \tilde{q}_{\ell}^{k}-\frac{\partial \hat{H}_{2}^{2}}{\partial \tilde{q}_{\ell}^{k}} / \partial \hat{p}_{\ell}^{k}, \quad \partial / \partial \tilde{p}_{j}^{i}\right]$
$\in \operatorname{span}\left\{\frac{\partial}{\partial} / \partial p_{i}^{2}, \ldots,{ }^{\partial} / \partial \vec{p}_{i-1}\right\}$
(by 4.2.5)

After straightforward computation exactly of the form of theorem 3.2.9 wi th

$$
\begin{gathered}
X_{\tilde{H}}=d \sigma(\hat{X}) \text { it is seen that } \\
\tilde{H}=\sum_{i=1}^{n} \tilde{p}_{i}^{\prime} \tilde{f}_{i}\left(\tilde{q}_{1}, \ldots, \bar{q}_{i-1}\right)+f_{0}\left(\tilde{q}_{1}, \ldots, \tilde{q}_{n}\right)
\end{gathered}
$$

similarly for $\tilde{H}_{u_{k}}$. Now using the moment map $\tau$, canonical coordinates $\left(q_{1}^{1}, \ldots, q_{r_{n}}^{n}, p_{1}^{1} \ldots, p_{r_{n}}^{n}\right) \quad$ can be found on $(M, \omega)$ wi th $H=\tau^{*} \widetilde{H}$ and $H_{u_{k}}=\tau^{*} \tilde{H}_{u_{k}}$
then

$$
H=\sum_{i=1}^{n} p_{i}^{\prime} f_{i}\left(q_{1}, \ldots, q_{i-1}\right)+f_{o}\left(q_{1}, \ldots, q_{n}\right)
$$

and also

$$
H_{u_{k}}=\sum_{i=1}^{n} p_{i}^{\prime} g_{i}^{k}\left(q_{1}, \ldots, q_{i-1}\right)+g_{o}^{k}\left(q_{1}, \ldots, q_{n}\right) \text { for } 1 \leqslant k \leqslant r
$$

as required.
Q.E.D.

Theorem 4.1.3 gave the following formula for calculating the Volterra series of a Hamiltonian system of the form (4.1.4) i.e.

$$
\begin{aligned}
& W_{n}\left(t, \sigma_{1}, \ldots, \sigma_{n}, x\right)=\left\{W_{n-1}\left(t, \sigma_{1}, \ldots, \sigma_{n-1}, x\right), h\left(\sigma_{n}\right)(x)\right\} \\
& W_{1}\left(t, \sigma_{1}, x\right)=\left\{h(t), h\left(\sigma_{1}\right)\right\}(x)
\end{aligned}
$$

where $h(\sigma)(x))=g_{0} \gamma_{f}(\sigma)(x)$
By use of the Campbell-Baker-Hausdorff formula and known identities of Hamiltonian vector fields on symplectic manifolds this can be rewritten in a form more amenable for calculation. The Campbell-Baker-Hausdorff formula is given by

$$
\gamma_{f}(-\sigma)_{*} X_{g}\left(\gamma_{f}(\sigma) x\right)=\sum_{i=0}^{\infty} \frac{\sigma^{i}}{i!} \operatorname{ad}_{X_{f}}^{i} X_{g}(x)
$$

where $\operatorname{ad}_{X_{f}}^{0} X_{g}=X_{g}, \quad \operatorname{ad}_{X_{f}}^{i} X_{g}=a d^{i-1} X_{f}\left[X_{f}, X_{g}\right]$ and $\gamma_{f}(t)$ is the flow of $X_{f}$. Also note,

$$
{ }^{a d} X_{f} X_{g}=\left[X_{f}, X_{g}\right]=X_{\{f, g\}}
$$

Now $\left\{\mathrm{h}(\mathrm{t})(\mathrm{x}), \mathrm{h}\left(\sigma_{1}\right)(\mathrm{x})\right\}$

$$
=\left\{g_{o} \gamma_{f}(t), g_{o} \gamma_{f}\left(\sigma_{1}\right)\right\}(x)
$$

$$
=\omega(x)\left(X_{g_{0} \gamma_{f}}(\tau) \quad(x), X_{g_{0} \gamma_{f}\left(\sigma_{1}\right)}(x)\right)
$$

$=\omega(x)\left(\gamma_{f}(-t)_{\star} X_{g}\left(\gamma_{f}(t) x\right), \gamma_{f}\left(-\sigma_{1}\right)_{\star} X_{g}\left(\gamma_{f}\left(\sigma_{1}\right) x\right)\right.$

$$
\begin{equation*}
=\sum_{i=0}^{\infty} \sum_{j=0}^{\infty} \frac{t^{j}{ }_{0}^{i}}{j!i!} \omega(x) \quad\left(a d_{X_{f}}^{j} X_{g}(x), \operatorname{ad}_{X_{f}}^{i} x_{g}(x)\right) \tag{4.2.7}
\end{equation*}
$$

The following detai led example of a strongly accessible, weakly observable linear analytic symplectic realization of a finite Volterra series shows the relevant computation.

Example 4.2.9 Consider the following system on $\left(\underset{\sim}{R},_{i=1}^{2} \sum_{i} \wedge q_{i}\right)$

$$
\begin{array}{rlr}
\dot{q}_{1}=u & q_{1}(0)=0 \\
\dot{q}_{2} & =q_{1}^{2}+q_{1}+u & q_{2}(0)=0 \\
-\dot{p}_{1} & =\left(2 q_{1}+1\right) p_{2}+q_{2} & p_{1}(0)=0 \\
-\dot{p}_{2} & =q_{1} & p_{2}(0)=0 \\
y(q, p)=p_{1}+p_{2} & \\
\text { Here } f & =\left(q_{1}^{2}+q_{1}\right) p_{2}+q_{1} q_{2} & \\
g & =p_{1}+p_{2} &
\end{array}
$$

and

$$
\begin{aligned}
& x_{f}=\left(q_{1}^{2}+q_{1}\right) \partial / \partial q_{2}-\left(2 q_{1}+1\right) p_{2} \partial / \partial p_{1}-q_{2} \partial / \partial p_{1}-q_{1} \partial / \partial p_{2} \\
& x_{g}=\partial / \partial q_{1}+\partial / \partial q_{2}
\end{aligned}
$$

This system is strongly accessible since

$$
\left[x_{g}, x_{f}\right]=\left(2 q_{1}+1\right) \frac{\partial}{\partial q_{2}}-2 p_{2} / \partial p_{1}-\partial / \partial p_{1}-\partial / \partial p_{2}
$$

$$
\left[x_{g},\left[X_{g}, x_{f}\right]\right]=2 \frac{\partial}{\partial q_{2}}
$$

$$
\left[x_{f},\left[x_{g}, x_{f}\right]\right]=2 q_{1} \frac{\partial}{\partial p_{1}}
$$

$$
\left[x_{g}, \quad\left[x_{f}\left[x_{g}, x_{f}\right]\right]\right]=2 \frac{\partial}{\partial p_{1}}
$$

and thus $S=\left\{\operatorname{ad}_{X_{f}}, X_{g}\right\}$ L.A. is such that $\operatorname{dim} S(x)=\operatorname{dim} T_{x}{\underset{R}{ }}^{4} \forall x \in{\underset{\sim}{R}}^{4}$ and $S$ is nilpotent.

It is weakly observable either by simple computation to check the algebraic conditions of Krener and Hermann [13] i.e.

$$
\begin{aligned}
& d g=d p_{1}+d p_{2} \\
& d L_{X_{f}} g=-d\left(\left(2 q_{1}+1\right) p_{2}+q_{2}+q_{1}\right) \\
& \\
& =-\left(2 q_{1} d p_{2}+2 p_{2} d q_{1}+d q_{2}+d q_{1}\right)
\end{aligned}
$$

$d L_{\left[X_{g}, X_{f}\right]} g=-d\left(2 p_{2}+2\right)$

$$
=-2 \mathrm{dp}_{2}
$$

$d L_{\left[X_{f}\left[X_{g}, X_{f}\right]\right]} g=d\left(2 q_{1}\right)$

$$
=2 \mathrm{dq} .
$$

and from this it is clear that $\operatorname{dim} d X(x)=\operatorname{dim} T_{x}^{*} R^{4} \quad \forall x \in R^{4}$ and thus the system is weakly observable. Alternatively, it is easy to show that a Hamiltonian system of the form (4.1.4) is strongly accessible if and only if it is weakly observable from the observation that if $h \in \mathcal{R}$

$$
X_{f} h=L_{X_{f}} h=\{f, h\}
$$

and $d L_{X_{f}} h=d\{f, h\}=\bar{\omega} \quad\left(\left[X_{f}, X_{h}\right]\right)$
where $\bar{\omega}: X \longmapsto w(x, \cdot)$. Thus if $\left\{\operatorname{ad}_{X_{f}}, X_{g}\right\}_{\text {L. }}$. spans so does $d \mathcal{X}$. For further details see Van der Schaft [31]. Moreover, if such a system is weakly observable by Crouch [5] it is observable.

The following computations of Poisson brackets are required

$$
\begin{aligned}
& \{f, g\}=\left(2 q_{1}+1\right) p_{2}+q_{2}+q_{1} \\
& \{g\{f, g\}\}=-2 p_{2}-1 \\
& \{f\{f, g\}\}=q_{1}\left(2 q_{1}+1\right)
\end{aligned}
$$

$$
\{\{f, g\},\{f\{f, g\}\}\}=0
$$

These are needed to calculate the Volterra series for this system.

$$
\begin{align*}
& W_{0}(t, x)=g\left(\gamma_{f}(t) x\right) \\
& W_{1}\left(t, \sigma_{1}, x\right)=\left\{h(t)(x), h\left(\sigma_{1}\right)(x)\right\} \tag{4.2.8}
\end{align*}
$$

where $h(\sigma)=g_{o} Y_{f}(\sigma)$
Using (4.2.7), (4.2.8) be comes

$$
\sum_{i=0}^{\infty} \sum_{j=0}^{\infty} \frac{t^{j} \sigma_{1}^{i}}{j!i!} \omega\left(a d_{X_{f}}^{j} X_{g}, a d_{X_{f}}^{i} X_{g}\right)(x)
$$

Now $\omega\left(X_{f}, X_{g}\right)=\{f, g\}$ and $\omega\left(X_{f},\left[X_{g}, X_{h}\right]\right)=\{f\{g, h\}\}$
so (4.2.8) $=\sigma_{1} \omega\left(X_{g}, a d_{X_{f}} X_{g}\right)(x)$

$$
+t \omega\left(a d_{X_{f}} X_{g}, X_{g}\right)(x)+\frac{t \sigma_{1}^{2}}{2} \omega\left(a d_{X_{f}} X_{G}, \operatorname{ad}_{X_{f}}^{2} X_{g}\right)(x)
$$

$$
+\frac{t^{2} \sigma_{1}}{2} \omega\left(a d_{X_{f}}^{2} X_{g}, a d_{X_{f}} X_{g}\right)(x)+\ldots
$$

$$
=\sigma_{1} \omega\left(X_{g}, a d_{X_{f}} X_{g}\right)(x)+t \omega\left(a d_{X_{f}} X_{g}, X_{g}\right)(x)
$$

(all other terms vanish)

$$
=-\sigma_{1}\left(2 p_{2}+1\right)+t\left(2 p_{2}+1\right)
$$

$$
=\left(t-\sigma_{1}\right)\left(2 p_{2}+1\right)
$$

By (4.2.8)

$$
W_{2}\left(t, \sigma_{1}, \sigma_{2}, x\right)=\left\{W_{1}\left(t, \sigma_{1}, x\right), h\left(\sigma_{2}\right)(x)\right\}
$$

$=\omega(x)\left(X_{W_{1}}\left(t, \sigma_{1}\right)(x), X_{h\left(\sigma_{2}\right)}(x)\right)$

Again using the Campbell-Baker-Hausdorff formula this can be rewritten as
$=\sum_{i=1}^{\infty} \frac{\sigma_{2}^{i}}{i!} \omega\left(X_{W_{1}}\left(t, \sigma_{1}\right) \quad, \quad a d_{X_{f}}^{i} X_{g}\right)(x)$
$X_{W_{1}}=2\left(t-\sigma_{1}\right)^{\partial} / \partial q_{2}$
$\left\{W_{1}, g\right\}=0$
$\left\{W_{1},\{f, g\}\right\}=2\left(t-o_{1}\right)$
$\left\{W_{1},\{f,\{f, g\}\}\right\}=0$
Computing (4.2.9)

$$
W_{2}\left(t, \sigma_{1}, \sigma_{2}, x\right)=2 \sigma_{2}\left(t-\sigma_{1}\right)
$$

Hence the Volterra series for this system corresponds to

$$
\begin{aligned}
y(t) & =\int_{0}^{t}\left(2 p_{2}+1\right)\left(t-\sigma_{1}\right) u\left(\sigma_{1}\right) d \sigma_{1} \\
& +\int_{0}^{t} \int_{0}^{\sigma_{1}} 2 \sigma_{2}\left(t-\sigma_{1}\right) u\left(\sigma_{2}\right) u\left(\sigma_{1}\right) d \sigma_{2} d \sigma_{1} .
\end{aligned}
$$

Notice that $W_{1}\left(t, \sigma_{1}, x\right)$ does depend on the state variable but $W_{2}\left(t, \sigma_{1}, \sigma_{2}, x\right)$ does not. This is as expected from lemna 4.2 .3 as $S$ is given by

$$
\begin{aligned}
& S=\left({ }^{\partial} / \partial q_{1}+\partial / \partial q_{2},\left(2 q_{1}+1\right)^{\partial} / \partial q_{2}-2 p{ }^{\partial} / \partial p_{1}-\partial / \partial p_{1}-\partial / \partial p_{2},\right. \\
& \\
& \left.\quad 2 q_{1} \quad \partial p_{1}\right\} \\
&
\end{aligned}
$$

which has step length 2 i.e. $a d_{X}^{3}=0 \forall X \in S$.

Further observe that $\Delta_{2}$ is spanned by the Hamiltonian vector fields $\left[X_{g}, X_{f}\right]$ and $\left[X_{g}\left[X_{f}\left[X_{g} X_{f}\right]\right]\right]$ and $\Delta_{1}$ is spanned by $\left[X_{g}\left[X_{f}\left[X_{g} X_{f}\right]\right]\right]$.

Corollary 4.2.10 With the same conditions as theorem 4.2.7 but with $S \notin \mathcal{L}$ and stationarity assumed then the canonical realization is of the form.

$$
\begin{array}{ll}
\dot{q}_{1}=A_{1} q_{1} & g_{1}^{i} \\
\dot{q}_{2}=A_{2} q_{2}+f_{2}\left(q_{1}\right) & g_{2}^{i}\left(q_{1}\right) \\
\vdots & +\sum_{i=1}^{r} u_{i} \\
\dot{q}_{n}=A_{n} q_{n}+f_{n}\left(q_{1}, \ldots, q_{n-1}\right) & g_{j}^{i}\left(q_{1}, \ldots, q_{j-1}\right) \\
\vdots & g_{n}\left(q_{1}, \ldots, q_{n-1}\right)
\end{array}
$$

$$
\begin{aligned}
& -\dot{p}_{1}=A_{1}^{\prime} p_{1}+\frac{\partial f_{2}^{\prime}}{\partial q_{1}} p_{2}+\ldots+\frac{\partial f_{n}^{\prime}}{\partial q_{1}} p_{n}+\frac{\partial f_{0}^{\prime}}{\partial q_{1}} \\
& +\sum_{i=1}^{r} u_{i} \quad \frac{\partial g_{2}^{i}}{\partial q_{1}^{i}} p_{1}+\ldots+\frac{\partial g_{n}^{i \prime}}{\partial q_{1}^{\prime}} p_{n}+\frac{\partial g_{0}^{i}}{\partial q_{1}^{\prime}} \\
& -\dot{p}_{n}=A_{n}^{\prime} p_{n}+\frac{\partial f_{0}^{\prime}}{\partial q_{n}}+\sum_{i=1}^{r} u_{1} \frac{\partial g_{0}^{i^{\prime}}}{\partial q_{n}} \\
& x=\left(q_{1}, \ldots, q_{n}, p_{1}, \ldots, p_{n}\right), \quad x(0)=x_{0} \\
& H=p_{1}^{\prime} A_{1} q_{1}+p_{2}^{\prime}\left(A_{2} q_{2}+f_{2}\left(q_{1}\right)\right)+\ldots+p_{n}^{\prime}\left(A_{n} q_{n}+f_{n}\left(q_{1}, \ldots, q_{n-1}\right)\right) \\
& +f_{0}\left(q_{1}, \ldots, q_{n}\right) \\
& g^{i}=H_{u_{i}}=p_{1}^{\prime} g_{1}^{i}+p_{2}^{\prime} g_{2}^{i}\left(q_{1}\right)+\ldots+p_{n}^{\prime} g_{n}^{i}\left(q_{1}, \ldots, q_{n-1}\right) \\
& +g_{0}^{i}\left(q_{1}, \ldots, q_{n}\right) \\
& 1 \leqslant i \leqslant r
\end{aligned}
$$

Proof. Follows on application of lemma 3.3 .25 by letting $F=\operatorname{expt} \hat{X}_{H}$ where $X_{H}$ is the free vector field on $(M, \omega), T\left(z_{o}\right)=f$ so Ad ${ }_{\operatorname{expt}}^{*} \tilde{X}_{H}=f$
follows Erom Ad* - equivariance as previously shown. Also as the commutar in this case for $F=\operatorname{expt} \bar{X}_{H}$ obviously vanishes, expt $\tilde{X}_{H}$ is a valid choice for an $F$ which satisfies the conditions of lemma 3.3.25. Furthemore,

$$
\operatorname{ad} \tilde{x}_{\mathrm{E}}: \tilde{s} \longrightarrow \tilde{s}
$$

and so clearly


Now $\tilde{S}_{f} c a n$ be interpreted as those vector fields which vanish at $f$, thus

$$
a \tilde{X}_{H}: \tilde{\mathrm{S}}_{\mathrm{F}} \longmapsto \tilde{\mathrm{~S}}_{\mathrm{E}}
$$

so


A sequence of distributions satisfying the conditions of theorem
3.2.8 can now be found exactly as in the previous case using the polarization of lemma 3.3.25.

As in theorem 4.3 of Crouch [5] there is a decomposition of
$G=V N$ where $V$ is the one parameter subgroup with generator $X_{H_{1}}$ such that $X_{H_{1}}\left(x_{0}\right)=0$ and $X_{H}=X_{H_{1}}+X_{H_{2}}, \quad X_{H_{2}} \in S$,

So as above

which is seen to be a linear endomorphism from Crouch [5].
Thus using the $s$ ame basis for $\Delta_{i}\left(A d_{g}^{\star} f\right)$ as in theorem 4.2.8., (4.2.11) gives that

$$
\begin{aligned}
& =d \sigma \operatorname{Ad}_{\bar{g}}\left(\left[\bar{X}_{\mathrm{H}_{1}}, \overline{\mathrm{Y}}_{\mathrm{j}}^{\mathrm{i}}\right]\right)\left(\operatorname{Ad}_{\mathrm{g}}^{\star} \mathrm{f}\right) \\
& =\operatorname{d\sigma }\left(\operatorname{Ad} \tilde{g}_{k=1}^{i}\left(\sum_{\ell=1}^{r_{k}} \tilde{a}_{k \ell} \dot{Y}_{\ell}^{k}\right)\right)\left(A{\underset{g}{g}}_{*}^{f}\right)
\end{aligned}
$$

Thus, with the introduction of canonical coordinates as in theorem 4.2.8 on ( $0, \omega_{0}$ ), this becomes
$\left[\sum_{t=1}^{n} \sum_{s=1}^{r_{t}} \frac{\partial \tilde{H}_{1}}{\partial \bar{p}_{s}^{t}} \quad \partial / \partial \hat{q}_{s}^{t}-\frac{\partial \tilde{H}_{1}}{\partial \tilde{q}_{s}^{t}} \quad \partial / \partial p_{s}^{2 t}, \partial / \partial \tilde{p}_{j}^{i}\right]$
$=\sum_{k=1}^{i} \sum_{\ell=1}^{r_{k}} \tilde{a}_{k \ell} \quad \frac{\partial}{\partial \tilde{p}_{\ell} k}$
$1 \leq i \leq n$
$1 \leqslant j \leqslant r_{i}$
$\tilde{a}_{k \ell}$ constants.
$\Rightarrow \tilde{H}_{1}=\tilde{p}, \tilde{A} \tilde{q}$
where $\tilde{q}=\left(\tilde{q}_{1}^{1}, \ldots, \tilde{q}_{r_{n}}^{n}\right) p=\left(\tilde{p}_{1}^{-1}, \ldots, \tilde{p}_{r_{n}}^{n}\right)$

As the other vector fields belong to $\widetilde{\mathrm{S}}$ their structure is determined by theorem 4.2.8. The required structure is obtained by observing that

$$
H=\tau^{*} \tilde{H}=\tau^{\star}\left(\tilde{H}_{1}+\tilde{H}_{2}\right)
$$

which means

$$
\begin{aligned}
& \quad H\left(q_{1}, \ldots, q_{n}, p_{1}, \ldots, p_{n}\right) \\
& =\sum_{i=1}^{n} p_{i}^{\prime}\left(A_{i} q_{i}+f_{i}\left(q_{1}, \ldots, q_{i-1}\right)+f_{0}\left(q_{1}, \ldots, q_{n}\right)\right. \\
& \text { with } f_{1}=0 \text { by gtationarity. }
\end{aligned}
$$

### 4.3. Interconnections

In this section the systems of 54.1 are shown to be intimately connected with systems that arise in the theory of electrical networks.

The results here depend on a system theoretical framework described by Willems [32] and Van der Schaft [33]. The basic motivation behind this approach is that it is often difficult to determine which space is actually the input space and which the output. Hence the following definition.

Definition 4.3.1 A differential system in state space form is described by:
(i) analytic manifolds X , W
(ii) an analytic bundle $\pi: B \longmapsto X$
(iii) an analytic map $f: B \mapsto T X \times W$ for which the diagram

commutes.

The system itself is defined by $\sum:=\{(x, w): \underset{\sim}{R} \longmapsto X \times W: x$ absolutely continuous and $(\hat{x}(t), W(t)) \in f\left(\pi^{-1}(x(t))\right.$ a.e.f and will be denoted by $\sum(X, W, B, f)$.

Obviously analyticity in the above definition may be replaced by a less stringent assumption but that is not necessary here.

To make such a system Hamiltonian the appropriate extra structure must be added to give,

Definition 4.3.2 A Hamiltonian system in state space form is described by:

commutes
(i) ( $M, W$ ) is a symplectic manifold of dimension $2 n$.
(ii) ( $\mathrm{W}, \omega^{e}$ ) is a symplectic manifold of dimension 2 m .

It can be seen, see for example Abrahan and Marsden [1], that TM $\times W$ can be made into a symplectic manifold by defining the symplectic form to be

$$
\Omega=\pi_{1}^{*} \dot{\omega}-\pi_{2}^{*} \omega^{e} \text { on } T M \times W
$$

where $\dot{\omega}$ is the induced symplectic form on the symplectic manifold $T M$ constructed from $M$, that is as $\omega$ is nondegenerate it defines a bundle isomorphism $\bar{\omega}$ from TM to $T^{*} M$ by setting

$$
\bar{\omega}(X)=i(X) \omega \text { for } X \subset T M .
$$

as $T^{*} M$ is symplectic and has a naturally defined 1 - form $\theta$ then $\bar{\omega}^{*} \theta$ is a 1 -form on $T M$ and let $\dot{\omega}=d \bar{\omega}^{*} \theta$ which is now clearly a symplectic forn on TM.

Definition 4.3.3 If $\sum(M, W, B, f)$ is a system with $M$ and $W$ symplectic manifolds then it is said to be (i) full Hamiltonian if $f(B)$ is a

Lagrangian submanifold of (TM $\times W, \Omega$ ) and (ii) degenerate Hamiltonian if there exists a full Hamiltonian system $\sum^{\prime}\left(M, W, B^{\prime}, f^{\prime}\right)$ such that $f(B)$ is a submanifold of $f^{\prime}\left(B^{t}\right)$.

Definition 4.3.4 Let $\left(W_{i}, \omega_{i}^{e}\right) i=1, \ldots, k$ be symp lectic manifolds. Then $\left(W_{1} x \ldots x W_{k}, \pi_{1}^{*} \omega_{1}^{e}+\ldots+\pi_{k}^{*} \omega_{k}^{e}\right)$ is a symplectic manifold. An interconnection of $\left(W_{i}\right)_{i=1, \ldots, k}$ is a submanifold $I \subset W_{1} x \ldots x W_{k}$. An interconnection is full (degenerate) Hamiltonian if $I$ is a Lagrangian (co-isotropic) submanifold of $\left(W_{1} x \ldots x W_{k}, \pi_{1}^{*} \omega_{1}^{e}+\ldots+\pi_{k}^{*} \omega_{k}^{e}\right)$.

After a number of restrictions are placed upon the above definition it is shown in Van der Shaft [33] that Kirchhoff's laws are a special type of interconnection. Also in Van der Schaft [33] it is shown that the following theorem holds.

Theorem 4.3.5 Let $\sum_{i}\left(M_{i}, W_{i}, B_{i}, f_{i}\right) \quad i=1, \ldots, k$ be Hamiltonian systems interconnected by a Hamiltonian interconnection $I \subset W_{1} \times \ldots x W_{k}$. Then the resulting system is a Hamiltonian system

$$
\sum_{I}\left(M_{1} x \ldots \times M_{k}, W_{1} x_{\ldots} \ldots x W_{k}, B_{I}, f_{I}\right) .
$$

Proof. The product system is constructed given by
$\sum\left(M_{1} x_{1} \ldots M_{k}, W_{1} \times \ldots x W_{k}, B_{x}, f_{x}\right)$ of the systems $\sum_{i}\left(M_{i}, W_{i}, B_{i}, f_{i}\right)$, $i=1, \ldots, k$. Let $x=\left(x_{1}, \ldots, x_{k}\right) \in M_{1} x \ldots M_{k}$. Since $\pi_{i}: B_{i} \longmapsto M_{i}$ are fibre bundles there exists neighbourhoods $U_{i} \subset M_{i}$ of $x_{i}$ such that $\pi_{i}^{-1}\left(U_{i}\right) \cong U_{i} \times F_{i}, F_{i}$ is the standard fibre, $B_{x}$ is now defined locally as $\pi_{x}:\left(U_{1} \times \ldots \times U_{k}\right) \times\left(F_{1} \times \ldots \times F_{k}\right) \mapsto U_{1} \times \ldots \times U_{k}$.

Next define $f_{x}: B_{x} \longmapsto T\left(M_{1} \times \ldots x M_{k}\right) \times W_{1} X \ldots \times W_{k}$ locally as $f_{x}=\left(f_{1}, \ldots, f_{k}\right)$. Obviously if $\sum_{i}\left(M_{i}, W_{i}, B_{i}, f_{i}\right)$ are Hamiltonian then [ $\left(M_{1} \times \ldots x M_{k}, W_{1} x \ldots x W_{k}, B_{x}, f_{x}\right)$ is Hamiltonian either full or degenerate. Also the interconnected system must satisfy $f_{I}\left(B_{I}\right) \subset f_{x}\left(B_{x}\right)$. Thus definition 4.2 .3 gives result.
Q.E.D.

The main result of this section may now be given.

Theorem 4.3.6 A stationary linear analytic symplectic realization of a stationary finite Volterra series can be decomposed into a set of linear Hamiltonian systems with a full Hamiltonian interconnection.

Proof. Theorem 4.1.7 gives the conditions on the Volterra series for a stationary linear analytic symplectic realization to exist. Furthermore it is clear from the constructions given in the proof of theorem 4.1.7 that a Volterra series has a stationary linear analytic symplectic realization if and only if it has a stationary bilinear symplectic realization. This bilinear realization is in the form of system (4.1.11) i.e.

.

都

$$
b=\left[\begin{array}{l}
b_{1} \\
0 \\
0 \\
0
\end{array}\right], \quad c=\left[\begin{array}{c}
0 \\
\cdot \\
0 \\
0 \\
c_{n}
\end{array}\right], \quad q=\left[\begin{array}{c}
q_{1} \\
\cdot \\
\cdot \\
q_{n} \\
q_{n}
\end{array}\right], \quad p=\left[\begin{array}{l}
p_{1} \\
\cdot \\
\cdot \\
p_{n}
\end{array}\right]
$$

This can be decomposed into $n$ linear Hamiltonian systems as follows,

$$
\text { (2) } \begin{aligned}
\dot{q}_{2} & =A_{2} q_{2}+b_{2} u_{1}^{2}, q_{2}(0) \\
-\dot{p}_{2} & =A_{2}^{\prime} p_{2}+c_{2} u_{2}^{2}, p_{2}(0)
\end{aligned}=0,\left[\begin{array}{l}
y_{1}^{2} \\
y_{2}^{2}
\end{array}\right]=\left[\begin{array}{ll}
0 & b_{2}^{\prime} \\
c_{2}^{\prime} & 0
\end{array}\right]\left[\begin{array}{l}
q_{2} \\
p_{2}
\end{array}\right]
$$

with interconnection I satisfying

$$
\begin{aligned}
& u_{1}^{1}=u \\
& u_{2}^{1}=u \quad y_{1}^{2} \\
& u_{1}^{2}=u \quad y_{2}^{1} \\
& u_{2}^{2}=u \quad y_{1}^{3} \\
& :
\end{aligned}
$$

$$
\begin{gathered}
: \\
u_{1}^{i}=u \quad y_{2}^{i-1} \\
u_{2}^{i}=u \quad y_{i}^{i+1} \\
\vdots \\
\vdots \\
u_{1}^{n}=u \quad y_{2}^{n-1} \\
u_{2}^{n}=u
\end{gathered}
$$

with the usual symplectic form on $W_{1} X_{\text {. . . }} \mathrm{xW}_{\mathrm{n}}$ given locally by $d y_{1}^{1} \wedge d u_{1}^{1}+d y_{2}^{1} \cdot d u_{2}^{1}+\ldots+d y_{1}^{n} \cdot d u_{1}^{n}+d y_{2}^{n} \cdot d u_{2}^{n}$. So on I this form is, with the usual assumption that $u \in \ell_{p, c .}$.
$u d y_{1}^{2} \cdot d y_{2}^{1}+u d y_{2}^{1} \cdot d y_{1}^{2}+u d y_{1}^{2} \cdot d y_{2}^{2}+u d y_{2}^{2} \cdot d y_{1}^{3}+\ldots+u d y_{1}^{n} \cdot d y_{2}^{n-1}$ $+\mathrm{udy}{ }_{2}^{\mathrm{n}-1} \cdot d y_{1}^{\mathrm{n}}=0$.

Hence $I$ is a Lagrangian submanifold since it is described by 2 N constraints on a 4 N - dimensional symplectic manifold and hence a full Hamiltonian interconnection.
Q.E.D.


A similar decomposition procedure may be obtained for a linear analytic Hamiltonian realization of the form appearing in Corollary 4.2.9. For simplicity assume only one control then this system may be written as n - linear Hamiltonian Systems of the form

$$
\text { (1) } \begin{aligned}
\dot{q}_{1} & =A_{1} q_{1}+u_{1}^{1}, q_{1}(0)
\end{aligned}=q_{1}^{o}, \quad\left[\begin{array}{l}
y_{1} \\
-\dot{p}_{1}
\end{array}\right] A_{1}^{\prime} p_{1}+u_{2}^{1}, p_{1}(0)=p_{1}^{o}, \quad\left[\begin{array}{ll}
0 & I_{r_{1}} \\
y_{2}
\end{array}\right] \quad\left[\begin{array}{l}
q_{1} \\
I_{r_{1}} \\
0
\end{array}\right]
$$

(2) $\dot{q}_{2}=A_{2} q_{2}+u_{1}^{2}, q_{2}(0)=q_{2}^{0}$,

$$
-\dot{p}_{2}=A_{2}^{\prime} p_{2}+u_{2}^{2}, p_{2}(0)=p_{2}^{0}
$$

$$
\left[\begin{array}{l}
y_{1}^{2} \\
y_{2}^{2}
\end{array}\right]=\left[\begin{array}{ll}
0 & I_{r_{2}} \\
I_{r_{2}} & 0
\end{array}\right]
$$

$$
\left[\begin{array}{l}
q_{2} \\
p_{2}
\end{array}\right]
$$

where $I_{r_{i}}$ is the $r_{i} \times r_{i}$ identity matrix with interconnection $I$ satisfying
$k u_{1}^{1}=\partial V / \partial k_{y_{1}}^{1}, \forall k \in\left\{1, \ldots, r_{1}, 1, \ldots, r_{2}, \ldots, 1, \ldots, r_{n}\right\}$
$k_{u_{2}}^{1}=\partial V / \partial k_{2}{ }_{2}$
$k u_{1}^{2}=\partial v / \partial k_{1}{ }_{1}^{2}$
$k u_{2}^{2}=\partial v / \partial k y_{2}^{2}$
$k_{u_{1}^{n}}^{n}=\partial V / \partial k_{1}^{n}$
$k_{u_{2}}^{n}=\partial V / \partial k_{y_{2}}^{n}$
where $u{ }_{j}^{i}=\left({ }^{1} u_{j}^{i}, \ldots,{ }^{\mathbf{r}}{ }_{u}{ }_{j}^{i}\right), \quad 1 \leqslant i \leqslant n$
$1 \leqslant j \leqslant 2$

$$
y_{j}^{i}=\left({ }^{1} y_{j}^{i}, \ldots{ }^{r_{i}} y_{j}^{i}\right)
$$

and
$V=y_{1}^{2} \quad f_{2}\left(y_{2}^{1}\right)+\ldots+y_{1}^{n} f_{n}\left(y_{2}^{1}, \ldots, y_{2}^{n-1}\right)+f_{o}\left(y_{2}^{1}, \ldots, y_{2}^{n}\right)$

$$
+u\left(y_{1}^{1} g_{1}+y_{1}^{2} g_{2}\left(y_{2}^{1}\right)+\ldots+y_{1}^{n} g_{n}\left(y_{2}^{1}, \ldots, y_{2}^{n-1}\right)+g_{0}\left(y_{2}^{1}, \ldots, y_{2}^{n}\right)\right)
$$

The symplectic form on $W_{1} x \ldots x W_{n}$ is given locally by,

$$
\sum_{\ell=1}^{n} \sum_{k=1}^{r_{\ell}} d^{k} y_{1}^{\ell} \wedge d^{k} u_{1}^{\ell}+d^{k} y_{2}^{\ell} \wedge d^{k} u_{2}^{\ell}
$$

and on this form I vanishes since

$$
\begin{aligned}
& \sum_{\ell=1}^{n} \sum_{k=1}^{r_{\ell}} d^{k} y_{1}^{\ell} \wedge d\left(\frac{\partial V}{\partial^{k} y_{1}^{\ell}}\right)+d^{k} y_{2}^{\ell} \cdot d\left(\frac{\partial V}{\partial^{k} y_{2}^{\ell}}\right) \\
& =\sum_{\ell=1}^{n} \sum_{k=1}^{r_{\ell}} \sum_{s=1}^{n} \sum_{j=1}^{r}\left(\frac{\partial^{2}}{\partial^{k} \frac{v}{\ell}} \partial_{1} j_{y} s \quad d_{1}^{k} y_{1}^{\ell} \sim d^{j} y_{1}^{s}+\frac{\partial^{2} v}{\partial^{k} y_{l}^{\ell}} \partial_{1} j_{y_{2}} d^{k} y_{1}^{\ell} \sim d^{j} y_{2}^{s}\right) \\
& +\sum_{\ell=1}^{n} \sum_{k=1}^{r_{\ell}} \sum_{s=1}^{n} \sum_{j=1}^{\mathbf{r}}\left(\frac{\partial^{2} v}{\partial^{k} y_{2}^{\ell} \partial^{2}} j_{y_{1}} \quad d^{k} y_{2}^{\ell} \wedge d^{j} y_{1}^{s}\right. \\
& \left.+\frac{\partial^{2} v}{\partial^{k} y_{2}^{\chi}} \partial^{j} y_{2}^{s} \quad d^{k} y_{2}^{\ell} \wedge d^{j} y_{2}^{s}\right)
\end{aligned}
$$

is clearly zero as the wedge product (, ) is anti-symmetric.

Furthermore, as $W_{i}={\underset{R}{ }}^{4} r_{i}, 1 \leqslant i \leqslant n$ and $I$ is described by
$2 \sum_{j=1}^{n} r_{j}=\frac{1}{2}$ dim $\left(W_{1} x \ldots W_{n}\right)$ constraints $I$ is a Lagrangian submanifold and thus, by definition, a full Hamiltonian interconnection.

## Chapter 5

## COMPLETE INTEGRABILITY OF NILPOTENT SYSTEMS

In this chapter a generalization of the concept of a transitive abelian group action on a Lagrangian submanifold of a symplectic manifold as in Liouville's theorem [2] is found by considering a nilpotent group action which is the first stage to a control theoretic version of Liouville's theorem. It is shown that such systems can also be integrated by quadratures.

### 5.1 Integrability of Hamiltonian Systems

In this section a feasible Eramework is constructed in order to look at families of Hamilconi an vector fields on a symplectic manifold ( $M, w$ ) which have a nilpotent Lie algebra structure. This, in some sense, extends the classical case of complete integrability as given by Liouville's theorem where, essentially, abelian Lie algebras are of incerest.

First, however, the symplectic vector space case is investigated.

Let $(V, w)$ be a $2 n$ - dimensional symplectic vector space and
$\mathcal{L} \subset \mathrm{gl}(\mathrm{V})$ a Lie $s u b a l g e b r a$ of Hamiltonian endomorphisms. Assume
$U_{1}=\left\{v \in V: H(v)=0\right.$ for all $\left.X_{H} \in \mathcal{L}\right\}$ is an isotropic subspace of $V$, and $\bar{W}$ is a Lagrangian subspace of $V$ transversal to the level surface $U_{1}$ i.e. $\bar{W} \cap U_{1}=\{0\} . U_{1}$ can be extended to a Lagrangian subspace, $W$, as follows. If $U_{1}$ is not already Lagrangian then

$$
\mathrm{U}_{1} \subsetneq_{\neq} \mathrm{u}_{1}^{\perp}
$$

It is claimed that $U^{\perp} \subset U_{1}+\bar{W}$ for if $U_{1}^{\perp} \subset U_{1}+\bar{W} \Rightarrow U_{1}^{\perp} \cap \bar{W}^{\perp} \subset U_{1}$ $\Rightarrow U_{1}^{\perp} \cap \bar{W} \subset U_{1}$, so $U_{1} \cap \bar{W}=\{0\} \Rightarrow U_{1}^{\perp} \cap \bar{W}=\{0\} \Rightarrow \operatorname{dim}_{1}^{\perp} \leq n$ $\Rightarrow \operatorname{dim} U_{1} \geqslant n$
which is a contradiction.
So there exists a $v \in U_{1}^{\perp}, v \notin U_{1}+\bar{W} . \quad$ Let $U_{2}=U_{1}+R v$ then $\operatorname{dim} U_{2}=\operatorname{dim} U_{1}+1$. All that is now required is to prove that $U_{2}$ is isotropic. So suppose
$w \in U_{2} \cap \bar{W} \Rightarrow w=v_{1}+\lambda v$ with $v_{1} \in U_{1} \Rightarrow \lambda v=w-v_{1} \subset U_{1}+\bar{W}$
but $v \not \& U_{1}+\bar{W}$ so $\lambda=0$, this implies that $w \in U_{1} \cap \bar{W}=\{0\} \Rightarrow w=0$. In fact this automatically proves that $\bar{W}$ is transversal to $U_{2}$ as well as $\mathrm{U}_{2}$ being isotropic.

Continue by induction until $U_{m}$ is reached with $\operatorname{dim} U_{m}=n=\frac{1}{2} \operatorname{dim} V$. Let $W=U_{m}$ which is a Lagrangian subspace.

Proposition 5.1.1 Let ( $V, \omega$ ) be a $2 n$ - dimensional symplectic vector space with $W$ and $\bar{W}$ as above. Suppose $X_{H} \cdot \bar{W} \subset \bar{W}$ and $X_{H} \cdot W \subset W$, $\forall X_{H} \in \mathcal{L}$ and that the restriction $X_{i}$ of ${X_{H_{i}}}, \forall X_{H_{i}} \in \mathcal{L}$, to $W$ is a nilpotent endomorphism on $W$. Then there exists a symplectic basis of $V$ such that in canonical coordinates relative to this basis each $X_{H} \in \mathcal{L}$ can be expressed by

$$
\mathrm{X}_{\mathrm{H}}=\left(\begin{array}{cc}
\mathrm{N} & 0 \\
0 & -\mathrm{N}
\end{array}\right)
$$

where $N \subset{\underset{\sim}{R}}^{\mathrm{n} \times n}$ is strictly lower triangular.

Proof. By Engel's theorem 3.1.2 there exists an ordered basis for W, say $e_{q_{1}}, \ldots, e_{q_{n}}$ such that in coordinates relative to this basis each $X_{i}$ can be expressed by a strictly lower triangular matrix. Then, as in theorem 3.1.11 a complementary basis may be found on $\bar{W}$, say $e_{p_{1}}, \ldots, e_{p_{n}}$ such that ( $e_{q_{1}}, \ldots, e_{q_{n}}, e_{p_{1}}, \ldots, e_{p_{n}}$ ) is a symplectic basis on $V$. Let $\left(q_{1}, \ldots, q_{n}, p_{1}, \ldots, p_{n}\right)$ be canonical coordinates relative to this basis then each $X_{H}$ takes the form

$$
X_{H}=\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right) \quad \begin{aligned}
& a, b, c, d=R^{n \times n} \\
& a=-d^{\prime}, b=b^{\prime}, c=c^{\prime}
\end{aligned}
$$

as each $X_{H} \in \mathcal{L}$ is by assumption infinitesimally syrmplectic.

$$
\text { Let }(0, \bar{w}) ' \in \bar{W} \text {, then since } X_{H} \cdot \bar{W} \in \bar{W}, \forall x_{H} \in L
$$

$$
X_{H} \cdot(0, \bar{w})^{\prime}=\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right)\binom{0}{\bar{w}}=\binom{b \bar{w}}{d \bar{w}} \subset \bar{w}
$$

which implies $b \equiv 0$. Thus

$\left.\begin{array}{c}0 \\ -a^{1}\end{array}\right)$
and so $H=p^{\prime} a q+\frac{1}{2} q^{\prime} c q$.
But $H=0$ on $W$, so $c=0$ on $W$, as $W$ is described by $p_{1}=\ldots=p_{n}=0$, and thus on $V$. Hence, $X_{H}=\left(\begin{array}{cc}a & 0 \\ 0 & -a^{\prime}\end{array}\right)$

But $X_{H}$ restricted to $W$ is strictly lower triangular, thus $a=N$ which gives the result.
Q.E.D.

The following example introduces the concept of transitivity of a nilpotent Lie group on a Lagrangian subspace of a symplectic vector space as well as providing an example of proposition 5.1.1.

Example 5.1.2 Consider the following bilinear Hamiltonian system on ( $\mathrm{R}^{4}, \mathrm{~J}$ ) with canonical coordinates $\left(\mathrm{q}_{1}, \mathrm{q}_{2}, \mathrm{P}_{1}, \mathrm{p}_{2}\right.$ )

$$
\dot{x}=X_{H}+u X_{H_{u}}, \quad x(0)=0, \quad x \quad\left(R^{4}, J\right) \text { where } x=\left(q_{1}, q_{2}, p_{1}, p_{2}\right)
$$ and $H=P_{2} q_{1}, H_{u}=p_{2} q_{1}+p_{1}$, This can be expanded as,

$$
\begin{align*}
& \dot{\mathrm{q}}_{1}=\mathrm{u} \\
& \dot{\mathrm{q}}_{2}=\mathrm{q}_{1}+\mathrm{u} \mathrm{q}_{1} \\
&-\mathrm{q}_{1}(0)=0  \tag{5.1.1}\\
& \dot{\mathrm{p}}_{1}=\mathrm{p}_{2}+\mathrm{u} \mathrm{p}_{2}
\end{align*} \mathrm{q}_{2}(0)=0
$$

Now consider the level surface given by

$$
L=\left\{v \in{\underset{\sim}{R}}^{4}: H_{i}(v)=0, \forall X_{H_{i}} \in \mathcal{L}=\left\{X_{H}, X_{H_{u}}\right\}_{L, A .}\right\} .
$$

Since $\left[\mathrm{X}_{\mathrm{H}_{\mathrm{u}}}, \mathrm{X}_{\mathrm{H}}\right]=\frac{\partial}{\partial \mathrm{q}_{2}}=\mathrm{X}_{\mathrm{P}_{2}}$ and all other brackets vanish, L may be described by $p_{1}=p_{2}=0$. Thus $L$ is a Lagrangian subspace of ${\underset{\sim}{R}}^{4}$ and further the restriction of the vector fieldsof $\mathcal{L}_{\text {to }} L$ span $L$ everywhere and are nilpotent.

Furthermore, by introducing, the new coordinates $q_{0}, p_{0}$, (5.1.1) can be rewritten as in chapter 4.2 as,
$\left[\begin{array}{l}\dot{q}_{0} \\ \dot{q}_{1} \\ \dot{q}_{2} \\ \dot{p}_{0} \\ \dot{p}_{1} \\ \dot{p}_{2}\end{array}\right]=\left[\begin{array}{llllll}0 & 0 & 0 & 0 & 0 & 0 \\ 0 & 0 & 0 & 0 & 0 & 0 \\ 0 & 1 & 0 & 0 & 0 & 0 \\ 0 & 0 & 0 & 0 & 0 & 0 \\ 0 & 0 & 0 & 0 & 0 & -1 \\ q_{2} \\ q_{1} \\ p_{0} \\ q_{0} \\ p_{1} & 0 & 0 & 0 & 0 & 0 \\ 0 & 1 & 0 & 0 & 0 & 0 \\ 0 & 0 & 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 0 & 0 & 0 \\ 0 & 0 & 0 & 0 & 0 & -1 \\ q_{1} \\ q_{2} \\ p_{2}\end{array}\right]\left[\begin{array}{l}q_{0} \\ q_{0} \\ p_{0} \\ p_{1} \\ p_{1} \\ 0\end{array}\right]$

on ( $\underline{R}^{6}, J^{e}$ ) with canonical coordinates ( $\left.q_{0}, q_{1}, q_{2}, p_{0}, p_{1}, p_{2}\right)$, $J^{e}=\left(\begin{array}{cc}0 & I_{3} \\ -I_{3} & 0\end{array}\right)$ and initial conditions $q_{0}(0)=1, q_{1}(0)=q_{2}(0)=0$, $p_{0}(0)=p^{0}, p_{1}(0)=p_{2}(0)=0$.

Now it is necessary to pass to the extended level surface $U_{1}$ on ${\underset{\sim}{R}}^{6}$ of $L$ on $\underline{R}^{4}$ in order to make a comparison with proposition 5.1.1 where

$$
U_{1}=\left\{v \in \mathbb{R}^{6}: H_{i}^{e}(v)=0, \forall X_{H_{i}^{e}} \in \mathcal{L}^{e}=\left\{X_{H}, X_{H_{u}}^{e}\right\}_{L, A .}\right\}
$$

Now $H^{e}=P_{2} q_{1}$

$$
H_{u}^{e}=p_{2} q_{1}+p_{1} q_{o}
$$

and $\left\{H_{u}^{e}, H^{e}\right\}=-P_{2} q_{0}$ with all other Poisson brackets vanishing. So $U_{1}$ is an isotropic subspace described by $p_{1}=p_{2}=0$. In this case the obvious extension is a Lagrangian subspace W is

$$
w=\left\{v \in \underline{R}^{6}: H_{i}^{e}(v)=0, \forall X_{H_{i}}^{e} \in L^{e} \text { and } p_{0}=0\right\}
$$

which may be described by $p_{0}=p_{1}=p_{2}=0$.
Obviously theorem 5.1.1 can be proved with nilpotent replaced by solvable and assuming the field over the Lie subalgebra $\mathcal{L} \subset g^{\ell}$ (V) is algebraically closed, then by using Lie's theorem 3.1.14 each element of [ has the matrix representation

$$
x_{H}=\left(\begin{array}{cc}
S & 0 \\
0 & -S^{\prime}
\end{array}\right) \quad \text { where } s \in \underline{R}^{\mathrm{nxn}} \text { is lower triangular. }
$$

The next part of this section considers a globalization of the above.
Let $\left\{\mathrm{X}_{\mathrm{H}_{\mathrm{o}}}, \mathrm{X}_{\mathrm{H}_{1}}, \ldots, \mathrm{X}_{\mathrm{H}_{\mathrm{m}}}\right.$ \} be a family of Hamiltonian vector fields on a $2 n$ - dimensional symplectic manifold $(M, \omega)$. Denote by $L$ the Lie algebra generated by these vector fields i.e. $\mathcal{L}=\left\{X_{H_{0}}, X_{H_{1}}, \ldots, X_{H_{m}}\right\}_{\text {Li. }}$. and suppose its corresponding lie group is given by G.

Let exp $X_{H}$ represent the flow of the Hamiltonian vector field $X_{H}$ on $M$, $\forall x_{H} \in \mathcal{L}$. Define a group action on $M$ as follows

$$
\oplus: G \times M \longmapsto M
$$

$\Phi \quad\left(\operatorname{expt} X_{1} X_{1} \quad \ldots \ldots \exp t_{k} X_{H_{k}}, x\right)$
$=\operatorname{expt}{ }_{1} X_{H_{1}} 0 \ldots 0 \exp _{k} X_{H_{k}} \cdot x$

Consider the level surface

$$
N=\left\{x \in M: H(x)=0, \quad \forall X_{H} \in \mathcal{R}\right\}
$$

For the rest of this section $N$ is assumed to be a Lagrangian submanifold of $M$.

By Kostant's corollary 3.2.5 in a neighbourhood $U$ of $N$ in $M, M$ can be identified with $T^{*} N$. As the results in the remainder of this section are local, without loss of generality, $M$ is assumed to be $T^{*} N$.

Let ( $\mathrm{q}_{1}, \ldots, q_{\mathrm{n}}, \mathrm{p}_{\mathrm{I}}, \ldots, \mathrm{p}_{\mathrm{n}}$ ) be canonical coordinates on $\mathrm{T}^{*} \mathrm{~N}$. $N$ can be thought of as being the zero section and thus defined by $p_{1}=\ldots=p_{n}=0$.

Let $\Delta$ be a polarization in $\mathrm{T}^{\star} \mathrm{N}$ transverse to the Lagrangian submanifold N. By proposition 3.2.7 locally $\Delta$ is spanned by $\partial / \partial p_{1}, \ldots ., \partial p_{n}$.

Therefore the corresponding condition here to that of $X_{H} \cdot \bar{W} \subset \bar{W}$ in theorem 5.1.1 is

$$
\begin{equation*}
\left[\mathrm{x}_{\mathrm{H}}, \Delta\right] \subset \Delta, \quad \forall \mathrm{x}_{\mathrm{H}} \in L \tag{5.1.2}
\end{equation*}
$$

Proposition 5.1.3 Each Hamiltonian function, H, of a Hamiltonian vector field on the symplectic manifold ( $T^{\star} N, \omega_{N}$ ) satisfying (5.1.2), locally has the form $H=\sum_{i=1}^{n} p_{i} f_{i}\left(q_{1}, \ldots, q_{n}\right)+f_{0}\left(q_{1}, \ldots, q_{n}\right)$

Proof. It is sufficient to perform the following calculation with respect to the canonical coordinates ( $q_{1}, \ldots, q_{n}, p_{1}, \ldots, p_{n}$ )

$$
\begin{aligned}
& {\left[X_{H}, f(q, p)^{\partial} / \partial p_{j}\right]=\left[\sum_{i=1}^{n} \frac{\partial H}{\partial p_{i}} \partial / \partial q_{i}-\frac{\partial H}{\partial q_{i}} \quad \partial / \partial p_{i}, f(q, p)^{\partial} / \partial p_{j}\right]} \\
& =\sum_{i=1}^{n} f \frac{\partial^{2} H}{\partial p_{j} \partial q_{i}} \quad \frac{\partial}{\partial p_{i}}-f \frac{\partial^{2} H}{\partial p_{j} \partial p_{i}} \quad \frac{\partial}{\partial q_{i}}+\frac{\partial H}{\partial p_{i}} \frac{\partial f}{\partial q_{i}} \quad \partial / \partial p_{j}-\frac{\partial H}{\partial q_{i}} \frac{\partial f}{\partial p_{i}} \partial / \partial p_{j} \\
& \subset \Delta \text { by (5.1.2) } \quad 1 \leqslant j \leqslant n
\end{aligned}
$$

This implies $\frac{\partial^{2} H}{\partial p_{i}} \partial p_{j}=0, \quad \psi \quad 1 \leqslant i_{i} \quad j \leqslant n$
$\Rightarrow H$ is linear in each $P_{i}, 1 \leqslant i \leqslant n$.
$\Rightarrow H=\sum_{i=1}^{n} p_{i} f_{i}\left(q_{1}, \ldots, q_{n}\right)+f_{o}\left(q_{1}, \ldots, q_{n}\right)$

This structure does appear naturally as the following proposition found in Abraham and Mars den [1] shows.

Proposition 5.1.4 Let expt $X_{H}$ be the flow of a Hamiltonian vector field $X_{H}$ on a symplectic manifold $\left(T^{*} N, \omega_{N}\right)$, where $\omega_{N}=\sum_{i=1}^{n} d q_{i}$. dp $i_{i}$ locally. Assume that expt $X_{H}: T^{\star} N \longmapsto T^{\hbar^{\prime}} N$ preserves the fibres of $T^{*} N$, then locally the Hamiltonian function, $H$, is given by

$$
H=\sum_{i=1}^{n} p_{i} f_{i}\left(q_{1}, \ldots, q_{n}\right)+f_{o}\left(q_{1}, \ldots, q_{n}\right)
$$

However, by assumption each $H=0$ on $N$ i.e. on $p_{1} \ldots \ldots=p_{n}=0$ which implies $f_{0}\left(q_{1}, \ldots, q_{n}\right)=0$ on $N$, and since it only depends on the $q_{i}$ coordinates, on $T^{*} N$. So each $H=p^{\prime} f(q)$ where $q=\left(q_{1}, \ldots, q_{n}\right)^{\prime}$, $p=\left(p_{1}, \ldots, p_{n}\right)^{\prime}$ for all $X_{H} \in \mathcal{L}$.

Again these structures appear quite naturally as the next proposition from Abraham and Mars den [1] shows.

Proposition 5.1.5 Suppose exp $X_{H}$ is the flow of a Hamiltonian vector field $X_{H}$ on the symplectic manifold ( $T^{\star} N, w_{N}$ ), and further assume that it preserves the natural $1-$ form $\theta_{N}, \omega_{N}=-\mathrm{d} \theta_{\mathrm{N}}$. Then, locally, H is of the form $H=\sum_{i=1}^{n} p_{i} f_{i}\left(q_{1}, \ldots, q_{n}\right)$.

The following structure theorem on nilpotent Lie algebras of Matsushima [34] is of importance.

Theorem 5.1.6 Let $G$ be a connected simply connected nilpotent group with corresponding lie algebra $\mathcal{L}$. Then a basis $X_{1}, \ldots, x_{n}$ of $\mathcal{L}$ can be chosen which has the property $\left\{X_{i+1}, \ldots, X_{n}\right\}$ is an ideal of $\left\{x_{i}, \ldots, X_{n}\right\}$ for $i=1, \ldots, n-1$, and hence every element in $G$ may be written uniquely in the form $\exp t_{1} X_{1} \ldots \exp _{n} X_{n}$.

The Hamiltonian vector fields on $M$, such that each Hamiltonian function $H=p^{\prime} f(q)$ can be decomposed as follows, let

$$
\begin{aligned}
& x_{H_{1}}=\bar{X}_{1}+x_{1} \\
& x_{H_{2}}=\bar{X}_{2}+x_{2} \\
& x_{i}=f_{1}^{i}\left(q_{1}, \ldots, q_{n}\right)^{\partial} / \partial q_{1}+\ldots+f_{n}^{i}\left(q_{1}, \ldots, q_{n}\right)^{\partial / \partial q_{n}}, i=1,2 \\
& \bar{x}_{i}=\sum_{k=1}^{n} \sum_{j=1}^{n} p_{k} \frac{\partial f_{k}^{i}}{\partial q_{j}}{ }^{\partial} / \partial p_{j} \quad, i=1,2
\end{aligned}
$$

$\Delta$ is the involutive distribution spanned by ${ }^{\partial} / \partial p_{1}, \ldots,{ }^{\partial} / /_{\partial \bar{p}_{n}}$ and
let $\Delta_{n}$ be the involutive distribution spanned by $\frac{\partial}{\partial q_{1}}, \ldots,{ }^{\partial} / \partial q_{n}$, then notice that

$$
\left[\bar{x}_{1}, \bar{x}_{2}\right] \in \Delta_{n},\left[\bar{x}_{1}, x_{2}\right] \in \Delta,\left[x_{1}, \bar{x}_{2}\right] \subset \Delta
$$

and $\left[x_{1}, x_{2}\right] \in \Delta$.

Let $\overline{\mathcal{L}}=\left\{\overline{\mathrm{X}}_{0}, \overrightarrow{\mathrm{X}}_{1}, \ldots, \overrightarrow{\mathrm{X}}_{\mathrm{m}}\right\}_{\text {L.A. }}$. with corresponding Lie group $\overline{\mathrm{G}}$ which is assumed to be nilpotent and acts transitively on N. Thus by Matsushima's theorem 5.1.4 there exists a basis $\bar{X}_{1}, \ldots, \bar{X}_{n}$ such that each element in $\vec{C}$ can be expressed by

$$
\operatorname{expt}_{1} \bar{x}_{1} \ldots \operatorname{expt}_{n} \bar{x}_{n} .
$$

 Also since $\bar{X}_{1}(x), \ldots, \bar{X}_{n}(x)$ spans $T_{x} N, \forall x \in N$ by transitivity, a dual basis $\bar{\omega}_{i}$ to $\bar{X}_{i}$ can be chosen, i.e.

$$
\bar{w}_{2}\left(X_{j}\right)=\delta_{i j} \text { on } N, \quad 1 \leqslant i, j \leqslant n \text { such that } \bar{w}_{i} \text { spans } T_{x}^{*}{ }_{N}
$$ $\forall x \in N, 1 \leq i \leq n$.

Lemma 5.1.7 Let $\theta$ be the natural 1 -form on $T^{\star} \mathrm{N}$, given locally by $\sum_{i=1}^{n} p_{i}{ }^{d q}{ }_{i}$, and let $\omega_{i}$ be the extension of $\bar{\omega}_{i}$ defined by

$$
\omega_{i}\left(X_{H_{j}}\right)=\bar{\omega}_{i}\left(\bar{X}_{j}\right) \text { where each } X_{\mathrm{H}} \text { satisfies the above conditions. }
$$

$$
\operatorname{Then} \sum_{i=1}^{n} p_{i} d q_{i}=\sum_{i=1}^{n} \omega_{i} H_{i}
$$

Proof. It is only necessary to consider $\Delta_{n}$, since both forms vanish on the complement i.e. on $\Delta$. Also, since $\bar{X}_{1}(x), \ldots, \bar{X}_{n}(x)$ spans $T_{x}^{*} N$, $\forall x \in N$, it is sufficient to consider the Hamiltonian vector fields $X_{H_{1}} \ldots . . X_{H_{n}}$. Thus

$$
\begin{equation*}
\sum_{i=1}^{n} p_{i} d q_{i}\left(X_{H_{j}}\right)=\sum_{i=1}^{n} p_{i} \frac{\partial H_{j}}{\partial p_{i}} \tag{5.1,3}
\end{equation*}
$$

and


But since $H_{j}=\sum_{i=1}^{n} p_{i} f_{i}^{j}(q), H_{j}=\sum_{i=1}^{n} p_{i} \frac{\partial H_{i}}{\partial p_{i}}$.

Notice that under poisson bracket considerations that $X_{H_{j}}\left(H_{k}\right)=\left\{H_{k}, H_{j}\right\}$ and that this vanishes on $N$ i.e.
$\left.X_{H_{j}}\left(H_{k}\right)\right|_{N}=0$
thus the flow of each $\bar{X} \in \overline{\mathcal{L}}$ does indeed stay on $N$.
Since $\bar{G}$ is a nilpotent transitive lie group action on $N$ the coordinate system of Croucli [5]. Theorem 2.1.10 can be used to reparameterize the $q-$ coordinates on $N$.

That is
$\left(q_{1}, \ldots, q_{n}\right) \longmapsto\left(I_{1}, \ldots, I_{n}\right)$
where

$$
\begin{aligned}
& \dot{I}_{1}=f_{1}=\text { constant } \\
& \dot{I}_{2}=f_{2}\left(I_{1}\right) \\
& \vdots \\
& \dot{I}_{n}=E_{n}\left(I_{1}, \ldots, I_{n-1}\right)
\end{aligned}
$$

where $f_{1}, \ldots, f_{n}$ are polynomials in the $I$ - coordinates.
In $I_{1}, \ldots, I_{n}$ coordinates $\bar{X}$ is given by

$$
\begin{equation*}
f_{1} \quad \partial / \partial I_{1}+f_{2}\left(I_{1}\right)^{\partial} / \partial I_{2}+\ldots+f_{n}\left(I_{1}, \ldots I_{n-1}\right)^{\partial} / \partial I_{n} \tag{5.1.5}
\end{equation*}
$$

Now introduce coordinates $\phi_{1}, \ldots, \phi_{n}$ such that $I_{1}, \ldots, I_{n}, \phi_{1}, \ldots, \phi_{n}$ form canonical coordinates on a neighbourhood of $\mathrm{T}^{\star} \mathrm{N}$. The aim is to show that there exists a local symplectomorphism from ( $q_{1}, \ldots, q_{n}, p_{1}, \ldots, p_{n}$ ) to ( $I_{1}, \ldots, I_{n}, \phi_{1}, \ldots, \phi_{n}$ )..

If $H$ is a Hamiltonian function on $\mathrm{T}^{\star} \mathrm{N}$ with $\mathrm{X}_{\mathrm{H}}=\overline{\mathrm{X}}+\mathrm{X}$ as before and $\bar{X}$ as in equation (5.1.5), then $H$ must be of the form

$$
H=\phi_{1} f_{1}+\phi_{2} f_{2}\left(I_{1}\right)+\ldots+\phi_{n} f_{n}\left(I_{1}, \ldots, I_{n-1}\right)+f_{0}\left(I_{1}, \ldots, I_{n}\right)
$$

Locally $N$ is defined by $\phi_{1} \ldots \ldots \phi_{n}=0$ thus $f_{0}\left(I_{1}, \ldots, I_{n}\right)=0$ on $T^{*} N$ since it only depends on $I_{i}, 1 \leqslant i \leqslant n$.

Theorem 5.1.8 Let $\left(U, \psi_{U}\right)$ be a chart on $T^{*} N$ with canonical coordinates $\left(q_{1}, \ldots, q_{n}, p_{1}, \ldots, p_{n}\right)$ and $x_{0} \in U$. Let $\left(v, \psi_{v}\right), x_{o} \in V$ be a chart on $T^{*} N$ with the above canonical coordinates ( $I_{1}, \ldots, I_{n}, \phi_{1}, \ldots, \phi_{n}$ ).

Then, after suitable restrictions of $U$ and $V$ there exists a local symplectomorphism $f: U^{\prime} \longmapsto V^{\prime}, U^{\prime} \subset U, V^{\prime} \subset V$.

Proof. It is sufficient to prove that

$$
\sum_{i=1}^{n} \omega_{i} H_{i}=\sum_{i=1}^{n} \phi_{i} d I_{i}
$$

and then applying lemma 5.7.1.
As

$$
\begin{align*}
& H_{i}=\sum_{j=1}^{n} \phi_{j} f_{j}^{i}\left(I_{1}, \ldots, I_{j-1}\right) \\
& \sum_{i=1}^{n} \omega_{i} H_{i}=\sum_{i=1}^{n} \sum_{j=1}^{n} \omega_{i} \phi_{j} f_{j}^{i}\left(I_{1}, \ldots, I_{j-1}\right) \tag{5.1.6}
\end{align*}
$$

Also dI ${ }_{j}\left(X_{H_{k}}\right)=\frac{\partial H_{k}}{\partial \phi_{j}}=f_{j}^{k}\left(I_{1}, \ldots, I_{j-1}\right)$
Thus, as $\omega_{i}\left(X_{H_{k}}\right)=\delta_{i k}$,

$$
\begin{equation*}
d I_{j}=\sum_{i=1}^{n} \omega_{i} f_{j}^{i}\left(I_{1}, \ldots, I_{j-1}\right) \tag{5.1.7}
\end{equation*}
$$

Combining (5.1.6) and (5.1.7) gives result.
Q.E.D.

The above structure can be understood more fully with the aid of an example.

Example 5.1.9 Take as the state space ( $\left.T^{*}{\underset{\sim}{R}}^{4}, \omega\right)$ where $\omega=\sum_{i=1} d I{ }_{i} \wedge \phi_{i}$. $\left(I_{1}, \ldots, I_{4}, \phi_{1}, \ldots, \phi_{4}\right)$ being canonical coordinates on $T^{*} \mathbb{R}^{4}$.

$$
\text { Suppose } \dot{x}=X_{H_{0}}+\sum_{i=1}^{3} u_{i} X_{H_{i}} \quad x(0)=X_{0}, x \in T^{*}{\underset{\sim}{R}}^{4} \text { where }
$$

$$
H_{0}=\phi_{1}
$$

$$
H_{1}=\phi_{2}+I_{1} \phi_{3}
$$

$$
\mathrm{H}_{2}=\phi_{3}+\mathrm{I}_{2} \phi_{4}
$$

$$
H_{3}=\phi_{4}
$$

so $\quad X_{H_{0}}={ }^{\partial} / \partial I_{1}$
$X_{H_{1}}=\partial / \partial I_{2}+I_{1}{ }^{\partial} / \partial I_{3}-\phi_{3}{ }^{\partial} / \partial \rho_{1}$
$X_{H_{2}}=\partial / \partial I_{3}+I_{2} \partial / \partial I_{4}-\phi_{4} \partial / \partial \phi_{2}$

$$
X_{H_{3}}=\partial /_{\partial I_{4}}
$$

$\operatorname{Let} \mathcal{L}=\left\{X_{H_{0}}, X_{H_{1}}, X_{H_{2}}, X_{H_{3}}\right\}$ L.A.

It is clear that $\mathcal{L}$ is a nilpotent Lie algebra since

$$
\left[\mathrm{X}_{\mathrm{H}_{0}}, \mathrm{X}_{\mathrm{H}_{1}}\right]=\frac{\partial}{\partial \mathrm{I}_{3}},\left[\mathrm{X}_{\mathrm{H}_{1}}, \mathrm{X}_{\mathrm{H}_{2}}\right]=\frac{\partial}{\partial \mathrm{I}_{4}} \text { and all other brackets }
$$

vanish.

Let $N=\left\{x \in T{ }^{*}{\underset{\sim}{R}}^{4}: H(x)=0, \forall X_{H} \in \mathcal{L}\right\}$
$\mathbb{N}$ is obviously Lagrangian defined by $\phi_{1}=\phi_{2}=\phi_{3}=\phi_{4}=0$ and

$$
\begin{aligned}
& \overline{\mathrm{X}}_{1}=\partial / \partial \mathrm{I}_{1} \\
& \overline{\mathrm{x}}_{2}=\partial / \partial \mathrm{I}_{2}+\mathrm{I}_{1} \partial / \partial \mathrm{I}_{3} \\
& \overline{\mathrm{x}}_{3}=\partial / \partial \mathrm{I}_{3}+\mathrm{I}_{2} \partial / \partial \mathrm{I}_{4} \\
& \overline{\mathrm{x}}_{4}=\partial / \partial \mathrm{I}_{4}
\end{aligned}
$$

These vector fields span $T_{x} N, \forall x<N$, so this system is in accordance with the previous theory. The dual forms $\omega_{i}, 1 \leqslant i \leqslant 4$, of lamina 5.1 .7

$$
\text { i.e. } \omega_{i}\left(X_{H_{j}}\right)=\delta_{i j} \text { are }
$$

$$
\mathrm{w}_{1}=\mathrm{dI} 1
$$

$$
w_{2}=d I_{2}
$$

$$
\omega_{3}=\mathrm{dI}_{3}-\mathrm{I}_{1} \mathrm{dI}_{2}
$$

$$
\omega_{4}=d I_{4}-I_{2} d I_{3}+I_{1} I_{2} d I_{2}
$$

$$
\text { Thus } \begin{aligned}
\sum_{i=1}^{4} \omega_{i} H_{i} & =\phi_{1} \mathrm{dI}_{1}+\left(\phi_{2}+\mathrm{I}_{1} \phi_{3}\right) \mathrm{dI} 2 \\
& +\left(\phi_{3}+\mathrm{I}_{2} \mathrm{~d} \phi_{4}\right)\left(\mathrm{dI}_{3}-\mathrm{I}_{1} \mathrm{dI}_{2}\right) \\
& +\phi_{4}\left(\mathrm{dI}_{4}-\mathrm{I}_{2} \mathrm{dI}_{3}+\mathrm{I}_{1} \mathrm{I}_{2} \mathrm{dI}_{2}\right) \\
& =\sum_{i=1}^{4} \phi_{i} \mathrm{dI}_{\mathrm{i}}
\end{aligned}
$$

Hence,

$$
\begin{aligned}
\phi_{1} d I_{1} & =\omega_{1} H_{1} \\
\phi_{2} d I_{2} & =\omega_{2} H_{2}-\phi_{3} I_{1} d I_{2} \\
\phi_{3} d I_{3} & =\omega_{3} H_{3}+\phi_{3} I_{1} d I_{2}-\phi_{4} I_{2} d I_{3}+\phi_{4} I_{1} I_{2} d I_{2} \\
\phi_{4} d I_{4} & =\omega_{4} H_{4}+\phi_{4} I_{2} d I_{3}-\phi_{4} I_{1} I_{2} d I_{2} \\
\text { Let } \mu_{1} & =\phi_{3} I_{1}{d I_{2}}, \mu_{2}=\phi_{4} I_{2}\left(I_{1} d I_{2}-d I_{3}\right) \\
\text { then } \phi_{1} d I_{1} & =\omega_{1} H_{1} \\
\phi_{2} d I_{2} & =\omega_{2} H_{2}-\mu_{1} \\
\phi_{3} d I_{3} & =\omega_{3} H_{3}+\mu_{1}+\mu_{2} \\
\phi_{4} d I_{4} & =\omega_{4} H_{4}-\mu_{2}
\end{aligned}
$$

So if $\zeta_{1}=0, \zeta_{2}=-\mu_{1}, \zeta_{2}=\mu_{1}+\mu_{2}, \zeta_{4}=-\mu_{2}$

$$
\sum_{i=1}^{4} \zeta_{i}=0 \text { and so } \sum_{i=1}^{4} \phi_{i} d I_{i}
$$

But each $\phi_{i} \mathrm{dI}_{\mathrm{i}}=\omega_{i} \mathrm{H}_{\mathrm{i}}+\zeta_{i}(\mathrm{I}, \phi)$

This example illustrates the difficulty of calculating the coordinates $\phi_{1}, \ldots, \phi_{n}$ in general, since the terms $\zeta_{1}, \cdots, \zeta_{n}$ are not necessarily closed, this makes integration normally impossible.

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