# On the Normal Forms of Differential Equations in the Neighborhood of an Equilibrium Point 

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## § 1. Introduction.

1. In this note we use the notations $\partial_{i} u$ and $\partial_{i j}^{2} u$ for $\frac{\partial}{\partial x_{i}} u$ and $\frac{\partial^{2}}{\partial x_{i} \partial x_{j}} u$ respectively. The vectors $\left(x_{1}, \cdots, x_{m}\right)$ and $\left(y_{1}, \cdots, y_{m}\right)$ in $R^{m}$ will be denoted briefly by $x$ and $y$ respectively.

Let $A=\left(a_{i j}\right)$ be a constant real ( $m, m$ )-matrix, all of whose characteristic roots $\lambda_{i}(i=1, \cdots, m)$ have non-zero real parts, and $f(x)=\left(f_{1}(x), \cdots, f_{m}(x)\right)$ a real vector function of class $C^{1}$ on some neighborhood of $x=0$, such that $f(0)=0$ and $\left|\partial_{x} f(x)\right| \leqq K \cdot|x|$ with a constant $K>0$ where

$$
|x|=\left(\sum_{i} x_{i}^{2}\right)^{\frac{1}{2}}, \quad\left|\partial_{x} f(x)\right|=\left\{\sum_{i, j}\left(\partial_{i} f_{j}(x)\right)^{2}\right\}^{\frac{1}{2}}
$$

We consider the autonomous systems

$$
\begin{equation*}
\frac{d x}{d t}=A \cdot x+f(x) \tag{1.1}
\end{equation*}
$$

and

$$
\begin{equation*}
\frac{d y}{d t}=A \cdot y \tag{1.2}
\end{equation*}
$$

regarding $x, y$ and $f(x)$ as the column-vectors. The purpose of this note is to show that, under some conditions on $\lambda_{i}(i=1, \cdots, m)$ and $f(x)$, the system (1.1) can be transformed into (1.2) by a change of variables

$$
\begin{equation*}
y=x+u(x) \tag{1.3}
\end{equation*}
$$

where $u(x)=\left(u_{1}(x), \cdots, u_{m}(x)\right)$ is a real vector function of class $C^{1}$, such that

$$
\left\{\begin{array}{l}
u(0)=0  \tag{1.4}\\
\left|\partial_{x} u(x)\right| \leqq L \cdot|x|
\end{array}\right.
$$

with some constant $L>0$.

When $f(x)$ is analytic regular in $x$, in order to show the existence of the transformation given by (1.3) with analytic regular $u(x)$, we must necessarily assume that there exist no relations of the form

$$
\begin{equation*}
\lambda_{i}=\sum_{j=1}^{m} n_{j} \cdot \lambda_{j} \tag{1.5}
\end{equation*}
$$

where $n_{j}(j=1, \cdots, m)$ are non-negative integers such that $\sum_{j=1}^{m} n_{j}>1$. As to this case, some results were obtained by H. Poincaré, C. L. Siegel, and others, while we obtain the present result for the real systems with a transformation of class $C^{1}$ under some weaker conditions.

## § 2. Main Theorem.

2. Theorem. Assumptions:
(i) $A$ is a constant real $(m, m)$-matrix, all of whose characteristic roots $\lambda_{i}(i=1, \cdots, m)$ have non-zero real parts: $\mathfrak{R}\left(\lambda_{i}\right) \neq 0(i=1, \cdots, m)$.
(ii) Let

$$
\begin{equation*}
f_{i}(x)=p_{i}(x)+q_{i}(x) \quad(i=1, \cdots, m) \tag{2.1}
\end{equation*}
$$

where $p_{i}(x)$ are polynomials in $x$ with real coefficients such that $p_{i}(0)=\partial_{j_{-}} 力_{i}(0)=0(i=1, \cdots, m ; j=1, \cdots, m)$, and $q_{i}(x)(i=1, \cdots, m)$ are real-valued functions of class $C^{1}$ satisfying

$$
\left\{\begin{array}{l}
q(0)=0  \tag{2.2}\\
\left|\partial_{x} q(x)\right| \leqq Q \cdot|x|^{n}
\end{array}\right.
$$

with some integer $h>0$ and some constant $Q>0$.
(iii) There exist no relations of the form

$$
\lambda_{i}=\sum_{j=1}^{m} n_{j} \cdot \lambda_{j}
$$

where $n_{j}(j=1, \cdots, m)$ are non-negative integers such that

$$
h>\sum_{j=1}^{m} n_{j}>1
$$

Conclusion: There exists a positive constant $h_{0}$, depending only on $\lambda_{i}(i=1, \cdots, m)$, with the following property: if $h>h_{0}$, there exist functions $u_{i}(x)(i=1, \cdots, m)$ of class $C^{1}$ satisfying (1.4), such that the system (1.1) is reduced to the form (1.2) by the substitution (1.3).
3. If (1.1) is transformed into (1.2) by (1.3), $u(x)$ must satisfy the system of partial differential equations

$$
\begin{equation*}
\sum_{i=1}^{m}\left(\sum_{j=1}^{m} a_{i j} x_{j}+f_{i}(x)\right) \cdot \partial_{i} u_{\nu}=\sum_{\mu=1}^{m} a_{\nu \mu} u_{\mu}-f_{\nu}(x) \quad(\nu=1, \cdots, m) . \tag{3.1}
\end{equation*}
$$

For we have, by operating $\frac{d}{d t}$ on both sides of (1.3),

$$
\frac{d y_{i}}{d t}=\frac{d x_{i}}{d t}+\sum_{\nu=1}^{m} \partial_{\nu} u_{i} \cdot \frac{d x_{\nu}}{d t} \quad(i=1, \cdots, m)
$$

from which (3.1) follows immediately by (1.1), (1.2) and (1.3). Conversely, if $u(x)$ is any function satisfying (3.1), then the substitution (1.3) will transform (1.1) into (1.2). Thus we have only to show the existence of $u(x)$ satisfying (1.4) and (3.1), if $h$ is sufficiently large.

## § 3. Auxiliary Theorem.

4. In this section we consider the system of semi-linear partial differential equations

$$
\begin{equation*}
\sum_{i=1}^{m} P_{i}(x) \cdot \partial_{i} u_{\nu}=Q_{\nu}(x, u) \quad(\nu=1, \cdots, l) \tag{4.1}
\end{equation*}
$$

where $x=\left(x_{1}, \cdots, x_{m}\right)$ and $u=\left(u_{1}, \cdots, u_{l}\right)$ denote real vectors in $R^{m}$ and $R^{l}$ respectively. Let $P_{i}(x)$ be real-valued functions of class $C^{1}$ in an open domain $D \subset R^{m}$, such that

$$
\begin{equation*}
\left(P_{1}(x), \cdots, P_{m}(x)\right) \neq(0, \cdots, 0) \quad(x \in D) . \tag{4.2}
\end{equation*}
$$

And $Q_{\nu}(x, u)$ be real-valued functions of class $C^{1}$ in

$$
\Omega=\left\{(x, u) \in R^{m+l}: x \in D,|u| \leqq \omega(x)\right\}
$$

where $\omega(x)$ is some positive-valued function of class $C^{1}$ in $D$. A curve $x=x(t)$ in $R^{m}$ is said to be a base characteristic of (4.1) if $x(t)$ satisfies the following system of ordinary differential equations:

$$
\begin{equation*}
\frac{d x_{i}}{d t}=P_{i}(x) \quad(i=1, \cdots, m) \tag{4.3}
\end{equation*}
$$

Let an ( $m-1$ )-dimensional manifold $M$ in $R^{m}$ be given by

$$
\begin{equation*}
M: x_{i}=A_{i}\left(s_{1}, \cdots, s_{m-1}\right) \quad(i=1, \cdots, m) \tag{4.4}
\end{equation*}
$$

where $A_{i}(s)$ are functions of class $C^{1}$ in some domain $S \subset R^{m-1}$ such that $A(s)=\left(A_{1}(s), \cdots, A_{m}(s)\right) \in D$ for $s=\left(s_{1}, \cdots, s_{m-1}\right) \in S$. We assume that

$$
\left|P_{i}(A(s)), \partial_{j} A_{i}(s) \quad \begin{array}{l}
i \downarrow 1, \cdots, m  \tag{4.5}\\
j \rightarrow 1, \cdots, m-1
\end{array}\right|^{1)} \neq 0 \quad \text { for } \quad s \in S
$$

and that any base characteristic

$$
\begin{equation*}
x=x(t, s) \tag{4.6}
\end{equation*}
$$

issuing from a point of $M$ so that $x(0, s)=A(s)$, exists on the interval: $0 \leqq t<\tau(s)$ where $\tau(s)$ is a continuous function on $S$, and that the set $X=\{x=x(t, s): 0 \leqq t<\tau(s), s \in S\}$ is filled up only onefold with all those curves $x=x(t, s)(s \in S)$, i.e. to any point $x \in X$ there corresponds just one $(t, s)$ such that $x=x(t, s), 0 \leqq t<\tau(s), s \in S$. Then we have easily

$$
\left.\left.\begin{array}{rl} 
& \frac{\partial\left(x_{1}, \cdots, x_{m}\right)}{\partial\left(t, s_{1}, \cdots, s_{m-1}\right)} \\
= & \mid P_{i}(A(s)), \partial_{j} A_{i}(s) \\
& i \downarrow 1, \cdots, m \\
j \rightarrow 1, \cdots, m-1
\end{array} \right\rvert\, \cdot \exp \left(\int_{0}^{t} \sum_{i=1}^{m} \partial_{i} P_{i}(x)_{x=x x t}, s\right) d t\right) \neq 0,
$$

which shows that the $1-1$ mapping (4.6) from $\{(t, s): s \in S, 0 \leqq t<\tau(s)\}$ onto $X$ and its inverse are both of class $C^{1}$.

By (4.6) the system (4.1) is reduced to the following system of ordinary differential equations, $s$ being a parameter :

$$
\begin{equation*}
\frac{d u_{\nu}}{d t}=Q_{\nu}(x(t, s), u) \quad(\nu=1, \cdots, l) \tag{4.7}
\end{equation*}
$$

We have then

$$
\partial_{t} \omega(x(t, s))=\left[\sum_{i=1}^{m} P_{i}(x) \cdot \partial_{i} \omega(x)\right]_{x=x(t, s)}
$$

and

$$
\partial_{t}|u(t, s)| \cdot|u(t, s)|=\sum_{\nu=1}^{t} u_{\nu}(t, s) \cdot Q_{\nu}(x(t, s), u(t, s))
$$

for any solution $u(t, s)$ of (4.7). Hence, we obtain easily the following auxiliary theorem which is our principal tool.

Auxiliary theorem. Under the conditions mentioned above, let the inequality

$$
\begin{equation*}
\sum_{i=1}^{m} P_{i}(x) \partial_{i} \omega(x) \geqq \frac{1}{\omega(x)} \sum_{\nu=1}^{m} Q_{\nu}(x, u) \cdot u_{\nu} \tag{4.8}
\end{equation*}
$$

hold for any $x \in X$ such that $|u|=\omega(x)$. Then, for any function $B(s)=\left(B_{1}(s), \cdots, B_{l}(s)\right)$ of class $C^{1}$ on $S$ such that

1) an ( $m, m$ )-determinant.

$$
\begin{equation*}
|B(s)| \leqq \omega(A(s)) \tag{4.9}
\end{equation*}
$$

there exists a unique solution $u(x)$ of (4.1) on $X$, such that

$$
u(A(s))=B(s)
$$

and

$$
\begin{equation*}
|u(x)| \leqq \omega(x) \tag{4.10}
\end{equation*}
$$

for $x \in X$.

## § 4. Estimation of $\boldsymbol{u}(\boldsymbol{x})$.

5. Consider the system of partial differential equations

$$
\begin{equation*}
\sum_{i=1}^{m}\left(\sum_{j=1}^{m} a_{i j} x_{j}+f_{i}(x)\right) \partial_{i} u_{\nu}=\sum_{\mu=1}^{m} a_{\nu \mu} \cdot u_{\mu}+g_{\nu}(x) \quad(\nu=1, \cdots, m) \tag{5.1}
\end{equation*}
$$

for which we have the following lemma.
Lemma. Let $A=\left(a_{i j}\right)$ and $f(x)$ satisfy the assumptions (i), (ii) and (iii) in the theorem. Let $g_{\nu}(x)(\nu=1, \cdots, m)$ be real-valued functions of class $C^{1}$ on some neighborhood of 0 , such that

$$
\left\{\begin{array}{l}
g(0)=0  \tag{5.2}\\
\left|\partial_{x} g(x)\right| \leqq G|x|^{p}
\end{array}\right.
$$

where $G$ and $p$ are positive constants.
Then there exists a constant $h_{0}>0$, which depends only on $\lambda_{i}(i=1, \cdots, m)$, with the following property: if $p>h_{0}$, the system (5.1) has a unique solution $u(x)$ in a neighborhood of 0 , such that

$$
\begin{equation*}
u(x)=0 \quad \text { on the cone } \sum_{i=1}^{k} x_{i}^{2}=\sum_{i=k+1}^{m} x_{i}^{2}{ }^{2)} \tag{5.3}
\end{equation*}
$$

and

$$
\begin{equation*}
\left|\partial_{x} u(x)\right| \leqq C \cdot G|x|^{p} \tag{5.4}
\end{equation*}
$$

where $C$ is a positive constant depending only on $\lambda_{i}(i=1, \cdots, m)$ and $p$.
6. Proof. By setting $P_{i}(x)=\sum_{j=1}^{m} a_{i j} x_{i}+f_{i}(x)$ and $Q_{\nu}(x, u)=\sum_{\mu=1}^{m} a_{\nu \mu} u_{\mu}$ $+g_{\nu}(x)$, the system (5.1) has the form (4.1). Without loss of generality we assume that $A=\left(a_{i j}\right)$ has the following form:

$$
\begin{array}{ll}
a_{i i}=\Re\left(\lambda_{i}\right) & (i=1, \cdots, m) \\
\Re\left(\lambda_{i}\right)>0 & \text { for } \quad i \leqq k  \tag{i}\\
\Re\left(\lambda_{i}\right)<0 & \text { for } \quad i>k
\end{array}
$$

[^0](ii)
\[

$$
\begin{array}{ll}
a_{i j}=0 & \text { for } i \leqq k \text { and } j>k, \\
a_{i j}=0 & \text { for } i>k \text { and } j \leqq k .
\end{array}
$$
\]

(iii)

$$
\begin{equation*}
\left|\sum_{i \neq j} a_{i j} x_{i} x_{j}\right| \leqq \delta|x|^{2} \tag{6.1}
\end{equation*}
$$

where $\delta$ is any prescribed positive number.
In what follows, we write $\sum_{\alpha}=\sum_{\alpha=1}^{k}$ and $\sum_{\beta}=\sum_{\beta=k+1}^{m}$. We suppose that $f$ and $g$ are functions of class $C^{1}$ on $U=\left\{x: \sum_{\alpha} x_{\alpha}^{2} \leqq r^{2}, \sum_{\beta} x_{\beta}^{2} \leqq r^{2}\right\}$ where $r$ is a positive constant. We consider the case $0<k<m$. Because, if $k=0$ or $k=m$, the proof of the lemma will be simpler.

With sufficiently small $\varepsilon>0$ we set ${ }^{3)}$

$$
S_{\varepsilon}(x)=\left\{\begin{array}{l}
\sum_{\alpha} x_{\alpha}^{2}-\sum_{\beta} x_{\beta}^{2} \quad \text { when } \quad \sum_{\alpha} x_{\alpha}^{2} \geq \varepsilon^{2}  \tag{6.2}\\
\sum_{\alpha} x_{\alpha}^{2}-\sum_{\beta} x_{\beta}^{2}-\frac{1}{2 \varepsilon^{2}}\left(\varepsilon^{2}-\sum_{\alpha} x_{\alpha}^{2}\right)^{2} \quad \text { when } \quad \sum_{\alpha} x_{\alpha}^{2}<\varepsilon^{2}
\end{array}\right.
$$

and define a bounded region $U_{\varepsilon}$ by $U_{\varepsilon}=\left\{x \in U: S_{\varepsilon}(x) \geq 0\right\}$.


$$
\begin{aligned}
& \mathrm{N}=(r, 0) \\
& \mathrm{P}=\left(\frac{\sqrt{6}-\sqrt{2}}{2} \varepsilon, 0\right) \\
& \mathrm{Q}=(\varepsilon, 0)
\end{aligned}
$$

First we consider the solution of (5.1) in $U_{\varepsilon}$ vanishing on the ( $m-1$ )-dimensional manifold $M_{\varepsilon}=\left\{x \in U: S_{\varepsilon}(x)=0\right\}$. For the base characteristic $x=x(t)$ of (5.1), we have

$$
\begin{aligned}
& \frac{1}{2} \cdot \frac{d}{d t} S_{\varepsilon}(x(t))=\sum_{\alpha} x_{\alpha} \cdot \dot{x}_{\alpha}-\sum_{\beta} x_{\beta} \cdot \dot{x}_{\beta} \\
& \quad=\sum_{\alpha}\left(\sum_{i=1}^{m} a_{\alpha i} x_{i}+f_{a}(x)\right) x_{a}-\sum_{\beta}\left(\sum_{i=1}^{m} a_{\beta i} x_{i}+f_{\beta}(x)\right) x_{\beta} \\
& \quad=\left(\sum_{\alpha} \sum_{i=1}^{m} a_{\alpha i} x_{\alpha i} x_{i}-\sum_{\beta} \sum_{i=1}^{m} a_{\beta i} x_{\beta} x_{i}\right)+\left(\sum_{\alpha} f_{\alpha}(x) x_{\alpha}-\sum_{\beta} f_{\beta}(x) x_{\beta}\right)>0
\end{aligned}
$$

[^1]when $\varepsilon^{2}<\sum_{\alpha} x_{\alpha}^{2} \leqq r^{2}$, by taking $r$ small enough, and also when $\sum_{\alpha} x_{\alpha}^{2} \leqq \varepsilon^{2}$
$$
\frac{1}{2} \cdot \frac{d}{d t} S_{\varepsilon}(x(t))=\left\{1+\frac{1}{\varepsilon^{2}}\left(\varepsilon^{2}-\sum_{\alpha} x_{\alpha}^{2}\right)\right\} \sum_{\alpha} x_{\infty} \cdot \dot{x}_{\omega}-\sum_{\beta} x_{\beta} \cdot \dot{x}_{\beta}>0 .
$$

From these inequalities we see that, if $r>0$ is taken small enough, every base characteristic of (5.1) meeting $M_{\varepsilon}$ is transverse to $M_{\varepsilon}$, and that (4.5) will hold for this case with $M=M_{\varepsilon}$. In addition, since we have

$$
\frac{1}{2} \frac{d}{d t} \sum_{\alpha}\left(x_{\alpha}(t)\right)^{2}=\sum_{\alpha}\left(\sum_{i} a_{\alpha i} x_{i}+f_{\infty}(x)\right) x_{\omega}>0
$$

for any base characteristic $x(t)$, when $r$ is small enough, we see that $U_{\mathrm{z}}$ is filled up only onefold with the base characteristics issuing from $M_{\mathrm{z}}$. Therefore, we apply the auxiliary theorem to this case, setting $U_{\mathrm{e}}=X$.

We set

$$
\begin{equation*}
\varphi(x)=(1+\gamma) \sum_{\alpha} x_{\alpha}^{2}-\sum_{\beta} x_{\beta}^{2} \tag{6.3}
\end{equation*}
$$

and

$$
\begin{equation*}
\omega(x)=W \cdot G \cdot \varphi(x)^{\frac{p+1}{2}} \tag{6.4}
\end{equation*}
$$

where $\gamma>0$ and $W>0$ will be determined later. Then

$$
\begin{equation*}
\frac{\gamma}{2}|x|^{2} \leqq \varphi(x) \leqq(1+\gamma)|x|^{2} \tag{6.5}
\end{equation*}
$$

for every $x \in U_{\mathrm{e}}$. Thus we obtain

$$
\begin{aligned}
& \sum_{i=1}^{m} P_{i}(x) \partial_{i} \omega(x) \\
& =(p+1) \cdot W \cdot G \cdot \varphi(x)^{\frac{p-1}{2}}\left\{(1+\gamma) \sum_{\alpha} \sum_{j} a_{a j} x_{\alpha} x_{j}-\sum_{\beta} \sum_{j} a_{\beta j} x_{\beta} x_{j}\right. \\
& \\
& \left.+(1+\gamma) \sum_{\alpha} x_{\alpha} f_{\alpha}(x)-\sum_{\beta} x_{\beta} f_{\beta}(x)\right\}
\end{aligned}
$$

and

$$
\sum_{\nu=1}^{m} Q_{\nu}(x, u) \cdot u_{\nu}=\sum_{\nu=1}^{m} \sum_{\mu=1}^{m} a_{\nu \mu} u_{\nu} u_{\mu}-\sum_{\nu=1}^{m} g_{\nu}(x) u_{\nu} .
$$

Now we set

$$
\begin{equation*}
\Lambda_{0}=\min _{1 \leqq i \leqq m}\left|\Re\left(\lambda_{i}\right)\right|, \quad \Lambda_{1}=\max _{1 \leqq i \leqq^{m}}\left|\Re\left(\lambda_{i}\right)\right| \tag{6.6}
\end{equation*}
$$

Then we have

$$
\begin{equation*}
\sum_{i=1}^{m} P_{i}(x) \cdot \partial_{i} \omega(x)>(p+1)\left(\Lambda_{0}-2 \delta\right) W \cdot G|x|^{2} \mathcal{P}(x)^{\frac{p-1}{2}} \tag{6.7}
\end{equation*}
$$

and

$$
\begin{equation*}
\sum_{\nu=1}^{m} Q_{\nu}(x, u) \cdot u_{\nu}<\left(\Lambda_{1}+\delta\right)|u|^{2}+|g(x)| \cdot|u| \tag{6.8}
\end{equation*}
$$

on $U_{\mathrm{s}}$ where $\delta$ is given by (6.1), taking $r$ small enough. From (6.8) it follows that, if $|u|=\omega(x)$ for some $x \in U_{\varepsilon}$, then

$$
\frac{1}{\omega(x)} \sum_{\nu=1}^{m} Q_{\nu}(x, u) u_{\nu}<\left(\Lambda_{1}+\delta\right) \omega(x)+G|x|^{p+1}
$$

and so

$$
\begin{equation*}
\frac{1}{\omega(x)} \sum_{\nu=1}^{m} Q_{\nu}(x, u) u_{\nu}<W \cdot G\left(\Lambda_{1}+\delta\right)\left\{1+\gamma+\frac{\left(\frac{2}{\gamma}\right)^{\frac{p-1}{2}}}{\Lambda_{1} \cdot W}\right\}|x|^{2} \mathcal{P}(x)^{\frac{p-1}{2}} \tag{6.9}
\end{equation*}
$$

by virtue of (6.4) and (6.5). Thus, if we assume

$$
\begin{equation*}
(p+1) \Lambda_{0}>\Lambda_{1} \tag{6.10}
\end{equation*}
$$

we have

$$
\begin{equation*}
\sum_{i=1}^{m} P_{i}(x) \partial_{i} \omega(x)>\frac{1}{\omega(x)} \sum_{\nu=1}^{m} Q_{\nu}(x, u) \cdot u_{\nu} \tag{6.11}
\end{equation*}
$$

for any $x \in U_{\varepsilon}$ such that $|u|=\omega(x)$, from (6.7) and (6.9), by taking $\delta>0$ and $\gamma>0$ small enough and then $W$ large enough.

Let us assume hereafter that (6.10) holds. Then it follows from the auxiliary theorem that, when $r$ is small enough, there exists a unique solution $u(x ; \varepsilon)$ of (5.1) on $U_{\varepsilon}$ which vanishes on $M_{\varepsilon}$, and that it satisfies

$$
\begin{equation*}
|u(x ; \varepsilon)| \leqq \omega(x)=W \cdot G \cdot \varphi(x)^{\frac{p+1}{2}} \leqq G \cdot K \cdot|x|^{p+1} \quad \text { for } \quad x \in U_{\varepsilon} \tag{6.12}
\end{equation*}
$$

where $K=W(1+\gamma)^{\frac{p+1}{2}}$. Notice that $W$ and $r$ are taken independent of $\varepsilon>0$ in the above arguement.

## § 5. Continuation of the Proof of the Lemma. Estimation of $\left|\partial_{x} u(x)\right|$.

7. Next, let us prove that the inequality (5.4) holds for $u=u(x ; \varepsilon)$ on $U_{\mathrm{s}}$ with some constant $C>0$ independent of $\varepsilon$. In this paragraph we fix $\varepsilon>0$ and write $u_{\nu}(x)$ in place of $u_{\nu}(x ; \varepsilon)$ for simplicity.

In order to estimate $\left|\partial_{x} u(x)\right|$ on the manifold $M_{\varepsilon}$, we reduce the system (5.1) into

$$
\begin{equation*}
\frac{d u_{\nu}}{d t}=\sum_{\mu=1}^{m} a_{\nu \mu} u_{\mu}+g_{\nu}(x(t, s)) \quad(\nu=1, \cdots, m) \tag{7.1}
\end{equation*}
$$

by the change of variables given by (4.6). Thus, if we set $u_{\nu}(t, s)$ $=u_{\nu}(x(t, s)), u=u(t, s)$ is the solution of (7.1) with the initial condition

$$
\begin{equation*}
u(0, s) \equiv 0 \tag{7.2}
\end{equation*}
$$

Therefore,

$$
\left\{\begin{array}{l}
\partial_{t} u(0, s) \equiv g(x(0, s))  \tag{7.3}\\
\partial_{s} u(0, s) \equiv 0
\end{array}\right.
$$

on $M_{\varepsilon}$. Let $\partial_{n} u(s)$ denote the normal derivative of $u(t, s)$ at $x=x(0, s)$ on $M_{\varepsilon}$, then we have

$$
\begin{equation*}
\left[\frac{\partial_{t} u(t, s)}{\left\{\sum_{i=1}^{m}\left(\partial_{t} x_{i}(t, s)\right)^{2}\right\}^{\frac{1}{2}}}\right]_{t=0}=\partial_{n} u(s) \cdot \cos \theta(s) \tag{7.4}
\end{equation*}
$$

where $\theta(s)$ represents the angle between the base characteristic $x=x(t, s)$ and the normal of $M_{\mathrm{e}}$ at $x(0, s)$, i.e.

$$
\begin{equation*}
\cos \theta(s)=\frac{(1+\sigma(x)) \cdot \sum_{\alpha} x_{\alpha} \cdot \partial_{t} x_{\alpha}-\sum_{\beta} x_{\beta} \cdot \partial_{t} x_{\beta}}{\left\{\sum_{i=1}^{m}\left(\partial_{t} x_{i}\right)^{2}\right\}^{\frac{1}{2}}\left\{(1+\sigma(x))^{2} \sum_{\alpha} x_{\alpha}^{2}+\sum_{\beta} x_{\beta}^{2}\right\}^{\frac{1}{2}}} \tag{7.5}
\end{equation*}
$$

where

$$
\sigma(x)= \begin{cases}0 & \text { when } \sum_{\alpha} x_{\alpha}^{2} \geq \varepsilon^{2} \\ \frac{1}{\varepsilon^{2}}\left(\varepsilon^{2}-\sum_{\alpha} x_{\alpha}^{2}\right) & \text { when } \sum_{\alpha} x_{\alpha}^{2}<\varepsilon^{2}\end{cases}
$$

Since
(7.6)

$$
(1+\sigma(x)) \cdot \sum_{\alpha} x_{a} \cdot \partial_{t} x_{\alpha}-\sum_{\beta} x_{\beta} \cdot \partial_{t} x_{\beta}>\frac{1}{2} \Lambda_{0}|x|^{2} \quad\left(x \in U_{\varepsilon}\right)
$$

and

$$
\begin{equation*}
(1+\sigma(x))^{2} \cdot \sum_{\alpha} x_{\alpha}^{2}+\sum_{\beta} x_{\beta}^{2}<2|x|^{2} \tag{7.7}
\end{equation*}
$$

when $r$ is sufficiently small, we obtain, from (7.3), (7.4), (7.5), (7.6) and (7.7),

$$
\left|\partial_{n} u\right|<\frac{2 \sqrt{2}}{\Lambda_{0}} \cdot \frac{|g(x)|}{|x|}
$$

Hence

$$
\begin{equation*}
\left|\partial_{n} u\right|<K^{\prime} \cdot G|x|^{p} \tag{7.8}
\end{equation*}
$$

on $M_{\varepsilon}$ with some constant $K^{\prime}>0$ independent of $\varepsilon$. Thus, from (7.3) and (7.8), we see

$$
\begin{equation*}
\left|\partial_{x} u(x ; \varepsilon)\right| \leqq K^{\prime} \cdot G|x|^{p} \quad\left(x \in M_{\varepsilon}\right) \tag{7.9}
\end{equation*}
$$

Now, operating $\partial_{\mu}$ on both sides of (5.1) and setting $\partial_{\mu} u_{\nu}=u^{\nu \mu}$, we have

$$
\begin{align*}
& \sum_{i=1}^{m}\left(\sum_{j=1}^{m} a_{i j} x_{j}+f_{i}(x)\right) \partial_{i} u^{\nu \mu}  \tag{7.10}\\
= & \sum_{j=1}^{m} a_{\nu j} u^{j \mu}-\sum_{i=1}^{m}\left(a_{i \mu}+\partial_{\mu} f_{i}(x)\right) u^{\nu i}+\partial_{\mu} g_{\nu}(x) \quad(\nu, \mu=1, \cdots, m)
\end{align*}
$$

which has also the form (4.1) with unknown functions $u^{\nu \mu}$. Let us assume first that $f$ and $g$ are functions of class $C^{2}$ in $U$ and apply the auxiliary theorem to (7.10). We set

$$
\begin{aligned}
& u^{\prime}=\left(u^{1,1}, u^{1,2}, \cdots, u^{m, m}\right) \quad\left(\in R^{m 2}\right) \\
& P_{i}(x)=\sum_{j=1}^{m} a_{i j} x_{j}+f_{i}(x) \\
& Q^{\nu \mu}\left(x, u^{\prime}\right)=\sum_{j=1}^{m} a_{\nu j} u^{j \mu}-\sum_{i=1}^{m}\left(a_{i \mu}+\partial_{\mu} f_{i}(x)\right) \cdot u^{\nu i}+\partial_{\mu} g_{\nu}(x) \quad(\nu, \mu=1, \cdots, m) \\
& \text { and } \quad \omega^{\prime}(x)=W^{\prime} \cdot G \cdot \phi^{\prime}(x)^{\frac{p}{2}} \quad \text { where } \quad \varphi^{\prime}(x)=\left(1+\gamma^{\prime}\right) \sum_{\alpha} x_{\alpha}^{2}-\sum_{\beta} x_{\beta}^{2} .
\end{aligned}
$$

Then, if we assume

$$
\begin{equation*}
p \Lambda_{0}>\tilde{\Lambda} \equiv \max _{i, j}\left|\Re\left(\lambda_{i}\right)-\Re\left(\lambda_{j}\right)\right| \tag{7.11}
\end{equation*}
$$

and if we take $r$ small enough, we have

$$
\sum_{i=1}^{m} P_{i}(x) \partial_{i} \omega^{\prime}(x)>\frac{1}{\omega^{\prime}(x)} \sum_{=1}^{m} \sum_{\mu=1}^{m} Q^{\nu \mu}\left(x, u^{\prime}\right) \cdot u^{\nu \mu}
$$

for $x$ such that $\omega^{\prime}(x)=\left|u^{\prime}\right|$ in $U_{\varepsilon}$, taking $W^{\prime}$ and $\gamma^{\prime}>0$ appropriately. By the auxiliary theorem and (7.9) we thus get

$$
\begin{equation*}
\left|\partial_{x} u(x ; \varepsilon)\right| \leqq C \cdot G|x|^{p} \tag{7.12}
\end{equation*}
$$

in $\dot{U}_{\varepsilon}$ where $C>0$ is a constant independent of $\varepsilon$. In the above consideration we can also choose $r$ independent of $\varepsilon$. Let us write

$$
h_{0}=\max \left(\frac{\Lambda_{1}}{\Lambda_{0}}-1, \tilde{\Lambda}\right)
$$

and assume $p>h_{0}$ hereafter, from which (6.10) and (7.11) follow.
We will study in 9 . as to the case that $f$ and $g$ are functions of class $C^{1}$.
8. Notice that $C$ of (7.12) and $r$ can be chosen independent of $\varepsilon>0$ which is sufficiently small. Now we consider $\varepsilon$ as a variable tending to zero. We see easily $U_{\varepsilon} \subset U_{\varepsilon}^{\prime}$ as $\varepsilon>\varepsilon^{\prime}>0$, and $v=u\left(x ; \varepsilon^{\prime}\right)-u(x ; \varepsilon)$ must satisfy on $U_{\text {e }}$

$$
\begin{equation*}
\sum_{i=1}^{m}\left(\sum_{j=1}^{m} a_{i j} x_{j}+f_{i}(x)\right) \partial_{i} v_{\nu}=\sum_{\mu=1}^{m} a_{\nu \mu} v_{\mu} \quad(\nu=1, \cdots, m) . \tag{8.1}
\end{equation*}
$$

From (6.12) and (7.12) we see easily that

$$
|v| \leqq 2^{\frac{p+3}{2}} K \cdot G \cdot \varepsilon^{p+1}
$$

and

$$
\left|\partial_{x} v\right| \leqq 2^{\frac{p+2}{2}} C \cdot G \cdot \varepsilon^{p}
$$

hold for $v=u\left(x ; \varepsilon^{\prime}\right)-u(x ; \varepsilon)$ on $M_{\varepsilon} . \quad$ Notice that $\min _{x \in U_{\varepsilon}}|x|=\frac{\sqrt{6}-\sqrt{2}}{2} \varepsilon$, and we have for any $q>0$

$$
\left\{\begin{array}{c}
|v| \leqq K_{0} \cdot \varepsilon^{p-q}|x|^{q+1}  \tag{8.2}\\
\left|\partial_{x} v\right| \leqq K_{1} \cdot \varepsilon^{p-q}|x|^{q}
\end{array}\right.
$$

on $M_{\varepsilon}$ where $K_{0}$ and $K_{1}$ are constants not depending on $\varepsilon$ and $\varepsilon^{\prime}$.
We now choose $q$ so that $p>q>h_{0}$. The system (8.1) is a special case of (5.1) with $g(x) \equiv 0$, and we get similarly as (6.12) and (7.12)

$$
\left\{\begin{array}{c}
\left|u\left(x ; \varepsilon^{\prime}\right)-u(x ; \varepsilon)\right| \leqq K^{\prime} \cdot \varepsilon^{p-q}|x|^{q+1} \\
\left|\partial_{x} u\left(x ; \varepsilon^{\prime}\right)-\partial_{x} u(x ; \varepsilon)\right| \leqq K^{\prime} \cdot \varepsilon^{p-q}|x|^{q}
\end{array}\right.
$$

in $U_{\mathrm{e}}$ where $K^{\prime}$ is some positive constant independent of $\varepsilon$ and $\varepsilon^{\prime}$. Thus we see that, as $\varepsilon \rightarrow 0, u_{\nu}(x ; \varepsilon)$ and $\partial_{\mu} u_{\nu}(x ; \varepsilon)$ tend to certain functions $u_{\nu}(x)$ and their derivatives $\partial_{\mu} u_{\nu}(x)$ respectively. Clearly this $u(x)$ is a solution of (5.1), vanishing on the manifold : $\sum_{\alpha} x_{\alpha}^{2}=\sum_{\beta} x_{\beta}^{2}$ and satisfying (5.4) in $U_{1}=\left\{x: \sum_{\alpha} x_{\alpha}^{2} \leqq r^{2}, \sum_{\beta} x_{\beta}^{2} \leqq \sum_{\alpha} x_{\alpha}^{2}\right\}$.

Quite similarly as above, we can prove the existence of a solution $u(x)$ of (5.1) vanishing on $\sum_{\alpha} x_{\alpha}^{2}=\sum_{\beta} x_{\beta}^{2}$ and satisfying (5.4) in $U_{2}$ $=\left\{x: \sum_{\beta} x_{\beta}^{2} \leqq r^{2}, \sum_{\alpha} x_{\alpha}^{2} \leqq \sum_{\beta} x_{\beta}^{2}\right\}$.
9. Now it remains to prove our lemma when $f$ and $g$ are functions of class $C^{1}$. We construct approximation sequences $\left\{f^{n}(x)\right\}_{n=1}^{\infty}$ and $\left\{g^{n}(x)\right\}_{n=1}^{\infty}$ of vector functions of class $C^{2}$ such that ${ }^{4)}$

$$
\left\{\begin{array}{l}
f^{n}(0)=g^{n}(0)=0  \tag{9.1}\\
f^{n}(x)=p(x)+q^{n}(x)^{5)} \\
\left|\partial_{x} f^{n}(x)-\partial_{x} f(x)\right| \leqq \frac{1}{n}|x|^{n} \\
\left|\partial_{x} g^{n}(x)-\partial_{x} g(x)\right| \leqq \frac{1}{n}|x|^{n}
\end{array}\right.
$$

4) $f^{n}(x)$ and $g^{n}(x)$ have only to be of class $C^{2}$ in $U$ excepting $x=0$.
5) c.f. (2.1).

Then there exists a system $u^{n}(x)=\left(u_{1}^{n}(x), \cdots, u_{m}^{n}(x)\right)$ of functions of class $C^{2}$ satisfying
(9.2) $\quad \sum_{i=1}^{m}\left(\sum_{j=1}^{m} a_{i j} x_{j}+f_{i}^{n}(x)\right) \partial_{i} u_{\nu}=\sum_{\mu=1}^{m} a_{\nu \mu} u_{\mu}+g_{\nu}^{n}(x) \quad(\nu=1, \cdots, m)$,
such that

$$
\left\{\begin{aligned}
\left|u^{n}(x)\right| & \leqq C \cdot K_{2}|x|^{p+1} \\
\left|\partial_{x} u^{n}(x)\right| & \leqq C \cdot K_{3}|x|^{p}
\end{aligned}\right.
$$

in $U$ where $K_{2}$ and $K_{3}$ are constants not depending on $n$. For $u^{n}(x)-u^{n^{\prime}}(x)$ ( $n \leqq n^{\prime}$ ) we have

$$
\sum_{i=1}^{m}\left(\sum_{j=1}^{m} a_{i j} x_{j}+f_{i}^{n}(x)\right) \partial_{i}\left(u_{\nu}^{n}-u_{\nu}^{n^{\prime}}\right)=\sum_{\mu=1}^{m} a_{\nu \mu}\left(u_{\mu}^{n}-u_{\mu}^{n^{\prime}}\right)+h_{\nu}^{n, n^{\prime}}(x),
$$

where

$$
h_{\nu}^{n, n^{\prime}}(x)=\left\{g_{\nu}^{n}(x)-g_{\nu}^{n^{\prime}}(x)\right\}+\partial_{i} u_{\nu}^{n^{\prime}}\left\{f_{\nu}^{n^{\prime}}(x)-f_{\nu}^{n}(x)\right\}
$$

and so

$$
h^{n, n^{\prime}}(0)=0 \quad \text { and } \quad\left|\partial_{x} h^{n, n^{\prime}}(x)\right| \leqq \frac{H}{n}|x|^{p}
$$

with a constant $H>0$ not depending on $n$ and $n^{\prime}$. Therefore we see that, as $n \rightarrow \infty, u^{n}(x)$ tend to the desired solution of (5.1). The proof of the lemma is thus completed.

## § 6. Proof of the Main Theorem.

10. Let us now turn to the system (3.1),

$$
\sum_{i=1}^{m}\left(\sum_{j=1}^{m} a_{i j} x_{j}+f_{i}(x)\right) \partial_{i} u_{\nu}=\sum_{\mu=1}^{m} a_{\nu \mu} u_{\mu}-f_{\nu}(x) \quad(\nu=1, \cdots, m)
$$

where $f_{\nu}(x)=p_{\nu}(x)+q_{\nu}(x)$. First, let us consider

$$
\begin{equation*}
\sum_{i=1}^{m}\left(\sum_{j=1}^{m} a_{i j} x_{j}+p_{i}(x)\right) \partial_{i} u_{\nu}=\sum_{\mu=1}^{m} a_{\nu \mu} u_{\mu}-p_{\nu}(x) \quad(\nu=1, \cdots, m) \tag{10.1}
\end{equation*}
$$

If there exist no relations of the form $\lambda_{i}=\sum n_{j} \lambda_{j}$ where $n_{j}$ are non-negative integers such that $\sum_{j} n_{j}>1$, we can construct infinite series of the form

$$
\underset{\substack{p_{i} \geqq 0 \\ p_{1}+\cdots+p_{m} \geqq 2}}{\sum c_{p_{1}, \cdots, p_{m}}^{\nu} \cdot x_{1}^{p_{1}} \cdots \cdots x_{m}^{p_{m}}}
$$

with real coefficients $c_{p_{1}, \cdots, p_{m}}^{\nu}$, such that $u_{\nu}=\sum c_{p_{1}, \cdots, p_{m}}^{\nu} \cdot x_{1}^{p_{1}} \cdots x_{m}^{p_{m}}$ ( $\nu=1, \cdots, m$ ) satisfy (10.1) formally. To see this, make a change of variables, $x=T \cdot y$ and $u=T \cdot w$, by a (complex) matrix $T$ transforming $A=\left(a_{i j}\right)$ into Jordan's canonical form $T^{-1} \cdot A \cdot T$, and consider about the (complex) system thus obtained.

## Setting

$$
\begin{equation*}
u_{\nu}(x)=\sum_{2 \leqq p_{1}+\cdots+p_{m} \leqq h} c_{p_{1} \cdots p_{m}}^{\nu} \cdot x_{1}^{p_{1}} \cdots x_{m}^{p_{m}} \quad(\nu=1, \cdots, m) \tag{10.2}
\end{equation*}
$$

and

$$
\begin{equation*}
\grave{u}_{\nu}=u_{\nu}-u_{\nu}(x) \quad(\nu=1, \cdots, m), \tag{10.3}
\end{equation*}
$$

we reduce (3.1) to the system

$$
\begin{equation*}
\sum_{i=1}^{m}\left(\sum_{j=1}^{m} a_{i j} x_{j}+f_{i}(x)\right) \partial_{i} \grave{u}_{\nu}=\sum_{\mu=1}^{m} a_{\nu \mu} \grave{u}_{\mu}+\tilde{f}_{\nu}(x) \quad(\nu=1, \cdots, m) \tag{10.4}
\end{equation*}
$$

where

$$
\tilde{f}_{\nu}(x)=\sum_{\mu=1}^{m} a_{\nu \mu} \dot{u}_{\mu}(x)-\sum_{i=1}^{m}\left(\sum_{j=1}^{m} a_{i j} x_{j}+f_{i}(x)\right) \partial_{i} \dot{u}_{\nu}(x)-f_{\nu}(x) .
$$

In order to define $\dot{u}(x)$ as above, we have only to assume that condition (iii) in the theorem (§2) is satisfied.

By (2.1) we have

$$
\begin{equation*}
\tilde{f}_{\nu}(x)=\tilde{p}_{\nu}(x)-\tilde{q}_{\nu}(x) \tag{10.5}
\end{equation*}
$$

where $\quad \tilde{p}_{\nu}(x)=\sum_{\mu=1}^{m} a_{\nu \mu} \hat{u}_{\mu}(x)-p_{\nu}(x)-\sum_{i=1}^{m}\left(\sum_{j=1}^{m} a_{i j} x_{j}+p_{i}(x)\right) \partial_{i} \hat{u}_{\nu}(x)$
and

$$
\tilde{q}_{\nu}(x)=q_{\nu}(x)+\sum_{i=1}^{m} q_{i}(x) \partial_{i} u_{\nu}(x),
$$

and it follows from the definition of $u(x)$ that $\tilde{p}_{v}(x)$, polynomials in $x_{i}(i=1, \cdots, m)$, contain no terms of degree $\leqq h$. Therefore, by (2.2)

$$
\left\{\begin{array}{l}
\tilde{f}_{\nu}(0)=0 \quad(\nu=1, \cdots, m) \\
\left|\partial_{x} \tilde{f}(x)\right| \leqq \tilde{K} \cdot|x|^{h}
\end{array}\right.
$$

where $\tilde{K}$ is a positive constant. Using the lemma, we thus see that there exists a solution $\grave{u}(x)=\left(\grave{u}_{1}(x), \cdots, \grave{u}_{m}(x)\right)$ of (10.4) in some neighborhood of $x=0$, such that

$$
\grave{u}_{i}(0)=0 \quad(i=1, \cdots, m)
$$

and

$$
\left|\partial_{x} \grave{u}(x)\right| \leqq C \cdot \tilde{K}|x|^{h}
$$

with some constant $C>0$. We set

$$
u_{\nu}(x)=u_{\nu}(x)+\grave{u}_{\nu}(x) \quad(\nu=1, \cdots, m)
$$

and obtain those functions $u_{\nu}(x)(\nu=1, \cdots, m)$ whose existence was claimed
in the theorem. The proof of the theorem is thus completed.
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[3] C. L. Siegel: Ueber die Normalform analytischer Differentialgleichungen in der Nähe einer Gleichgewichtslösung. Nachr. Akad. Wiss. Göttingen. Math.Phys. Kl. (1952), 21-30.


[^0]:    2) Cf. (i) $k=0$ means $\mathfrak{R}\left(\lambda_{i}\right)<0$ for all $i$, and $k=m$ means $\mathfrak{R}\left(\lambda_{i}\right)>0$ for all $i$.
[^1]:    3) For the case $k=0$ or $k=m$, we have to set $S_{\varepsilon}(x)=\sum_{i=1}^{m} x_{i}^{2}$.
