# Implicit Hamiltonian systems and singular curves of distributions 

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## §1. Introduction

The subject of the sub-Riemannian geometry is to study a triple $(M, \mathcal{D}, g)$ of a manifold $M$, a distribution $\mathcal{D}$ on $M$ and a bi-linear positive definite form $g$ on $\mathcal{D}$, which is called a sub-Riemannian manifold. Here, a distribution is a sub-bundle of the tangent bundle $T M$ of $M$. The object appears naturally as a collapsing Riemannian manifold and is a generalization of a Riemannian manifold. However properties of sub-Riemannian manifolds are much different from those of Riemannian manifolds. In fact, not much is known about the properties of exponential maps and even about smoothness of minimizers. For example there is a problem which is open for decades: "Are all locally minimizers smooth on any sub-Riemannian manifold?"

In this article we give a new clue to study of minimizers on subRiemannian manifolds.

For a distribution $\mathcal{D}$, there is an important class of curves called horizontal curves. A horizontal curve is an absolutely continuous curve $\gamma: I \rightarrow M$ such that $\dot{\gamma}(t)$ is a measurable and bounded map which satisfies $\dot{\gamma}(t) \in \mathcal{D}_{\gamma(t)}$ for almost every $t \in I$.

According to Chow-Rashevsky's theorem, if a distribution $\mathcal{D}$ on a connected manifold $M$ satisfies Hörmander's condition, every two points are connected by a horizontal curve. For a sub-Riemannian manifold

[^0]which satisfies Hörmander's condition, we may define a distance
\[

$$
\begin{aligned}
d_{C C}(p, q):=\inf _{\gamma}\{L(\gamma) & \left.:=\int_{[a, b]} \sqrt{g(\dot{\gamma}(t), \dot{\gamma}(t)}\right) d t \mid \\
& \gamma:[a, b] \rightarrow M: \text { horizontal }, \gamma(a)=p, \gamma(b)=q\}
\end{aligned}
$$
\]

which is called a Carnot-Carathéodory (or sub-Riemannian) distance. It is known that the topology from Carnot-Carathéodory distance agrees with the original one (ball-box theorem[10]).

A horizontal curve $\gamma$ connecting $p$ and $q$ is called a minimizer if $d_{C C}(p, q)=L(\gamma)$. A horizontal curve $\gamma: I \rightarrow M$ is a local minimizer if for any $t_{0} \in I$, there exists $\varepsilon>0$ such that for all closed sub-interval $J$ of $\left[t_{0}-\varepsilon, t_{0}+\varepsilon\right],\left.\gamma\right|_{J \cap I}$ is a minimizer between the end points. Note that any minimizer is necessarily a local minimizer.

To consider local minimizers, we classify horizontal curves on a subRiemannian manifold ( $M, \mathcal{D}, g$ ) by using the end-point mapping.

For a bounded measurable curve $c:[0, T] \rightarrow \mathcal{D}$, if a curve $\gamma:=$ $\pi_{\mathcal{D}} \circ c:[0, T] \rightarrow M$ satisfies $\dot{\gamma}(t)=c(t)$ for almost everywhere on $[0, T]$, then $\gamma$ is a horizontal curve and $c$ is called an admissible velocity. Here $\pi_{\mathcal{D}}: \mathcal{D} \rightarrow M$ is the canonical projection.

Expressing $c(t)$ as $c(t)=u_{1}(t) X_{1}(\gamma(t))+\cdots+u_{k}(t) X_{k}(\gamma(t))$ with respect to a local framing $\left\{X_{1}, \ldots, X_{k}\right\}$ of the distribution $\mathcal{D}$ on an open subset $O \subset M$, the condition is written as

$$
\dot{\gamma}(t)=u_{1}(t) X_{1}(\gamma(t))+\cdots+u_{k}(t) X_{k}(\gamma(t))
$$

with the fibre coordinates $\left(u_{1}, \ldots, u_{k}\right)$.
For any point $q_{0}$ on $M$, the set of admissible velocities

$$
\mathcal{V}_{q_{0}}:=\left\{c \mid c:[0, T] \rightarrow \mathcal{D}: \text { admissible velocity, } \gamma(0)=q_{0}\right\}
$$

will be a Banach manifold. The map

$$
\operatorname{End}\left(q_{0}\right): \mathcal{V}_{q_{0}} \rightarrow M, c \mapsto \gamma(T)
$$

is called an end-point mapping and is differentiable by means of Fréchet derivative. A singular (resp. regular) point of the end-point mapping is called a singular (resp. regular) velocity. The trajectory corresponds to singular velocity is called a singular (resp. regular) curve, i.e. the differential map

$$
d \operatorname{End}\left(q_{0}\right)_{c}: T_{c} \mathcal{V}_{q_{0}} \rightarrow T_{\gamma(T)} M
$$

is not surjective (resp. surjective) at singular (resp. regular) velocity $c$. Since every curve is either regular or singular, every minimizer is also either a regular curve or a singular curve.

For a sub-Riemannian manifold, there is a geodesic equation given as Hamiltonian formulation, not as Lagrangian. We give a function

$$
H_{E}(x, p)=-\frac{1}{2} \sum_{i, j} g^{i j}(x)\left\langle p, X_{i}(x)\right\rangle\left\langle p, X_{j}(x)\right\rangle
$$

on $T^{*} M$, where $g_{i j}=g\left(X_{i}, X_{j}\right)$ and $\left(g^{i j}\right)_{i, j}$ is the inverse matrix of $\left(g_{i j}\right)_{i, j}$. We may consider a Hamiltonian vector field associated to $H_{E}$ with canonical symplectic form on the cotangent bundle $T^{*} M$. An ordinary differential equation related to the Hamiltonian vector field is given by

$$
\dot{x}(t)=\frac{\partial H_{E}}{\partial p}(x(t), p(t)), \quad \dot{p}(t)=-\frac{\partial H_{E}}{\partial x}(x(t), p(t))
$$

with Darboux coordinates $(x, p)$ of $T^{*} M$. A solution of this equation (which is called the geodesic equation) is called a normal bi-extremal and its projection to $M$ is called a normal extremal (or a normal geodesic).

It is known that regular local minimizers are normal geodesics, and so they are smooth since they are solutions of the geodesic equation. A singular minimizer is sometimes also a normal geodesic depending on a metric $g$. A singular (local) minimizer which is not a normal geodesic is called a strictly singular (local) minimizer. Examples of singular minimizers which are not normal on Martinet distribution are given by R. Montgomery in 1994 [8].

Although the examples given by R. Montgomery guarantee the existence of strictly singular minimizers, we know few other examples. So, we are interested in finding more examples of strictly singular minimizers. As a strategy, we divide the problem into two problems; to detect singular horizontal curves which are not normal geodesics and to find local minimizers among them. In this paper we concentrate on the former one and give new examples of singular horizontal curves which are not normal geodesics.

Now, we introduce a result of foothold in the study of singular curves. Take a function $H: T^{*} M \times_{M} \mathcal{D} \rightarrow \mathbb{R}, H(x, p, u):=\langle p, u\rangle$ for $x \in M, p \in T_{x}^{*} M$ and $u \in \mathcal{D}_{x}$, here $\times_{M}$ is a fibre product. Let $\left\{X_{1}, \ldots, X_{k}\right\}$ be a local framing of $\mathcal{D}$ on an open neighborhood $U_{x_{0}}$ of $x_{0}$ in $M$ and $\left(u_{1}, \ldots, u_{k}\right)$ the fibre coordinates related to the local framing. Then we have locally

$$
H(x, p, u)=\sum_{i=1}^{k} u_{i}\left\langle p, X_{i}(x)\right\rangle
$$

Then singular curves are characterized by a constrained Hamiltonian system;

Proposition 1.1 ([6], p.567). A horizontal curve $x(t)$ on $M$ with rank $k$ distribution $\mathcal{D}$ is a singular curve if and only if there exist a positive number $\varepsilon>0$, a curve $p(t)$ on $T_{x(t)}^{*} M \backslash\{0\}$ and $u(t) \in \mathcal{D}_{x(t)}$ such that the curve $(x(t), p(t), u(t))$ satisfies the following equation (which is called the constrained Hamiltonian system) for almost all $t \in[0, \varepsilon)$;

$$
\left\{\begin{array}{l}
\dot{x}(t)=\frac{\partial H}{\partial p}(x(t), p(t), u(t)) \\
\dot{p}(t)=-\frac{\partial H}{\partial x}(x(t), p(t), u(t)) \\
\frac{\partial H}{\partial u_{i}}(x(t), p(t), u(t))=0(1 \leq i \leq k)
\end{array}\right.
$$

A solution $(x(t), p(t))$ on $T^{*} M$ of the constrained Hamiltonian system in Proposition 1.1 is called an abnormal bi-extremal and the projection of an abnormal bi-extremal to $M$ is called an abnormal extremal (or a singular curve). It is known that a local minimizer is either a normal extremal or an abnormal extremal; and the two possibilities are not mutually exclusive. Abnormal extremals are not local minimizers in general.

We may consider a constrained Hamiltonian system for any distribution. However it is not known whether there exists a solution or does not, in general. To study the existence of solutions, we intend to apply the theory of solvability of implicit differential systems. The theory, which is introduced in $\S 2$, is given by T. Fukuda and S. Janeczko. We give a refinement of Fukuda-Janeczko's theory in the case of two parameters in $\S 3$. In $\S 4$, an application of the results in $\S 3$ to rank two distribution is given. Then we have the following main result in this paper.

Theorem 1. Let $(M, \mathcal{D}, g)$ be a sub-Riemannian smooth manifold with a distribution $\mathcal{D}$ of rank two. Suppose that $\mathcal{D}_{1}:=\mathcal{D}+[\mathcal{D}, \mathcal{D}]$ is a sub-bundle of rank three and $\mathcal{D}_{2}:=\mathcal{D}_{1}+\left[\mathcal{D}, \mathcal{D}_{1}\right]$ is a sub-bundle of rank four. Then for any point $q$ in $M$, there exist an open neighborhood $U_{q}$ of $q$ in $M$ and a $C^{\infty}$ immersive singular curve $x(t)$ which is not a normal geodesic in $U_{q}$ defined on a small interval.

Theorem 1 is proved in $\S 4$ as an application of the results in $\S 3$. It is not known whether the singular curve in Theorem 1 is a minimizer or not.

## §2. Preliminary

### 2.1. Implicit differential systems and their solvability

An implicit differential system is a generalization of a differential equation defined by a vector field. We define some basic notions. Let $\pi: T M \rightarrow M$ be a canonical projection and $N$ a submanifold of $M$.

Definition 2.1. An implicit differential system on $M$ is a subset $S$ of tangent bundle $T M$.

A $C^{1}$ curve $\gamma:(a, b) \rightarrow N$ is called a solution of $S$ over $N$ if $(\gamma(t), \dot{\gamma}(t)) \in S \cap \pi^{-1}(N)$ for all $t \in(a, b)$.

A point $\left(x_{0}, \dot{x}_{0}\right) \in S$ is a solvable point of $S$ over $N$ if there exist a positive number $\varepsilon>0$ and a solution $\gamma:(-\varepsilon, \varepsilon) \rightarrow N$ such that $(\gamma(0), \dot{\gamma}(0))=\left(x_{0}, \dot{x}_{0}\right)$.

From now on we consider the case where $S$ is a smooth submanifold of $T M$. A point $\left(x_{0}, \dot{x}_{0}\right) \in S$ is a smoothly solvable point of $S$ over $N$ if there exist an open neighborhood $W$ in $S \times \mathbb{R}$ of $\left(x_{0}, \dot{x}_{0}, 0\right)$ and $C^{\infty}$ map $\bar{\gamma}: W \rightarrow N$ such that $\gamma_{(x, \dot{x})}(t):=\bar{\gamma}(x, \dot{x}, t)$ is a solution of $S$ over $N$ with $(\gamma(0), \dot{\gamma}(0))=(x, \dot{x})$ for all $(x, \dot{x}) \in \pi_{1}(W)$, where $\pi_{1}: S \times \mathbb{R} \rightarrow S$ is a natural projection.

An implicit differential system $S$ over $N$ is called a smoothly solvable submanifold over $N$ if $S$ consists only of smoothly solvable points of $S$ over $N$.

When a submanifold $N$ is $M$ itself, definitions above are given in Fukuda and Janeczko's papers [4][5]. A solution of $S$ over $M$ is just called a solution of $S$. We say, simply, an implicit differential system $S$ is (smoothly) solvable if $S$ is (smoothly) solvable over $M$.

### 2.2. Implicit Hamiltonian systems

Let $(M, \omega)$ be a symplectic manifold. Then there is the induced symplectic structure $\dot{\omega}$ on the tangent bundle $T M$ : We have an bundle isomorphism induced from interior product $b: T M \rightarrow T^{*} M, b_{x}\left(v_{q}\right)=$ $\iota_{v_{q}} \omega_{q}$ for each point $q \in M$. The symplectic structure $\dot{\omega}$ is given by the pullback of the Liouville form $\theta$ on $T^{*} M$, i.e., $\dot{\omega}:=b^{*} d \theta$. The induced symplectic structure $\dot{\omega}$ is locally written by

$$
\dot{\omega}=\sum_{i=1}^{n} d \dot{p}_{i} \wedge d x_{i}-d \dot{x}_{i} \wedge d p_{i}
$$

with the canonical coordinates $(x, p, \dot{x}, \dot{p})$ of tangent bundle $T M$ related to Darboux coordinates $(x, p)=\left(x_{1}, \ldots, x_{n}, p_{1}, \ldots, p_{n}\right)$ for the standard symplectic form $\omega=\sum_{i=1}^{n} d p_{i} \wedge d x_{i}$ of $M$.

We will define the notion of implicit Hamiltonian systems as Lagrangian submanifolds of $(T M, \dot{\omega})$. Then they are regarded as a generalization of Hamiltonian vector fields, i.e., of Hamiltonian dynamical systems. In what follows we set $M=\mathbb{R}^{2 n}$ with the standard symplectic form $\omega$ as above.

Definition $2.2([4,5])$. A Lagrangian submanifold $L$ of $\left(T \mathbb{R}^{2 n}, \dot{\omega}\right)$ (i.e., $\operatorname{dim} L=2 n$ and $\left.\omega\right|_{L}=0$ ) is called an implicit Hamiltonian system.

There is a well-known result that a Lagrangian submanifold is locally generated by a Morse family;

Theorem 2.3 ([1]). Let $L$ be a Lagrangian submanifold of $T \mathbb{R}^{2 n}$ and $\left(q_{0}, \dot{q}_{0}\right)=\left(x_{0}, p_{0}, \dot{x}_{0}, \dot{p}_{0}\right) \in L$. Suppose

$$
\operatorname{corank} d\left(\left.\pi\right|_{L}\right)\left(q_{0}, \dot{q}_{0}\right)=k>0
$$

Then there exist an open neighborhood $O$ of $\left(q_{0}, \dot{q}_{0}\right)$ in $T \mathbb{R}^{2 n}$, an open neighborhood $W$ of $\left(q_{0}, 0\right) \in \mathbb{R}^{2 n} \times \mathbb{R}^{k}$ and a smooth function $F: W \rightarrow \mathbb{R}$ such that

$$
\begin{array}{r}
L \cap O=\left\{\left(x_{0}, p_{0}, \dot{x}_{0}, \dot{p}_{0}\right) \in O \mid \exists u \in \mathbb{R}^{k} \text { s.t. }(x, p, u) \in W, \frac{\partial F}{\partial u_{l}}(x, p, u)=0,\right. \\
\left.\dot{x}_{i}=\frac{\partial F}{\partial p_{i}}(x, p, u), \dot{p}_{i}=-\frac{\partial F}{\partial x_{i}}(x, p, u), 1 \leq i \leq n, 1 \leq l \leq k\right\}
\end{array}
$$

and that
$\operatorname{rank}\left(\frac{\partial^{2} F}{\partial x_{i} \partial u_{l}}\left(q_{0}, 0\right), \frac{\partial^{2} F}{\partial p_{i} \partial u_{l}}\left(q_{0}, 0\right)\right)_{1 \leq i \leq n, 1 \leq l \leq k}=k, \frac{\partial^{2} F}{\partial u_{r} \partial u_{s}}\left(q_{0}, 0\right)=0$
for $1 \leq r, s \leq k$.
Recall that the family of functions $F: \mathbb{R}^{2 n} \times \mathbb{R}^{k} \rightarrow \mathbb{R}$ with $2 n$ parameters $\left(x_{1}, \ldots, x_{n}, p_{1}, \ldots, p_{n}\right)$ on $\left(\mathbb{R}^{k} ; u_{1}, \ldots, u_{k}\right)$ is called a Morse family if $0 \in \mathbb{R}^{k}$ is a critical point of the map $F\left(q_{0}, u\right)$ and the map

$$
\left(\frac{\partial F}{\partial u_{1}}, \ldots, \frac{\partial F}{\partial u_{k}}\right): \mathbb{R}^{2 n} \times \mathbb{R}^{k} \rightarrow \mathbb{R}^{k}
$$

is submersive at $\left(q_{0}, 0\right)$. We denote $L_{F}$ a Lagrangian submanifold generated by Morse family $F: \mathbb{R}^{2 n} \times \mathbb{R}^{k} \rightarrow \mathbb{R}$. That is, for the catastrophe set

$$
C(F)=\left\{(x, p, u) \in \mathbb{R}^{2 n} \times \mathbb{R}^{k} \left\lvert\, \frac{\partial F}{\partial u_{i}}(x, p, u)=0\right., i=1, \ldots, k\right\}
$$

of $F$ and $C^{\infty} \operatorname{map} \phi_{F}: \mathbb{R}^{2 n} \times \mathbb{R}^{k} \rightarrow T \mathbb{R}^{2 n}$ defined by

$$
\phi(x, p, u)=\left(x, p, \frac{\partial F}{\partial p_{i}}(x, p, u),-\frac{\partial F}{\partial x_{i}}(x, p, u)\right)
$$

we set $L_{F}=\phi_{F}(C(F))$.
The following propositions are given in [4] which are a necessary condition and a sufficient condition for $L_{F}$ to be (smoothly) solvable in the sense of Definition 2.1.

Proposition 2.4 ([4]). Let $(x, p, \dot{x}, \dot{p})$ be a solvable point of $L_{F}$. We set $\phi_{F}(x, p, u)=(x, p, \dot{x}, \dot{p})$. Then there exists a real vector $\mu=$ $\left(\mu_{1}, \ldots, \mu_{k}\right)$ in $\mathbb{R}^{k}$ such that
$\left(\begin{array}{ccc}\frac{\partial^{2} F}{\partial u_{1} \partial u_{1}}(x, p, u) & \cdots & \frac{\partial^{2} F}{\partial u_{1} \partial u_{k}}(x, p, u) \\ \vdots & \ddots & \vdots \\ \frac{\partial^{2} F}{\partial u_{k} \partial u_{1}}(x, p, u) & \cdots & \frac{\partial^{2} F}{\partial u_{k} \partial u_{k}}(x, p, u)\end{array}\right)\left(\begin{array}{c}\mu_{1} \\ \vdots \\ \mu_{k}\end{array}\right)=\left(\begin{array}{c}\left\{\frac{\partial F}{\partial u_{1}}, F\right\}(x, p, u) \\ \vdots \\ \left\{\frac{\partial F}{\partial u_{k}}, F\right\}(x, p, u)\end{array}\right)$.
Here the bracket $\{$,$\} is the Poisson bracket associated to the symplectic$ form $\omega$.

Proposition 2.5 ([4]). A point ( $x, p, \dot{x}, \dot{p}$ ) in $L_{F}$ is (smoothly) solvable if a linear equation

$$
\left(\begin{array}{ccc}
\frac{\partial^{2} F}{\partial u_{1} \partial u_{1}}(x, p, u) & \cdots & \frac{\partial^{2} F}{\partial u_{1} \partial u_{k}}(x, p, u) \\
\vdots & \ddots & \vdots \\
\frac{\partial^{2} F}{\partial u_{k} \partial u_{1}}(x, p, u) & \cdots & \frac{\partial^{2} F}{\partial u_{k} \partial u_{k}}(x, p, u)
\end{array}\right)\left(\begin{array}{c}
\mu_{1}(x, p, u) \\
\vdots \\
\mu_{k}(x, p, u)
\end{array}\right)=\left(\begin{array}{c}
\left\{\frac{\partial F}{\partial u_{1}}, F\right\}(x, p, u) \\
\vdots \\
\left\{\frac{\partial F}{\partial u_{k}}, F\right\}(x, p, u)
\end{array}\right)
$$

has a (smooth) solution on a neighborhood of $(x, p, u)=\phi_{F}(x, p, \dot{x}, \dot{p})$ in $C(F)$.

The differences between the necessary condition and the sufficient condition appear as those of the domain and smoothness of the solution $\mu$.

Now we consider a Morse family of particular type:

$$
F: \mathbb{R}^{2 n} \times \mathbb{R}^{k} \rightarrow \mathbb{R}, \quad F(x, p, u)=\sum_{j=1}^{k} a_{j}(x, p) u_{j}+b(x, p)
$$

Note that functions $a_{1}, \ldots, a_{k}$ are independent, i.e., differential one forms $d a_{1}(x, p), \ldots, d a_{k}(x, p)$ are linearly independent in $T_{(x, p)}^{*} \mathbb{R}^{2 n}$ at each point $(x, p) \in \mathbb{R}^{2 n}$ because $F$ is a Morse family. The catastrophe set of $F$ is given by $C(F)=K \times \mathbb{R}^{k}$ with

$$
K:=\left\{(x, p) \in \mathbb{R}^{2 n} \mid a_{i}(x, p)=0, i=1, \ldots, k\right\}
$$

Applying Proposition 2.4 and Proposition 2.5, we know that $L_{F}$ is smoothly solvable if and only if $\left\{a_{i}, a_{j}\right\}(x, p)=0,\left\{b, a_{i}\right\}(x, p)=0$ on $K,(1 \leq i, j \leq k)([5])$. Considering conditions in the propositions, we set
$\widetilde{S_{F}}:=\left\{(x, p, u) \in C(F) \mid \sum_{j=1}^{k}\left\{a_{i}, a_{j}\right\}(x, p) u_{j}=\left\{b, a_{i}\right\}(x, p), 1 \leq i \leq k\right\}$, $S_{F}:=\phi_{F}\left(\widetilde{S_{F}}\right)$.

Then we see that every smoothly solvable submanifold of $L_{F}$ is contained in $S_{F}([5])$. Moreover $S_{F}$ itself may be smoothly solvable:

Theorem 2.6 ([5]). $S_{F}$ is a smoothly solvable submanifold of $L_{F}$ if

$$
\begin{aligned}
& \operatorname{rank}\left(\left\{a_{i}, a_{j}\right\}(x, p)\right)_{1 \leq i, j \leq k}=r(\text { constant }) \quad \text { and } \\
& \left(\begin{array}{c}
\left\{b, a_{1}\right\}(x, p) \\
\vdots \\
\left\{b, a_{k}\right\}(x, p)
\end{array}\right) \in \operatorname{Im}\left(\left\{a_{i}, a_{j}\right\}(x, p)\right)_{1 \leq i, j \leq k}
\end{aligned}
$$

holds for every $(x, p) \in K=\left\{(x, p) \in \mathbb{R}^{2 n} \left\lvert\, \frac{\partial F}{\partial u_{i}}(x, p, u)=0\right.,1 \leq i \leq k\right\}$.

## §3. Main results

Now we pose a question; for which submanifold $S$ of $L_{F}$ does there exist a submanifold $A$ of $K$ such that $S$ is smoothly solvable over $A$ ? In this section, we will give an answer to this question in the case $k=2$.

Let $F: \mathbb{R}^{2 n} \times \mathbb{R}^{2} \rightarrow \mathbb{R}$ be a Morse family which is defined by

$$
F(x, p, u)=a_{1}(x, p) u_{1}+a_{2}(x, p) u_{2} .
$$

We consider solvability of the Lagrangian submanifold $L_{F}$ which is generated by $F$. Moreover we consider solvability of submanifolds of $L_{F}$. In detail, we consider a map $\phi_{F}: \mathbb{R}^{2 n} \times \mathbb{R}^{2} \rightarrow T\left(\mathbb{R}^{2 n}\right)$ which is defined by

$$
\phi_{F}(x, p, u)=\left(x, p, \frac{\partial F}{\partial p}(x, p, u),-\frac{\partial F}{\partial x}(x, p, u)\right)
$$

and the catastrophe set

$$
C(F)=\left\{(x, p, u) \mid a_{1}(x, p)=a_{2}(x, p)=0\right\}=K \times \mathbb{R}^{2}
$$

of $F$ for $K=\left\{(x, p) \mid a_{1}(x, p)=a_{2}(x, p)=0\right\} \subset \mathbb{R}^{2 n}$, then we define $L_{F}$ by the image $\phi_{F}(C(F))$. According to Fukuda-Janeczko's Theorem $2.6, L_{F}$ is smoothly solvable if and only if $\left\{a_{1}, a_{2}\right\}=0$ locally on $K$.

Now we consider the cases where the assumptions of Theorem 2.6 are not fulfilled. We consider the family of vector fields $X_{u}: K \rightarrow T\left(\mathbb{R}^{2 n}\right)$ along $K$ with parameter $u$

$$
X_{u}(x, p)=\left(x, p, \frac{\partial F}{\partial p}(x, p, u),-\frac{\partial F}{\partial x}(x, p, u)\right)
$$

and detect submanifolds of $K$ on which $X_{u}$ are tangent to $K$. We are going to give smoothly solvable submanifolds over submanifolds of $K$ in the following series of Propositions 3.1-3.6.

Let $A_{0}=K$. The vector field $X_{u}$ is tangent to $A_{0}$ if and only if

$$
X_{u}\left(a_{1}\right)=X_{u}\left(a_{2}\right)=0 \quad \text { i.e. } u_{1}\left\{a_{1}, a_{2}\right\}=u_{2}\left\{a_{1}, a_{2}\right\}=0 .
$$

Hence $\phi_{F}\left(A_{0} \times \mathbb{R}^{2}\right)$ is smoothly solvable if and only if $\left\{a_{1}, a_{2}\right\}=0$ on $A_{0}=K$. This fact is also obtained as a corollary of Theorem 2.6.

Then we consider a submanifold $A_{1}$ of $A_{0}$ consisting of points at which $X_{u}$ is tangent to $A_{0}$;

$$
A_{1}:=\left\{(x, p) \mid a_{1}(x, p)=a_{2}(x, p)=\left\{a_{1}, a_{2}\right\}(x, p)=0\right\} .
$$

We assume that the functions $a_{1}, a_{2},\left\{a_{1}, a_{2}\right\}$ are independent. The vector field $X_{u}$ is tangent to $A_{1}$ if and only if
$X_{u}\left(\left\{a_{1}, a_{2}\right\}\right)=0 \quad$ i.e. $u_{1}\left\{a_{1},\left\{a_{1}, a_{2}\right\}\right\}(x, p)+u_{2}\left\{a_{2},\left\{a_{1}, a_{2}\right\}\right\}(x, p)=0$
on $A_{1}$. Note that we have $X_{u}\left(a_{1}\right)=X_{u}\left(a_{2}\right)=0$ from the definition of $A_{1}$. Let $\mathcal{E}_{\mathbb{R}^{2 n}, q_{0}}$ be the $\mathbb{R}$-algebra of $C^{\infty}$ function germs at $q_{0}$ on $\mathbb{R}^{2 n}$. We denote by $\left\langle a_{1}, a_{2},\left\{a_{1}, a_{2}\right\}\right\rangle_{\mathcal{E}_{\mathbb{R}^{2 n}, q_{0}}}$ the $\mathcal{E}_{\mathbb{R}^{2 n}, q_{0}}$-module generated by $a_{1}, a_{2}$ and $\left\{a_{1}, a_{2}\right\}$. We set

$$
\xi_{1}:=\left\{a_{1},\left\{a_{1}, a_{2}\right\}\right\}, \xi_{2}:=\left\{a_{2},\left\{a_{1}, a_{2}\right\}\right\} .
$$

Then the vector field $X_{u}$ is tangent to $A_{1}$ if the functions $\xi_{1}$ and $\xi_{2}$ belong to the $\mathcal{E}_{\mathbb{R}^{2 n}, q_{0}}$-module $\left\langle a_{1}, a_{2},\left\{a_{1}, a_{2}\right\}\right\rangle_{\mathcal{E}^{2 n}, q_{0}}$ for any point $q_{0} \in A_{1}$.

Proposition 3.1. Assume that $a_{1}, a_{2}$, and $\left\{a_{1}, a_{2}\right\}$ are independent. Then $\phi_{F}\left(A_{1} \times \mathbb{R}^{2}\right)$ is a smoothly solvable submanifold of $L_{F}$ over $A_{1}$ if and only if $\xi_{1}, \xi_{2} \in\left\langle a_{1}, a_{2},\left\{a_{1}, a_{2}\right\}\right\rangle_{\mathbb{R}^{2 n}, q_{0}}$ for any point $q_{0}$ in $A_{1}$.

To find smoothly solvable submanifolds of $L_{F}$ over submanifolds of $A_{1}$, we construct fiber bundles as follows. Let

$$
C_{(x, p)}:=\left\{\left(u_{1}, u_{2}\right) \mid u_{1}\left\{a_{1},\left\{a_{1}, a_{2}\right\}\right\}(x, p)+u_{2}\left\{a_{2},\left\{a_{1}, a_{2}\right\}\right\}(x, p)=0\right\}
$$

for $(x, p) \in K$ and define line bundles

$$
\begin{aligned}
& {\overline{A_{2}^{1}}}^{1}:=\left\{(x, p, u) \mid u \in C_{(x, p)}^{1},(x, p) \in A_{2}^{1}\right\}, \\
& {\overline{A_{2}^{1}}}^{2}:=\left\{(x, p, u) \mid u \in C_{(x, p)}^{2},(x, p) \in A_{2}^{1}\right\}, \\
& {\overline{A_{2}^{2}}}^{1}:=\left\{(x, p, u) \mid u \in C_{(x, p)}^{1},(x, p) \in A_{2}^{2}\right\}, \\
& {\overline{A_{2}^{2}}}^{2}:=\left\{(x, p, u) \mid u \in C_{(x, p)}^{2},(x, p) \in A_{2}^{2}\right\}, \\
& {\overline{A_{1,1}}}^{2}:=\left\{(x, p, u) \mid u \in C_{(x, p)}^{2},(x, p) \in A_{1,1}\right\}, \\
& {\overline{A_{1,2}}}^{1}:=\left\{(x, p, u) \mid u \in C_{(x, p)}^{1},(x, p) \in A_{1,2}\right\}, \\
& {\overline{A_{1,(1,2)}}}^{1,2}:=\left\{(x, p, u) \mid u \in C_{(x, p)}^{1,2},(x, p) \in A_{1,(1,2)}\right\},
\end{aligned}
$$

with

$$
\begin{array}{lr}
A_{2}^{1}:=A_{1} \cap\left\{(x, p) \mid \xi_{1}=0\right\}, & C_{(x, p)}^{1}=\left\{\left(u_{1}, 0\right) \in C_{(x, p)}\right\}, \\
A_{2}^{2}:=A_{1} \cap\left\{(x, p) \mid \xi_{2}=0\right\}, & C_{(x, p)}^{2}=\left\{\left(0, u_{2}\right) \in C_{(x, p)}\right\}, \\
A_{1,1}:=A_{1} \cap\left\{(x, p) \mid \xi_{1} \neq 0\right\}, & C_{(x, p)}^{1,2}=C_{(x, p)} \backslash\{0\}, \\
A_{1,2}:=A_{1} \cap\left\{(x, p) \mid \xi_{2} \neq 0\right\}, & \\
A_{1,(1,2)}:=A_{1} \cap\left\{(x, p) \mid \xi_{1} \neq 0, \xi_{2} \neq 0\right\} . &
\end{array}
$$

Let us consider the case that one of $\xi_{1}$ and $\xi_{2}$ belongs to the $\mathcal{E}_{\mathbb{R}^{2 n}, q_{0}}-$ module $\left\langle a_{1}, a_{2},\left\{a_{1}, a_{2}\right\}\right\rangle_{\mathcal{E}_{\mathbb{R}^{2 n}, q_{0}}}$ and the other does not.

Proposition 3.2. Assume that $a_{1}, a_{2}$ and $\left\{a_{1}, a_{2}\right\}$ are independent. Assume also that

$$
\xi_{2} \in\left\langle a_{1}, a_{2},\left\{a_{1}, a_{2}\right\}\right\rangle_{\mathcal{E}_{\mathbb{R}^{2 n}, q_{0}}} \text { and } \xi_{1} \notin\left\langle a_{1}, a_{2},\left\{a_{1}, a_{2}\right\}\right\rangle_{\mathcal{E}_{\mathbb{R}^{2 n}, q_{0}}}
$$

at every point $q_{0}$ of $A_{1}$. Then the followings hold.
(1) $\phi_{F}\left({\overline{A_{1,1}}}^{2}\right)$ is a smoothly solvable submanifold of $L_{F}$ over $A_{1,1}$.
(2) Assume, furthermore, that $\xi_{1}, a_{1}, a_{2},\left\{a_{1}, a_{2}\right\}$ are independent.
(a) $\phi_{F}\left({\overline{A_{2}^{1}}}^{2}\right)$ is a smoothly solvable submanifold of $L_{F}$ over $A_{2}^{1}$.
(b) $\quad \phi_{F}\left(A_{2}^{1} \times \mathbb{R}^{2}\right)$ is a smoothly solvable submanifold of $L_{F}$ over $A_{2}^{1}$ if $\left\{a_{1}, \xi_{1}\right\} \in\left\langle a_{1}, a_{2},\left\{a_{1}, a_{2}\right\}, \xi_{1}\right\rangle_{\mathcal{E}_{\mathbb{R}^{2 n}, q_{0}}}$ for any point $q_{0}$ in $A_{2}^{1}$.

Proof. (1): $A_{1,1}$ is an open submanifold of $A_{1}$ from the definition. Since there exist $\beta_{1}, \beta_{2}$ and $\beta_{3} \in \mathcal{E}_{\mathbb{R}^{2 n}, q_{0}}$ such that $\xi_{2}=\beta_{1} a_{1}+\beta_{2} a_{2}+$ $\beta_{3}\left\{a_{1}, a_{2}\right\}$, the vector $X_{u}(x, p)$ with $u \in C_{(x, p)}^{2}$ is tangent to $A_{1,1}$ at each
point $(x, p) \in A_{1,1}$ because

$$
\begin{aligned}
& X_{u}\left(a_{1}\right)=u_{2}\left\{a_{2}, a_{1}\right\}=0 \\
& X_{u}\left(a_{2}\right)=0 \cdot\left\{a_{1}, a_{2}\right\}=0 \\
& X_{u}\left(\left\{a_{1}, a_{2}\right\}\right)=0 \cdot \xi_{1}+u_{2} \xi_{2}=u_{2} \xi_{2}=\beta_{1} a_{1}+\beta_{2} a_{2}+\beta_{3}\left\{a_{1}, a_{2}\right\}=0
\end{aligned}
$$

on $A_{1,1}$.
(2)-(a): We check the condition that $X_{u}(x, p)$ is tangent to $A_{2}^{1}$ with $u \in C_{(x, p)}^{2}$ at each point $(x, p) \in A_{2}^{1}$. Note that $\left\{a_{1}, \xi_{2}\right\}=\left\{a_{2}, \xi_{1}\right\}$ from Jacobian identity:

$$
\begin{aligned}
\left\{a_{1},\left\{a_{2},\left\{a_{1}, a_{2}\right\}\right\}\right\} & =\left\{a_{2},\left\{a_{1},\left\{a_{1}, a_{2}\right\}\right\}\right\}+\left\{\left\{a_{1}, a_{2}\right\},\left\{a_{1}, a_{2}\right\}\right\} \\
& =\left\{a_{2},\left\{a_{1},\left\{a_{1}, a_{2}\right\}\right\}\right\} .
\end{aligned}
$$

Then

$$
\left.\left.\begin{array}{l}
\quad X_{u}\left(\left\{a_{1}, a_{2}\right\}\right)=0 \cdot \xi_{1}+u_{2} \xi_{2}=0, \\
X_{u}\left(\xi_{1}\right)=u_{1}\left\{a_{1}, \xi_{1}\right\}+u_{2}\left\{a_{2}, \xi_{1}\right\}=u_{1}\left\{a_{1}, \xi_{1}\right\}+u_{2}\left\{a_{1}, \xi_{2}\right\} \quad \ldots \\
= \\
=u_{2}\left(\left\{a_{1}, \beta_{1} a_{1}\right\}+\left\{a_{1}, \beta_{2} a_{2}\right\}+\left\{a_{1}, \beta_{3}\left\{a_{1}, a_{2}\right\}\right\}\right) \\
= \\
=u_{2}\left(a_{1}\left\{a_{1}, \beta_{1}\right\}+\beta_{2}\left\{a_{1}, a_{2}\right\}+a_{2}\left\{a_{1}, \beta_{2}\right\}\right. \\
=0
\end{array} \quad+\beta_{3} \xi_{1}+\left\{a_{1}, a_{2}\right\}\left\{a_{1}, \beta_{3}\right\}\right)\right\}
$$

on $A_{2}^{1}$.
(2)-(b): From an equality ( $\star$ ) we have

$$
X_{u}\left(\xi_{1}\right)=u_{1}\left\{a_{1}, \xi_{1}\right\}+u_{2}\left\{a_{2}, \xi_{1}\right\}=u_{1}\left\{a_{1}, \xi_{1}\right\}+u_{2}\left\{a_{1}, \xi_{2}\right\} .
$$

Since $\xi_{2} \in\left\langle a_{1}, a_{2},\left\{a_{1}, a_{2}\right\}\right\rangle_{\mathcal{E}^{2} n, q_{0}}$, it holds that

$$
\left\{a_{1}, \xi_{2}\right\} \in\left\langle a_{1}, a_{2},\left\{a_{1}, a_{2}\right\}, \xi_{2}\right\rangle_{\mathcal{E}_{\mathbb{R}^{2 n}, q_{0}}} .
$$

Consequently, by using $\left\{a_{1}, \xi_{1}\right\} \in\left\langle a_{1}, a_{2},\left\{a_{1}, a_{2}\right\}, \xi_{2}\right\rangle_{\mathcal{E}_{\mathbb{R}^{2 n}, q_{0}}}$, we obtain $X_{u}\left(\xi_{1}\right)=0$ on $A_{2}^{1}$.
Q.E.D.

In the same way we have the counterpart of Proposition 3.2.
Proposition 3.3. Assume that $a_{1}, a_{2}$ and $\left\{a_{1}, a_{2}\right\}$ are independent. Assume also that

$$
\xi_{1} \in\left\langle a_{1}, a_{2},\left\{a_{1}, a_{2}\right\}\right\rangle_{\mathcal{E}_{\mathbb{R}^{2 n}, q_{0}}} \text { and } \xi_{2} \notin\left\langle a_{1}, a_{2},\left\{a_{1}, a_{2}\right\}\right\rangle_{\mathcal{E}_{\mathbb{R}^{2 n}, q_{0}}}
$$

at every point $q_{0}$ of $A_{1}$. Then the followings hold.
(1) $\phi_{F}\left({\overline{A_{1,2}}}^{1}\right)$ is a smoothly solvable submanifold of $L_{F}$ over $A_{1,2}$.
(2) Assume, furthermore, that $\xi_{2}, a_{1}, a_{2},\left\{a_{1}, a_{2}\right\}$ are independent.
(a) $\phi_{F}\left({\overline{A_{2}^{2}}}^{1}\right)$ is a smoothly solvable submanifold of $L_{F}$ over $A_{2}^{2}$.
(b) $\quad \phi_{F}\left(A_{2}^{2} \times \mathbb{R}^{2}\right)$ is a smoothly solvable submanifold of $L_{F}$ over $A_{2}^{2}$ if $\left\{a_{2}, \xi_{2}\right\} \in\left\langle a_{1}, a_{2},\left\{a_{1}, a_{2}\right\}, \xi_{2}\right\rangle_{\mathcal{E}_{\mathbb{R}^{2 n}, q_{0}}}$ for any point $q_{0}$ in $A_{2}^{2}$.
In the case $\xi_{1}, \xi_{2} \notin\left\langle a_{1}, a_{2},\left\{a_{1}, a_{2}\right\}\right\rangle_{\mathcal{E}^{2 n}, q_{0}}$ we have
Proposition 3.4. Assume that $a_{1}, a_{2}$ and $\left\{a_{1}, a_{2}\right\}$ are independent. Then $\phi_{F}\left({\overline{A_{1,(1,2)}}}^{1,2}\right)$ is a smoothly solvable submanifold of $L_{F}$ over $A_{1,(1,2)}$ if $\xi_{1}, \xi_{2} \notin\left\langle a_{1}, a_{2},\left\{a_{1}, a_{2}\right\}\right\rangle_{\mathcal{E}^{2 n}, q_{0}}$ for every point $q_{0}$ in $A_{1,(1,2)}$.

Proof. $A_{1,(1,2)}$ is an open submanifold of $A_{1}$ from the definition. The vector $X_{u}(x, p)$ with $u \in C_{(x, p)}^{1,2}$ is tangent to $A_{1,(1,2)}$ at each point $(x, p) \in A_{1,(1,2)}$ since

$$
\begin{aligned}
& X_{u}\left(a_{1}\right)=u_{2}\left\{a_{2}, a_{1}\right\}=0 \\
& X_{u}\left(a_{2}\right)=u_{1}\left\{a_{1}, a_{2}\right\}=0 \\
& X_{u}\left(\left\{a_{1}, a_{2}\right\}\right)=u_{1} \xi_{1}+u_{2} \xi_{2}=0
\end{aligned}
$$

on $A_{1,(1,2)}$.
Q.E.D.

In Proposition 3.1-3.4, we gave sufficient conditions for existence of smoothly solvable submanifolds of $L_{F}$ over $A_{1,1}, A_{2}^{1}, A_{1,2}, A_{2}^{2}$ and $A_{1,(1,2)}$ and examples of smoothly solvable submanifolds of $L_{F}$ over them. The following two propositions give different sufficient conditions for existence of smoothly solvable submanifolds of $L_{F}$ over $A_{2}^{1}$ and $A_{2}^{2}$ and examples of smoothly solvable submanifolds over them respectively.

Proposition 3.5. Assume that $a_{1}, a_{2},\left\{a_{1}, a_{2}\right\}$ and $\xi_{1}$ are independent. Then $\phi_{F}\left({\overline{A_{2}^{1}}}^{1}\right)$ is a smoothly solvable submanifold of $L_{F}$ over $A_{2}^{1}$ if $\left\{a_{1}, \xi_{1}\right\} \in\left\langle a_{1}, a_{2},\left\{a_{1}, a_{2}\right\}, \xi_{1}\right\rangle_{\mathcal{R}^{2 n}, q_{0}}$ for any point $q_{0}$ in $A_{2}^{1}$.

Proof. Since $a_{1}, a_{2},\left\{a_{1}, a_{2}\right\}$ and $\xi_{1}$ are independent, $A_{2}^{1}$ is a submanifold of $K$. For the vector field $X_{u}$ along $A_{2}^{1}$ with $\left(u_{1}, 0\right)$,

$$
\begin{aligned}
& X_{u}\left(a_{1}\right)=0 \cdot\left\{a_{1}, a_{2}\right\}=0 \\
& X_{u}\left(a_{2}\right)=u_{1}\left\{a_{2}, a_{1}\right\}=0 \\
& X_{u}\left(\left\{a_{1}, a_{2}\right\}\right)=u_{1} \xi_{1}=0
\end{aligned}
$$

hold on $A_{2}^{1}$. Since $\left\{a_{1}, \xi_{1}\right\} \in\left\langle a_{1}, a_{2},\left\{a_{1}, a_{2}\right\}, \xi_{1}\right\rangle_{\mathcal{E}_{\mathbb{R}^{2 n}, q_{0}}}$, there exist $\beta_{1}, \beta_{2}, \beta_{3}$ and $\beta_{4} \in \mathcal{E}_{\mathbb{R}^{2 n}, q_{0}}$ such that $\left\{a_{1}, \xi_{1}\right\}=\beta_{1} a_{1}+\beta_{2} a_{2}+\beta_{3}\left\{a_{1}, a_{2}\right\}+$ $\beta_{4} \xi_{1}$. Hence we have

$$
X_{u}\left(\xi_{1}\right)=u_{1}\left\{a_{1}, \xi_{1}\right\}=u_{1}\left(\beta_{1} a_{1}+\beta_{2} a_{2}+\beta_{3}\left\{a_{1}, a_{2}\right\}+\beta_{4} \xi_{1}\right)=0
$$

on $A_{2}^{1}$. Thus the vector field $X_{u}$ with $\left(u_{1}, 0\right)$ is tangent to $A_{2}^{1}$. Q.E.D.
In the same way we have
Proposition 3.6. Assume that $a_{1}, a_{2},\left\{a_{1}, a_{2}\right\}$ and $\xi_{2}$ are independent. Then $\phi_{F}\left({\overline{A_{2}^{2}}}^{2}\right)$ is a smoothly solvable submanifold of $L_{F}$ over $A_{2}^{2}$ if $\left\{a_{2}, \xi_{2}\right\} \in\left\langle a_{1}, a_{2},\left\{a_{1}, a_{2}\right\}, \xi_{2}\right\rangle_{\mathcal{E}_{\mathbb{R}^{2 n}, q_{0}}}$ for any point $q_{0}$ in $A_{2}^{2}$.

## §4. Application

We apply the results obtained in $\S 3$, to study distributions and its singular curves and we prove Theorem 1. Now we recall some basic notations for the study of distributions. The Lie flag of a distribution $\mathcal{D}$ is the sequence $\mathcal{D}_{0} \subset \mathcal{D}_{1} \subset \cdots$ defined inductively by

$$
\mathcal{D}_{0}:=\mathcal{D}, \quad \mathcal{D}_{i+1}:=\mathcal{D}_{i}+\left[\mathcal{D}_{0}, \mathcal{D}_{i}\right], \quad i \geq 0
$$

The small growth vector of a distribution $\mathcal{D}$ at $q \in M$ is the sequence of the dimension of Lie flags;

$$
\left(\operatorname{dim} \mathcal{D}_{0}(q), \operatorname{dim} \mathcal{D}_{1}(q), \operatorname{dim} \mathcal{D}_{2}(q), \ldots\right)
$$

For example a contact distribution $\mathcal{D} \subset T M$ on a manifold $M$ of dimension $2 n+1$, has small growth vector $(2 n, 2 n+1)$. An Engel distribution $\mathcal{D}$ on 4 dimensional manifold $M$ has small growth vector $(2,3,4)$.

In this section we consider distributions with small growth vector $(2,3,4, \ldots)$. The following lemma plays an essential role throughout the section.

Lemma 4.1. Let $\mathcal{D}$ be a rank two distribution with small growth vector $(2,3,4, \ldots)$ at any point in an open neighborhood of $q \in M$ and $g$ a bi-linear positive definite form on $\mathcal{D}$. Then there exist an open neighborhood $U_{q}$ and local orthonormal frame $X_{1}, X_{2}$ of $\mathcal{D}$ on $U_{q}$ such that

$$
X_{1}, X_{2},\left[X_{1}, X_{2}\right],\left[X_{1},\left[X_{1}, X_{2}\right]\right]
$$

are linearly independent at $q$ and $\left[X_{2},\left[X_{1}, X_{2}\right]\right]$ is a functional linear combination of $X_{1}, X_{2}$ and $\left[X_{1}, X_{2}\right]$ on $U_{q}$.

Proof. $\quad \mathcal{E}_{M, q}$ denotes the $\mathbb{R}$-algebra of $C^{\infty}$ function germs at $q$ on $M$. Let $X_{1}, X_{2}$ be any local frame of $\mathcal{D}$ around $q$. We may suppose that $X_{1}, X_{2},\left[X_{1}, X_{2}\right],\left[X_{1},\left[X_{1}, X_{2}\right]\right]$ are linearly independent at $q$. From the assumption, $\left[X_{2},\left[X_{1}, X_{2}\right]\right] \in\left\langle X_{1}, X_{2},\left[X_{1}, X_{2}\right],\left[X_{1},\left[X_{1}, X_{2}\right]\right]\right\rangle_{\mathcal{E}_{M, q}}$. Then there exists a $\lambda \in \mathcal{E}_{M, q}$ such that

$$
\left[X_{2},\left[X_{1}, X_{2}\right]\right] \equiv \lambda\left[X_{1},\left[X_{1}, X_{2}\right]\right] \quad \bmod \left\langle X_{1}, X_{2},\left[X_{1}, X_{2}\right]\right\rangle_{\mathcal{E}_{M, q}}
$$

Set $\tilde{X}_{2}=X_{2}-\lambda X_{1}$. Then $\left(X_{1}, \tilde{X}_{2}\right)$ is a local frame of $\mathcal{D}$ around $q$. Then

$$
\begin{aligned}
{\left[\tilde{X}_{2},\left[X_{1}, \tilde{X}_{2}\right]\right] } & =\left[\tilde{X}_{2},\left[X_{1}, X_{2}-\lambda X_{1}\right]\right] \\
& =\left[\tilde{X}_{2},\left[X_{1}, X_{2}\right]\right]-\left[\tilde{X}_{2}, X_{1}(\lambda) X_{1}\right] \\
& =\left[\tilde{X}_{2},\left[X_{1}, X_{2}\right]\right]-X_{1}(\lambda)\left[\tilde{X}_{2}, X_{1}\right]-\tilde{X}_{2}\left(X_{1}(\lambda)\right) X_{1} \\
& \equiv\left[\tilde{X}_{2},\left[X_{1}, X_{2}\right]\right]=\left[X_{2}-\lambda X_{1},\left[X_{1}, X_{2}\right]\right] \\
& \equiv 0 \quad \bmod \left\langle X_{1}, \tilde{X}_{2},\left[X_{1}, \tilde{X}_{2}\right]\right\rangle_{\mathcal{E}_{M, q}} .
\end{aligned}
$$

For functions

$$
g_{11}:=g\left(X_{1}, X_{1}\right), g_{12}:=g\left(X_{1}, \tilde{X}_{2}\right) \text { and } g_{22}:=g\left(\tilde{X}_{2}, \tilde{X}_{2}\right)
$$

we set

$$
X_{1}^{\prime}=\frac{\sqrt{g_{22}}}{\sqrt{g_{11} g_{22}-g_{12}^{2}}}\left(X_{1}-\frac{g_{12}}{g_{22}} \tilde{X}_{2}\right), X_{2}^{\prime}=\frac{1}{\sqrt{g_{22}}} \tilde{X}_{2} .
$$

Then $\left(X_{1}^{\prime}, X_{2}^{\prime}\right)$ is a local orthonormal basis of $\mathcal{D}$ around $q$ i.e.,

$$
g\left(X_{1}^{\prime}, X_{1}^{\prime}\right)=1, g\left(X_{1}^{\prime}, X_{2}^{\prime}\right)=0, g\left(X_{2}^{\prime}, X_{2}^{\prime}\right)=1
$$

and

$$
X_{1}^{\prime}, X_{2}^{\prime},\left[X_{1}^{\prime}, X_{2}^{\prime}\right],\left[X_{1}^{\prime},\left[X_{1}^{\prime}, X_{2}^{\prime}\right]\right]
$$

are linearly independent since $X_{1}^{\prime}$ and $X_{2}^{\prime}$ is a functional linear combination of $X_{1}$ and $\tilde{X}_{2}$. Moreover $\left(X_{1}^{\prime}, X_{2}^{\prime}\right)$ satisfies

$$
\left[X_{2}^{\prime},\left[X_{1}^{\prime}, X_{2}^{\prime}\right]\right] \equiv 0 \quad \bmod \left\langle X_{1}^{\prime}, X_{2}^{\prime},\left[X_{1}^{\prime}, X_{2}^{\prime}\right]\right\rangle_{\mathcal{E}_{M, q}}
$$

because of the following for functions

$$
\alpha_{1}=\frac{\sqrt{g_{22}}}{\sqrt{g_{11} g_{22}-g_{12}^{2}}}, \alpha_{2}=-\frac{\sqrt{g_{22}}}{\sqrt{g_{11} g_{22}-g_{12}^{2}}} \frac{g_{12}}{g_{22}} \text { and } \alpha_{3}=\frac{1}{\sqrt{g_{22}}} ;
$$

$$
\begin{aligned}
& {\left[X_{2}^{\prime},\left[X_{1}^{\prime}, X_{2}^{\prime}\right]\right]=} {\left[\alpha_{3} \tilde{X}_{2},\left[\alpha_{1} X_{1}, \alpha_{3} \tilde{X}_{2}\right]\right]-} \\
&= {\left[\alpha_{3} \tilde{X}_{2},\left[\alpha_{2} \tilde{X}_{2}, \alpha_{3} \tilde{X}_{2}\right]\right] } \\
& \quad-\left[\alpha_{3}, \alpha_{1} \alpha_{3}\left[X_{1}, \tilde{X}_{2}\right]\right]+\left[\alpha_{3} \tilde{X}_{2}, \alpha_{1} X_{1}\left(\alpha_{3}\right) \tilde{X}_{2}\left(\alpha_{1}\right) X_{1}\right] \\
& \equiv {\left[\alpha_{3} \tilde{X}_{2}, \alpha_{1} \alpha_{3}\left[X_{1}, \tilde{X}_{2}\right]\right] \quad \bmod \left\langle X_{1}^{\prime}, X_{2}^{\prime},\left[X_{1}^{\prime}, X_{2}^{\prime}\right]\right\rangle_{\mathcal{E}_{M, q}} } \\
& \equiv \alpha_{1} \alpha_{3}^{2}\left[\tilde{X}_{2},\left[X_{1}, \tilde{X}_{2}\right]\right] \quad \bmod \left\langle X_{1}^{\prime}, X_{2}^{\prime},\left[X_{1}^{\prime}, X_{2}^{\prime}\right]\right\rangle_{\mathcal{E}_{M, q}} \\
& \equiv 0 \bmod \left\langle X_{1}^{\prime}, X_{2}^{\prime},\left[X_{1}^{\prime}, X_{2}^{\prime}\right]\right\rangle_{\mathcal{E}_{M, q}}
\end{aligned}
$$

Q.E.D.

Let $\left\{X_{1}, X_{2}\right\}$ be a local frame of $\mathcal{D}$ on $U_{q}$ for any $q \in M$ with the property of Lemma 4.1 and define a function $H: T^{*} U_{q} \times_{U_{q}} \mathcal{D} \rightarrow \mathbb{R}$ for the distribution $\mathcal{D}$ locally by

$$
H(x, p, u)=u_{1}\left\langle p, X_{1}(x)\right\rangle+u_{2}\left\langle p, X_{2}(x)\right\rangle .
$$

For functions $a_{1}(x, p):=\left\langle p, X_{1}(x)\right\rangle$ and $a_{2}(x, p):=\left\langle p, X_{2}(x)\right\rangle$, Proposition 3.2-(2)-(a) can be applied and we obtain the following

Proposition 4.2. For a rank 2 distribution $\mathcal{D}$ with small growth vector $(2,3,4, \ldots)$ at each point $q$ in $M$, there exist an open neighborhood $U_{q}$ of $q$, a frame $\left\{X_{1}, X_{2}\right\}$ of $\mathcal{D}$ on $U_{q}$ and an abnormal bi-extremal $(x(t), p(t))$ in $T^{*} U_{q} \backslash\{o\}$ such that

$$
\dot{x}(t)=X_{2}(x(t)), \quad \dot{p}(t)=-\frac{\partial\left\langle p, X_{2}(x)\right\rangle}{\partial x}(x(t), p(t))
$$

and

$$
\begin{array}{cl}
\left\langle p(t), X_{1}(x(t))\right\rangle=0, & \left\langle p(t), X_{2}(x(t))\right\rangle=0 \\
\left\langle p(t),\left[X_{1}, X_{2}\right](x(t))\right\rangle=0, & \left\langle p(t),\left[X_{1},\left[X_{1}, X_{2}\right]\right](x(t))\right\rangle=0
\end{array}
$$

Proof. From the property in Lemma 4.1 we take a local frame $\left\{X_{1}, X_{2}\right\}$ of $\mathcal{D}$ on an open neighborhood $U_{q}$ of $q$ in $M$ such that

$$
X_{1}, X_{2},\left[X_{1}, X_{2}\right],\left[X_{1},\left[X_{1}, X_{2}\right]\right]
$$

are linearly independent and $\left[X_{2},\left[X_{1}, X_{2}\right]\right]$ is a functional linear combination of $X_{1}, X_{2}$ and $\left[X_{1}, X_{2}\right]$. We take the function $H: T^{*} U_{q} \times_{U_{q}} \mathcal{D} \rightarrow \mathbb{R}$ as

$$
H(x, p, u)=\left\langle p, X_{1}(x)\right\rangle u_{1}+\left\langle p, X_{2}(x)\right\rangle u_{2}
$$

and set $a_{1}(x, p):=\left\langle p, X_{1}(x)\right\rangle$ and $a_{2}(x, p):=\left\langle p, X_{2}(x)\right\rangle$. Then $H$ is a Morse family because $X_{1}$ and $X_{2}$ are linearly independent at each point.

In conformity with the property of vector fields $X_{1}$ and $X_{2}$, functions $a_{1}, a_{2},\left\{a_{1}, a_{2}\right\}$ and $\left\{a_{1},\left\{a_{1}, a_{2}\right\}\right\}$ are independent and we have

$$
\left.a_{1}, a_{2},\left\{a_{1}, a_{2}\right\},\left\{a_{2},\left\{a_{1}, a_{2}\right\}\right\} \in\left\langle a_{1}, a_{2},\left\{a_{1},\left\{a_{1}, a_{2}\right\}\right\}\right\rangle\right\rangle_{\mathcal{E}_{T^{*} U_{q}, Q}}
$$

for any $Q \in T^{*} U_{q}$.
According to the proof of Proposition 3.2, $X_{u}$ is a tangent vector field to a submanifold

$$
\begin{aligned}
A_{2}^{1}=\left\{(x, p) \in T^{*} U_{q} \mid a_{1}(x, p)=a_{2}(x, p)\right. & =\left\{a_{1}, a_{2}\right\}(x, p) \\
& \left.=\left\{a_{1},\left\{a_{1}, a_{2}\right\}\right\}(x, p)=0\right\}
\end{aligned}
$$

of $T^{*} U_{q}$ for $u=(0,1)$. Thus there exists an integral curve $(x(t), p(t))$ of $X_{u}$ starting from a point in $A_{2}^{1}$, that is, $(x(t), p(t))$ satisfies ordinary differential equations

$$
\begin{aligned}
\dot{x}(t) & =\frac{\partial H}{\partial p}(x(t), p(t),(0,1))=\frac{\partial\left\langle p, X_{2}(x)\right\rangle}{\partial p}(x(t), p(t))=X_{2}(x(t)) \\
\dot{p}(t) & =-\frac{\partial H}{\partial x}(x(t), p(t),(0,1))=-\frac{\partial\left\langle p, X_{2}(x)\right\rangle}{\partial x}(x(t), p(t))
\end{aligned}
$$

and following conditions;

$$
\begin{array}{cl}
a_{1}(x(t), p(t))=\left\langle p(t), X_{1}(x(t))\right\rangle=0, & \left\langle p(t), X_{2}(x(t))\right\rangle=0, \\
\left\langle p(t),\left[X_{1}, X_{2}\right](x(t))\right\rangle=0, & \left\langle p(t),\left[X_{1},\left[X_{1}, X_{2}\right]\right](x(t))\right\rangle=0 .
\end{array}
$$

Q.E.D.

Theorem 4.3. Let $M$ be a smooth manifold and $\mathcal{D}$ be a rank two distribution. Suppose that the distribution $\mathcal{D}$ has small growth vector $(2,3,4, \ldots)$ everywhere in an open neighborhood of any point $q$ in $M$. Then for any point $q$ in $M$, there exist an open neighborhood $U_{q}$ of $q$ in $M$ and a $C^{\infty}$ immersive singular curve $x(t)$ in $U_{q}$ which is defined on a small interval.

Proof. From Proposition 4.2, for any point $q$ in $M$, there exist an open neighborhood $U_{q}$ of $q$ and an abnormal bi-extremal $(x(t), p(t))$ on $A_{2}^{1} \subset T^{*} U_{q}$. Therefore the projection of $(x(t), p(t))$ by canonical projection $\pi_{M}: T^{*} M \rightarrow M$ is a singular curve $x(t)$ with admissible velocity directed to $X_{2}$ in $\pi_{M}\left(A_{2}^{1}\right)$.
Q.E.D.

We now prove Theorem 1.
Proof of Theorem 1. Let $q$ be a point of $M$ and let $x(t)$ be the $C^{\infty}$ immersive singular horizontal curve in a neighborhood $U_{q}$ of $q$ obtained
in Theorem 4.3. Let $(x(t), p(t))$ be the abnormal bi-extremal considered in the proof of Theorem 4.3 which is obtained in Proposition 4.2. We are going to prove that this curve $x(t)$ is not a normal extremal. From Lemma 4.1 we may take a local orthonormal frame $\left\{X_{1}, X_{2}\right\}$ of $\mathcal{D}$ on $U_{q}$. We consider the Hamiltonian function in terms of the orthonormal frame $\left\{X_{1}, X_{2}\right\}$;

$$
H_{E}(x, p)=-\frac{1}{2} \sum_{i=1}^{2}\left\langle p, X_{i}(x)\right\rangle^{2}
$$

Suppose that $x(t)$ is a normal extremal. Then there must exist a normal bi-extremal of the form $(x(t), \tilde{p}(t))$ which satisfies the following differential equation;

$$
\begin{aligned}
\dot{x}(t) & =\frac{\partial H_{E}}{\partial p}(x(t), \tilde{p}(t))=-\sum_{i=1}^{2} X_{i}(x(t)) \\
\dot{\tilde{p}}(t) & =-\frac{\partial H_{E}}{\partial x}(x(t), \tilde{p}(t))=\sum_{i=1}^{2} \frac{\partial\left\langle p, X_{i}(x)\right\rangle}{\partial x}(x(t), \tilde{p}(t)) .
\end{aligned}
$$

Since the abnormal extremal $x(t)$ satisfies $\dot{x}(t)=X_{2}(t)$ by Proposition 4.2,

$$
X_{2}(x(t))=\dot{x}(t)=-X_{1}(x(t))-X_{2}(x(t))
$$

holds. Thus $X_{1}(x(t))+2 X_{2}(x(t))=0$ holds. This is a contradiction to $\left\{X_{1}, X_{2}\right\}$ being a local frame of $\mathcal{D}$.
Q.E.D.

Remark 4.4 ([2], Theorem 2.8). It is known that there are no singular minimizer for a generic sub-Riemannian manifold (the genericity is used for the distribution as map-germs) with rank greater than 2.

Remark 4.5 ([7], Proposition 11). There is a result given by LiuSussmann which is similar type to Theorem 1, however the method for the proof is different from ours.

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