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## Numdam

# On the proof of Kuranishi's embedding theorem (*) 

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Abstract. - We prove a local holomorphic embedding theorem for a formally integrable, strictly pseudoconvex CR manifold $M$ with dim $\mathrm{M}=2 n-1 \geqq 7$. This embedding is obtained as the limit of a sequence of approximate embeddings into complex $n$-space, which is constructed and shown to converge by the methods of Nash and Moser. The linearized problem is solved using the explicit integral operators constructed by Henkin. With estimates wich we have previously obtained for these operators, we show that if $\mathbf{M}$ is of class $\mathbf{C}^{m}$, then it admits a $\mathbf{C}^{k}$ embedding provided $21 \leqq k, 6 k+5 n-2 \leqq m$. Our argument is much shorter and simpler than previous arguments, which were based on the Neumann operator and carried out in the $\mathrm{C}^{\infty}$ category.

Résumé. - Nous démontrons un théorème de plongement holomorphe local pour une variété CR M, intégrable et strictement pseudoconvexe, si $\operatorname{dim} \mathbf{M}=2 n-1 \geqq 7$. Ce plongement est obtenu comme limite d'une suite de plongements approximatifs dans l'espace $\mathbb{C}^{n}$. Nous construisons cette suite et démontrons sa convergence par les méthodes de Nash et Moser. Pour le problème linéarisé nous utilisons les opérateurs construits par Henkin et les bornes que nous avons obtenu auparavant. Si M est de classe $\mathrm{C}^{m}$, le plongement est de classe $\mathrm{C}^{k}$ pourvu que $k \geqq 21$ et $m \geqq 6 k+5 n-2$.

Mots clés : Plongement CR, existence locale, méthodes de Nash-Moser.
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## INTRODUCTION

We consider the local holomorphic embedding problem for a formally integrable CR structure of real hypersurface type with positive definite Levi form. It was first posed by Kohn [3] and has been solved in large part by Kuranishi [4] in the $\mathrm{C}^{\infty}$ case. His proof is rather long and technical, involving a delicate study in $\mathrm{L}^{2}$-spaces of the Neumann operator for solving the tangential Cauchy-Riemann equations. This was used in conjunction with a Nash-Moser iteration scheme to produce an embedding. Due to its central importance and the difficult nature of the original proof, the problem merits further study and better understanding.

In broad outline our approch here is similar to Kuranishi's ([4], III), but differs significantly in several important details. We also set up a sequence of approximate holomorphic embeddings and show convergence using the methods of Nash and Moser [5]. But rather than using the Neumann operator to solve the "linearized problem", we use the totally explicit integral operators of Henkin [2] on approximating real hypersurfaces in $\mathbf{C}^{n}$. The necessary estimates, which are given in [10], are much simpler than those in ([4], I, II). Working entirely in $\mathbf{C}^{k}$-spaces, we are able to prove the following version of Kuranishi's theorem.

Theorem. - Let M be a ( $2 n-1$ )-dimensional CR manifold of differentiability class $\mathrm{C}^{m}$. Then M admits, locally near each point, a holomorphic embedding of class $\mathbf{C}^{k}$, provided

$$
\begin{gather*}
n \geqq 4, \quad 2 n-1 \geqq 7,  \tag{0.1}\\
m \geqq 6 k+5 n-2, \quad k \geqq 21 . \tag{0.2}
\end{gather*}
$$

Specifically, if the vector fields defining the structure locally have coefficients of class $\mathbf{C}^{m}$, then they annihilate $n$ independent complex valued functions of class $\mathrm{C}^{k}$.

Kuranishi [4] requires $n \geqq 5,2 n-1 \geqq 9$. Recently, Akahori [1] has given an argument assuming $n \geqq 4$. He also uses $\mathrm{L}^{2}$-methods and the Neumann operator. In view of the examples of Nirenberg [8], only the case $n=3$ is still open. Aside from this the question of regularity remains to be settled. The inequalities ( 0.2 ) can probably be improved, even with the present methods. The derivative loss, $m-k$, comes mainly from small denominators occuring in the estimates of [10].

One aspect of our Nash-Moser argument should be pointed out. Typically in such an argument one controls derivatives up to some order $k$ at the expense of allowing those of order $k+\mu$, say, to become unbounded. For $\mu>k$, some of the constants in section 5, e.g. in (5.3), must involve negative powers of the inner radius of the domain. As the domain shrinks at each stage of iteration, this would cause a problem. We are able to carry out our argument with $\mu<k$, thus avoiding this difficulty and greatly
reducing the derivative loss. However, this requires the lower bound on $k$ in (0.2).
In section 1 we give a simple Taylor series argument for finding an approximate holomorphic embedding. This follows from general principles and might have been omitted. However, it is an integral part of the argument and contributes to the greater part of the derivative loss. We set up the "homotopy formula" for the tangential CR complex in section 2 , using the results of [2] and [10]. Section 3, which shows how to alter an embedding to make it more nearly holomorphic, contains the core of our argument. Sections 4 and 6 are dedicated to the technical details of controlling the embedding. Also in section 5 we state the results from [10] and make some minor modifications of them. Finally, in section 7 we present the convergence argument for the theorem. Aside from those considerations required by smoothing most of the ideas of our proof already occur in greatly simplified form in [9].
This work was brought to conclusion at the Mathematics Institute at E.T.H., Zürich. Many thanks are due to the staff for their kind assistance, and to J. Moser and E. Zehnder for their interest in this work and for several helpful discussions.

## 1. INITIAL NORMALIZATION AND APPROXIMATE HOLOMORPHIC EMBEDDING

Let D be a neighborhood of 0 in $\mathbb{R}^{2 n-1}$, with coordinates $z^{\alpha}=x^{\alpha}+i y^{\alpha}$, $1 \leqq \alpha \leqq n-1$, and $u=\bar{u}$; on which are given $n-1$ complex vector fields $X_{\alpha}$ of class $C^{m}$. We assume that the $X_{\alpha}$ together with their complex conjugates, $\bar{X}_{z} \equiv X_{\bar{\alpha}}$, are pointwise linearly independent and satisfy the integrability condition, $\left[X_{\alpha}, X_{\beta}\right]=\Gamma_{\alpha \beta}^{\gamma} X_{\gamma}$. Where convenient we use the convention that greek indices run from 1 to $n-1$, latin ones from 1 to $n$, and repeated indices are summed. Our vector fields are determined up to a frame change $X_{\alpha} \rightarrow C_{\alpha}^{\beta} X_{\beta}$.
Initially we choose coordinates so that $X_{\alpha} u(0)=X_{\alpha} z^{\beta}(0)=0$, $X_{z} z^{\beta}(0)=\delta_{\alpha}^{\beta}$, and adjust $C$ so that ( $\partial_{\alpha}=\partial / \partial z^{\alpha}$, etc.)

$$
\begin{equation*}
\mathrm{X}_{\alpha}=\partial_{\alpha}+\mathrm{A}_{\alpha}^{\bar{\gamma}} \partial_{\bar{\gamma}}+\mathrm{B}_{\alpha} \partial_{u}, \quad \mathrm{~A}(0)=0, \quad \mathbf{B}(0)=0 \tag{1.1}
\end{equation*}
$$

For such a frame the commutators vanish, $\left[X_{\alpha}, X_{\beta}\right]=0$, giving

$$
\begin{equation*}
\mathbf{X}_{\alpha} \mathbf{A}_{\beta}^{\bar{\gamma}}=\mathbf{X}_{\beta} \mathbf{A}_{\alpha}^{\bar{\gamma}}, \quad \mathbf{X}_{\alpha} \mathbf{B}_{\beta}=\mathbf{X}_{\beta} \mathbf{B}_{\alpha} . \tag{1.2}
\end{equation*}
$$

The Levi form is the hermitian matrix $g$ defined

$$
\begin{gather*}
{\left[\mathrm{X}_{\alpha}, \mathrm{X}_{\bar{\beta}}\right]=-i g_{\alpha \bar{\beta}} \partial_{u}, \quad \bmod \left\{\mathbf{X}_{\gamma}, \mathbf{X}_{\bar{\gamma}}\right\},} \\
g_{\alpha \bar{\beta}}=i\left(\mathbf{X}_{\alpha} \mathbf{B}_{\bar{\beta}}-\mathbf{X}_{\bar{\beta}} \mathbf{B}_{\alpha}\right) . \tag{1.3}
\end{gather*}
$$

We assume that $g$ is positive definite.
We make a coordinate change $F$ on $\mathbb{R}^{2 n-1}$ to simplify $A$ and $B$,

$$
\begin{gathered}
\mathrm{F}: \quad z^{\prime \beta}=\mathrm{F}^{\beta}(z, u)=z^{\beta}+f^{\beta}, \quad f^{\beta}=O(2), \\
u^{\prime}=\mathrm{F}^{n}(z, u)=u+f^{n}, \quad f^{n}=\bar{f}^{n}=O(2), \\
\mathbf{X}_{\alpha}=\mathbf{X}_{\alpha} \mathrm{F}^{\gamma} \partial_{\gamma}^{\prime}+\mathbf{X}_{\alpha} \mathrm{F}^{\prime} \partial_{\bar{\gamma}}^{\prime}+\mathbf{X}_{\alpha} \mathrm{F}^{n} \partial_{u}^{\prime} \equiv \mathrm{C}_{\alpha}^{\beta} \mathbf{X}_{\beta}^{\prime}, \\
\mathbf{X}_{\alpha}^{\prime}=\partial_{\alpha}^{\prime}+\mathbf{A}_{\alpha}^{\prime \gamma} \partial_{\bar{\gamma}}^{\prime}+\mathbf{B}_{\alpha}^{\prime} \partial_{u}^{\prime} .
\end{gathered}
$$

It follows that

$$
X_{\beta} F^{\bar{\gamma}}=\left(X_{\beta} F^{\rho}\right)\left(A_{\rho}^{\prime} \bar{\gamma} \circ F\right), \quad X_{\beta} F^{n}=\left(X_{\beta} F^{\rho}\right)\left(B_{\rho}^{\prime} \circ F\right)
$$

or more explicity,

$$
\begin{align*}
\mathbf{A}_{\beta}^{\bar{\gamma}}+\mathbf{X}_{\beta} f^{\bar{\gamma}} & =\left(\delta_{\beta}^{\rho}+X_{\beta} f^{\rho}\right)\left(\mathbf{A}_{\rho}^{\bar{\gamma}} \circ \mathbf{F}\right),  \tag{1.4}\\
\mathbf{B}_{\beta}+X_{\beta} f^{n} & =\left(\delta_{\beta}^{\rho}+X_{\beta} f^{\rho}\right)\left(\mathbf{B}_{\rho}^{\prime} \circ \mathbf{F}\right) . \tag{1.5}
\end{align*}
$$

We first take $f^{\gamma}=-z^{\bar{\sigma}} \mathbf{A}_{\bar{\sigma}}^{\gamma}(0, u), f^{n}=0$ to get $\mathrm{A}^{\prime}\left(0, u^{\prime}\right)=0$, and then $f^{\gamma}=0$, $f^{n}=-z^{\sigma} \mathbf{B}_{\sigma}(0, u)-z^{\bar{\sigma}} \mathbf{B}_{\bar{\sigma}}(0, u)$ to get $\mathbf{B}^{\prime}(0, u)=0$. However, $F$ would be of class $\mathbf{C}^{m}$ and $\mathrm{A}^{\prime}, \mathrm{B}^{\prime}$ only of class $\mathrm{C}^{m-1}$. Therefore, we replace $\mathrm{A}(0, u)$, $\mathbf{B}(0, u)$, and other functions of $u$ appearing in our transformation formulae by appropriate Taylor polynomials in $u$. This will result in normalization at $(z, u)=(0,0)$. Next we take

$$
f^{\gamma}=-z^{\rho} z^{\bar{\sigma}} \mathbf{X}_{\boldsymbol{\rho}} \mathbf{A}_{\bar{\sigma}}^{\gamma}(0, u)-z^{\bar{\rho}} z^{\bar{\sigma}} \mathbf{X}_{\bar{\rho}} \mathbf{A}_{\bar{\sigma}}^{\gamma}(0, u), \quad f^{n}=0
$$

By (1.2) the coefficient of $z^{\bar{\rho}} z^{\bar{\sigma}}$ is symmetric in $\rho$ and $\sigma$, and hence may be removed by this change. If we differentiate (1.4) in the first instance by $\mathrm{X}_{\alpha}$, in the second by $\mathrm{X}_{\bar{\alpha}}$, and set $z=0$, we see that $\mathrm{A}^{\prime}\left(z^{\prime}, u^{\prime}\right)=O\left(\left|z^{\prime}\right|^{2}\right)$. We next consider $B_{\alpha}$,

$$
\mathbf{B}_{\alpha}=\mathrm{B}_{\alpha \beta}(u) z^{\beta}+\mathrm{B}_{\alpha \bar{\beta}}(u) z^{\bar{\beta}}+\mathrm{B}_{\alpha}^{*}(z, u), \quad \mathrm{B}^{*}=O\left(|z|^{2}\right) .
$$

Again by (1.2), $B_{\alpha \beta}=B_{\beta \alpha}$, so that we may set

$$
f^{\gamma}=0, \quad f^{0}=-\operatorname{Re}\left(\mathrm{B}_{\alpha \beta} z^{\alpha} z^{\beta}\right)-\frac{1}{2}\left(\mathbf{B}_{\alpha \bar{\beta}}+\overline{\mathbf{B}}_{\beta \bar{\alpha}}\right) z^{\alpha} z^{\bar{\beta}} .
$$

After substituting this into (1.5) and taking $X_{\alpha}$-and $X_{\bar{\alpha}}$-derivatives along $z=0$, we find that $\mathrm{B}_{\alpha \beta}^{\prime}\left(u^{\prime}\right)=0$ and that $\mathrm{B}_{\alpha \bar{\beta}}^{\prime}\left(u^{\prime}\right)$ is skew hermitian. Thus we may arrange that $\mathrm{B}_{\alpha}-i g_{\alpha \bar{\beta}}(u) z^{\bar{\beta}}=O\left(|z|^{2}\right)$, where $g$ is hermitian and by (1.3) positive definite. Finally, a change $z^{\prime \alpha}=z^{\beta} W_{\beta}^{\alpha}(u), u^{\prime}=u$ makes $g_{\alpha \bar{\beta}}(u)=\delta_{\alpha \bar{\beta}}$.

This argument shows that with a polynomial change of coordinates we may achieve

$$
\left.\begin{array}{c}
\mathrm{A}_{\alpha}^{\bar{\beta}}(z, u)=O(2) \equiv O\left(|(z, u)|^{2}\right),  \tag{1.6}\\
\mathrm{B}_{\alpha}(z, u)=i \delta_{\alpha \bar{\beta}} z^{\bar{\beta}}+\mathrm{B}_{\alpha}^{*}(z, u), \quad \mathrm{B}^{*}=O(2) .
\end{array}\right\}
$$

With $\mathrm{H}=\delta_{\alpha \bar{\beta}} z^{\alpha} z^{\bar{\beta}} \equiv z \cdot \bar{z}$, we define

$$
\begin{equation*}
\mathrm{Z}(z, u)=\left(z, z^{n}\right), \quad z^{n}=u+i \mathrm{H}(z, u) \tag{1.7}
\end{equation*}
$$

which gives an approximate holomorphic embedding for our structure since $\mathrm{X}_{\bar{\alpha}} z^{\beta}=O(2), \mathrm{X}_{\bar{\alpha}} z^{n}=O$ (2). We shall modify Z so that, in addition $\mathrm{X}_{\bar{\alpha}} \mathrm{Z}=O(m)$. Inductively, we assume $\mathrm{X}_{\bar{\alpha}} \mathrm{Z}=O(s), 2 \leqq s \leqq m$, so that

$$
\begin{aligned}
& \mathrm{X}_{\bar{\alpha}} z^{j}=\sum \mathrm{E}_{\bar{\alpha}}{ }^{j} a \bar{b} z^{a} \bar{z}^{b} u^{i}+O(s+1), \\
& z^{a}=z^{\alpha_{1}} \ldots z^{\alpha_{p}}, \quad \bar{z}^{b}=\bar{z}^{\beta_{1}} \ldots \bar{z}^{\beta^{\beta_{q}}} .
\end{aligned}
$$

Here the coefficients E are symmetric separately in the multi-indices $a$ and $b$, and the sum is over all indices with $p+q+i=s$. Since $X_{\bar{\alpha}} X_{\bar{\beta}} z^{j}=X_{\bar{\beta}} X_{\bar{\alpha}} z^{j}$, they are also symmetric in ( $\alpha b$ ). Therefore, we may make the change-of-embedding

$$
z^{j}=z^{j}-\sum \mathrm{E}_{\bar{a}}^{j} a \bar{b} \bar{i} \bar{z}^{\alpha} z^{a} \bar{z}^{b} u^{i}
$$

to get $\mathrm{X}_{\bar{\alpha}} \mathrm{Z}=O(s+1)$. Since the term subtracted is $O(3),(1.6)$ is preserved.
We shall consider approximate holomorphic embeddings (1.7) of D onto a real hypersurface $\mathbf{M}=\mathbf{Z}(\mathrm{D})$ of the following more general form. With $z^{n}=u+i v, \mathrm{M}$ is given by

$$
\begin{gather*}
r \equiv-v+\mathbf{H}(z, u)=0, \quad \mathbf{H}=b(z)+h(z, u), \\
h=\bar{h} \in \mathrm{C}^{m}(\mathrm{D}), \quad h=O(3),  \tag{1.8}\\
b=\bar{b}^{t}=\mathrm{I}+\bar{b}, \quad|\bar{b}(z)| \leqq \frac{1}{2}|z|^{2} .
\end{gather*}
$$

Following Kuranishi [4], for $\rho>0$ we define

$$
\left.\begin{array}{c}
\mathrm{D}_{\rho} \equiv \mathrm{D}(\rho) \equiv \mathrm{D}(\rho, \mathrm{H})=\left\{(z, u) \in \mathrm{D}: u^{2}+\mathrm{H}-\mathrm{H}^{2}<\rho\right\},  \tag{1.9}\\
\mathrm{M}_{\rho}=\mathrm{Z}\left(\mathrm{D}_{\rho}\right) .
\end{array}\right\}
$$

As we shall see in section $4, D_{\rho}$ and $M_{\rho}$ are smoothly bounded with compact closure if $\rho$ is sufficiently small. We introduce locally defined vector fields on $\mathbb{C}^{n}$ by

$$
\left.\begin{array}{c}
\mathrm{Y}_{\alpha}=\partial_{\alpha}-\left(r_{\alpha} / r_{n}\right) \partial_{n}, \quad 1 \leqq \alpha \leqq n-1 \quad\left(r_{j}=\partial_{j} r\right),  \tag{1.10}\\
\mathrm{T}=\left(1+i h_{u}\right) \partial_{n}+\left(1-i h_{u}\right) \partial_{\bar{n}} .
\end{array}\right\}
$$

They are tangent to M , and restricted to $\mathrm{M}, \mathrm{T}=\partial_{u}, \mathrm{Y}_{\alpha}=\partial_{\alpha}-\left(r_{\alpha} / 2 r_{n}\right) \partial_{u}$. Also, $\mathrm{Y}_{\alpha} r \equiv 0, \mathrm{~T} r \equiv 0, \mathrm{~T} z^{\alpha} \equiv 0, \operatorname{ReT} z^{n} \equiv 1$. Changing our notation from (1.1), we put

$$
\begin{equation*}
\mathbf{X}_{\alpha}=\mathbf{Y}_{\alpha}+\mathrm{A}_{\alpha}{ }^{\bar{\beta}} \mathbf{Y}_{\bar{\beta}}+\mathrm{B}_{\alpha} \mathbf{T} \tag{1.11}
\end{equation*}
$$

along M so that

$$
\begin{equation*}
\mathrm{A}_{\bar{\alpha}}{ }^{\beta}=\mathrm{X}_{\bar{\alpha}} z^{\beta}, \quad \mathrm{B}_{\bar{\alpha}}=\left(1+i h_{u}\right)^{-1}\left\{\mathrm{X}_{\bar{\alpha}} z^{n}+\mathrm{X}_{\bar{\alpha}} z^{\beta}\left(r_{\beta} / r_{n}\right)\right\} . \tag{1.12}
\end{equation*}
$$

In particular Z is a holomorphic embedding if and only if $\mathrm{A}=0, \mathrm{~B}=0$. Vector fields X belonging to the original CR structure which satisfy (1.11) will be said to form a frame adapted to the embedding. They are uniquely determined by the condition

$$
\begin{equation*}
\mathrm{X}_{\alpha} z^{\beta}=\delta_{\alpha}^{\beta} . \tag{1.13}
\end{equation*}
$$

## 2. TANGENTIAL CAUCHY-RIEMANN EQUATIONS

For a function $f$ on $\mathbf{M}$ as in (1.8) we define (summing greek indices from 1 to $n-1$ )

$$
\begin{equation*}
\bar{\partial}_{\mathrm{M}} f=\mathrm{Y}_{\bar{\alpha}} f d z^{\bar{\alpha}} \tag{2.1}
\end{equation*}
$$

For a "tangential" ( $0, q$ )-form $\varphi$ on M ,

$$
\left.\begin{array}{c}
\varphi=\sum_{|\mathrm{C}|=q} \varphi_{\mathrm{C}} d \bar{z}^{\mathbf{C}}, \quad \mathrm{C}=\left(\gamma_{1}, \ldots, \gamma_{q}\right), \quad 1 \leqq \gamma_{j} \leqq n-1  \tag{2.2}\\
\bar{\partial}_{\mathrm{M}} \varphi=\sum \bar{\partial}_{\mathrm{M}} \varphi_{\mathrm{C}} \wedge d \bar{z}^{\mathrm{C}}
\end{array}\right\}
$$

Similarly, we define the operator

$$
\begin{equation*}
\bar{\partial}_{\mathbf{x}} f=\mathrm{X}_{\bar{\alpha}} f d z^{\bar{\alpha}}, \quad \bar{\partial}_{\mathbf{x}} \varphi=\sum \bar{\partial}_{\mathbf{x}} \varphi_{\mathrm{C}} \wedge d \bar{z}^{\mathbf{c}} \tag{2.3}
\end{equation*}
$$

Since both the $Y_{\alpha}$ 's commute and the $X_{\alpha}$ 's commute, we have

$$
\begin{equation*}
\left(\bar{\partial}_{M}\right)^{2}=0, \quad\left(\bar{\partial}_{\mathrm{X}}\right)^{2}=0 \tag{2.4}
\end{equation*}
$$

The ultimate goal is to find $n$ independent functions $f^{j}$ satisfying $\bar{\partial}_{\mathbf{x}} f^{j}=0$.
We extend the fields $Y_{\alpha}$ in (1.10) to a (1,0)-frame on $\mathbb{C}^{n}$ by setting $\mathrm{Y}_{n}=\left(i / r_{n}\right) \partial_{n}$. The dual coframe is $d z^{\alpha}, \theta \equiv-i \partial r, \partial f=\mathrm{Y}_{\alpha} f d z^{\alpha}+\mathrm{Y}_{n} f \theta$. A $(0, q)$-form on $\mathbb{C}^{n}$ has the decomposition $\psi=\psi^{\prime}+\psi^{\prime \prime} \wedge \theta$, where $\psi^{\prime}$ and $\psi^{\prime \prime}$ are tangential $(0, q)$ - and $(0, q-1)$-forms relative to $M$. Since $\bar{\partial} \bar{\theta}=0$, we have

$$
\begin{gathered}
\bar{\partial} \psi=\bar{\partial}_{\mathrm{M}} \psi^{\prime}+\left(\bar{\partial}_{v} \psi^{\prime}+\bar{\partial}_{\mathrm{M}} \psi^{\prime \prime}\right) \wedge \bar{\theta} \\
\bar{\partial}_{v} \psi^{\prime} \equiv(-1)^{q} \sum_{\mathrm{C}} \mathrm{Y}_{n} \psi_{\mathrm{C}}^{\prime} d \overline{z^{\mathrm{C}}}
\end{gathered}
$$

By definition $\bar{\partial}_{b} \varphi$, for a tangential $(0, q)$-form $\varphi$ on M , is given by first extending the coefficients of $\varphi \equiv \varphi^{\prime}$ to be independent of $v$, applying $\bar{\partial}$, and then restricting back to M . Thus,

$$
\begin{equation*}
\bar{\partial}_{b} \varphi=\bar{\partial} \varphi=\bar{\partial}_{M} \varphi+\bar{\partial}_{v} \varphi \wedge \bar{\theta} \tag{2.5}
\end{equation*}
$$

In [2] Henkin has constructed solution operators $\mathbf{P}, \mathrm{Q}$ for $\bar{\partial}_{b}$ on the hypersurfaces $\mathbf{M}_{\rho}$ in (1.9). For a (tangential) ( $0, q$ )-form $\varphi, 1 \leqq q \leqq n-3$, they satisfy

$$
\begin{align*}
& \varphi=\widehat{\partial}_{b} \mathbf{P} \varphi+\mathbf{Q} \widehat{\partial}_{b} \varphi=\left(\bar{\partial}_{b} \mathbf{P} \varphi\right)^{\prime}+\left(\mathbf{Q} \widehat{\partial}_{b} \varphi\right)^{\prime} \\
& =\bar{\partial}_{M} P^{\prime} \varphi+\mathrm{Q}^{\prime}\left(\bar{\partial}_{M} \varphi+\bar{\partial}_{v} \varphi \wedge \bar{\theta}\right), \tag{2.6}
\end{align*}
$$

where the prime denotes the tangential part. It is only at this point that we must restrict to $n \geqq 4$. As noted in [10], $P$ and $Q$ annihilate the ideal generated by $\bar{\theta}$. Thus, setting $P_{M}=P^{\prime}$ and $Q_{M}=Q^{\prime}$, we have the homotopy formula

$$
\begin{equation*}
\varphi=\bar{\partial}_{M} P_{M} \varphi+Q_{M} \bar{\partial}_{M} \varphi \tag{2.7}
\end{equation*}
$$

The operators $\bar{\partial}_{M}, \bar{\partial}_{\mathbf{x}}, \mathbf{P}_{\mathbf{M}}, \mathrm{Q}_{\mathbf{M}}$ will be applied component-wise when $\varphi$ or $f$ is a $\mathbb{C}^{n}$-valued or function on $\mathbf{M}_{\boldsymbol{p}}$.

## 3. ALTERATION OF THE EMBEDDING

We begin with Newton's method. Given an embedding $Z$ as in (1.7), (1.8) with $\bar{\partial}_{\mathrm{x}} \mathrm{Z}$ sufficiently small, we set

$$
\begin{equation*}
Z_{*}=Z(z, u)+F(z, u), \tag{3.1}
\end{equation*}
$$

and try to choose $F$ so that $\bar{\partial}_{x} Z_{*}$ is even smaller. With $\varphi=\bar{\partial}_{x} Z$ in (2.7) we get

$$
\begin{align*}
& \bar{\partial}_{\mathrm{x}} \mathrm{Z}_{*}=\bar{\partial}_{\mathrm{X}} \mathrm{Z}+\bar{\partial}_{\mathrm{M}} \mathrm{~F}+\left(\bar{\partial}_{\mathrm{x}}-\bar{\partial}_{\mathrm{M}}\right) \mathrm{F} \\
&=\bar{\partial}_{\mathrm{M}}\left(\mathrm{P}_{\mathrm{M}} \bar{\partial}_{\mathrm{x}} \mathrm{Z}+\mathrm{F}\right)+\mathrm{Q}_{\mathrm{M}} \bar{\partial}_{\mathrm{M}} \bar{\partial}_{\mathrm{X}} \mathrm{Z}+\left(\bar{\partial}_{\mathrm{X}}-\bar{\partial}_{\mathrm{M}}\right) \mathrm{F} . \tag{3.2}
\end{align*}
$$

At first, we would choose $F=-P_{M} \bar{\partial}_{\mathrm{x}} \mathrm{Z}$. Then, roughly (suitable norms will be introduced later), $\|\mathrm{F}\|=\left\|\mathrm{P}_{\mathrm{M}} \bar{\partial}_{\mathrm{x}} \mathrm{Z}\right\| \leqq\left\|\bar{\partial}_{\mathrm{x}} \mathrm{Z}\right\|$. From (1.11)

$$
\begin{equation*}
\left(\bar{\partial}_{\mathbf{x}}-\bar{\partial}_{\mathrm{M}}\right) \mathrm{F}=\left(\mathrm{A}_{\bar{\alpha}}^{\gamma} \mathbf{Y}_{\gamma}+\mathrm{B}_{\bar{\alpha}} \mathbf{T}\right) \mathrm{F} d z^{\bar{\alpha}}, \tag{3.3}
\end{equation*}
$$

so that (1.12) gives $\left\|\left(\bar{\partial}_{\mathrm{X}}-\bar{\partial}_{\mathrm{M}}\right) \mathrm{F}\right\| \leqq\left\|\bar{\partial}_{\mathrm{x}} \mathrm{Z}\right\|^{2}$. Furthermore,

$$
\begin{equation*}
\bar{\partial}_{\mathrm{M}} \bar{\partial}_{\mathrm{x}} \mathrm{Z}=\left(\bar{\partial}_{\mathrm{M}}-\bar{\partial}_{\mathrm{x}}\right) \bar{\partial}_{\mathrm{x}} \mathrm{Z}=-\left(\mathrm{A}_{\bar{\alpha}}^{\gamma} \mathrm{Y}_{\gamma}+\mathrm{B}_{\bar{\alpha}} \mathrm{T}\right) \mathrm{X}_{\bar{\beta}} \mathrm{Z} d z^{\bar{\alpha}} \wedge d z^{\bar{\beta}} . \tag{3.4}
\end{equation*}
$$

Thus, $\left\|\mathrm{Q}_{\mathrm{M}} \bar{\partial}_{\mathrm{M}} \bar{\partial}_{\mathrm{X}} \mathrm{Z}\right\| \leqq\left\|\bar{\partial}_{\mathrm{M}} \bar{\partial}_{\mathrm{X}} \mathrm{Z}\right\| \leqq\left\|\bar{\partial}_{\mathrm{X}} \mathrm{Z}\right\|^{2}$, and so $\left\|\bar{\partial}_{\mathrm{X}} Z_{*}\right\| \leqq\left\|\bar{\partial}_{\mathrm{X}} \mathrm{Z}\right\|^{2}$.
To develope a rigorous argument, we first observe that the above choice of F will not work in an iteration scheme, since $\mathrm{P}_{\mathrm{M}}$ does not regain fully the derivative lost in applying $\bar{\partial}_{\mathbf{x}}$ to Z . Therefore, we set

$$
\begin{equation*}
\mathrm{F}=-\mathrm{S}_{t} \mathrm{P}_{\mathrm{M}} \bar{\partial}_{\mathrm{x}} \mathrm{Z} \tag{3.5}
\end{equation*}
$$

where $S_{t}$ is a smoothing operator to be described later. From (3.2) we get

$$
\begin{gather*}
\bar{\partial}_{\mathrm{X}} \mathrm{Z}_{*}=\mathrm{I}_{1}+\mathrm{I}_{2}+\mathrm{I}_{3}, \\
\mathrm{I}_{1}=\bar{\partial}_{\mathrm{M}}\left(\mathrm{I}-\mathrm{S}_{\mathrm{t}}\right) \mathrm{P}_{\mathrm{M}} \bar{\partial}_{\mathrm{X}} \mathrm{Z},  \tag{3.6}\\
\mathrm{I}_{2}=\left(\bar{\partial}_{\mathrm{M}}-\bar{\partial}_{\mathrm{X}}\right) \mathrm{S}_{\mathrm{t}} \mathrm{D}_{\mathrm{X}} \mathrm{Z}, \\
\mathrm{I}_{3}=\mathrm{Q}_{\mathrm{M}} \bar{\partial}_{\mathrm{M}} \partial_{\mathrm{X}} \mathrm{Z} .
\end{gather*}
$$

While we may still expect $F$ to be the size of $\bar{\partial}_{\mathrm{x}} \mathrm{Z}$, and $\bar{\partial}_{\mathrm{x}} \mathrm{Z}_{*}$ to be smaller, the pertubation (3.1), (3.5) destroys the form (1.8) of M , which is needed to establish (2.7) and the estimates of [10]. For certain constant vectors $K$ we have

$$
\begin{gather*}
\mathrm{F}(z, u)=\mathrm{K}_{0}+z^{\mathrm{\beta}} \mathrm{~K}_{\beta}+z^{\bar{\beta}} \mathrm{K}_{\bar{\beta}}+u \mathrm{~K}_{n}+\mathrm{F}_{2}(z, u) \\
=\mathrm{K}_{0}+z^{j} \mathrm{~K}_{j}+z^{\bar{\beta}} \mathrm{K}_{\overline{\bar{\beta}}}-i \mathrm{HK}_{n}+\mathrm{F}_{2},  \tag{3.7}\\
\mathrm{~F}_{2}=O(2) .
\end{gather*}
$$

Since $\bar{\delta}_{\mathrm{x}} \mathrm{Z}(0)=0$ initially, we have

$$
\bar{\partial}_{\mathrm{x}} \mathrm{Z}_{\boldsymbol{*}}(0)=\mathrm{K}_{\overline{\mathrm{\beta}}} d z^{\bar{\beta}},
$$

so that the coefficients $K_{\bar{\beta}}$ are dominated by $\left\|\bar{\partial}_{\mathbf{x}} \mathbf{Z}_{\boldsymbol{*}}\right\|$. Therefore, we may set

$$
\left.\begin{array}{c}
\mathrm{Z}^{\prime}=\mathrm{Z}_{*}+\mathrm{E}  \tag{3.8}\\
\mathrm{E}(z, u)=-\mathrm{K}_{0}-z^{j} \mathrm{~K}_{j}-z^{\bar{\beta}} \mathrm{K}_{\bar{\beta}}+\mathrm{E}_{2}(z, u),
\end{array}\right\}
$$

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where $\mathrm{E}_{2}=\boldsymbol{O}(2)$ is a term holomorphic in $z^{j}$, to be chosen shortly. We shall then have $\bar{\partial}_{x} Z^{\prime}(0)=0$ and (roughly)

$$
\left\|\bar{\partial}_{\mathbf{x}} Z^{\prime}\right\| \leqq 2\left\|\bar{\partial}_{\mathrm{x}} \mathrm{Z}_{*}\right\| \leqq 2\left\|\bar{\partial}_{\mathrm{x}} \mathbf{Z}\right\|^{2}
$$

More precisely, we take $E_{2}^{\alpha}=0$ so that

$$
\begin{gather*}
z^{\prime \alpha}=z^{\alpha}+f_{2}^{\alpha}, \quad f_{2}^{\alpha}=\mathrm{F}_{2}^{\alpha}-i \mathrm{HK}_{n}^{\alpha},  \tag{3.9}\\
z^{\prime n}=z^{n}-i \mathrm{HK}_{n}^{n}+\mathrm{F}_{2}^{n}+\mathrm{E}_{2}^{n} . \tag{3.10}
\end{gather*}
$$

With $\mathrm{K}_{n}^{n}=\mathrm{K}_{n}^{\prime n}+i \mathrm{~K}_{n}^{\prime \prime n}$, (3.10) gives

$$
\begin{gather*}
u^{\prime}=u+f_{2}^{n}, \quad f_{2}^{n}=\operatorname{Re}\left(\mathrm{F}_{2}^{n}+\mathrm{E}_{2}^{n}\right)+\mathrm{HK}_{n}^{\prime \prime n},  \tag{3.11}\\
v^{\prime}=\left(1-\mathrm{K}_{n}^{\prime \prime n}\right) \mathrm{H}+\operatorname{Im}\left(\mathrm{F}_{2}^{n}+\mathrm{E}_{2}^{n}\right) . \tag{3.12}
\end{gather*}
$$

We may write

$$
\begin{align*}
\mathrm{F}_{2}^{n}=\mathrm{K}_{\alpha \beta} z^{\alpha} z^{\beta}+\mathrm{K}_{\bar{\alpha} \bar{\beta}} z^{\bar{\alpha}} z^{\bar{\beta}}+\mathrm{K}_{\alpha \bar{\beta}} z^{\alpha} z^{\bar{\beta}} &  \tag{3.13}\\
& +\mathrm{K}_{\alpha n} z^{\alpha} u+\mathrm{K}_{\bar{\alpha} \bar{n}} z^{\bar{\alpha}} u+\mathrm{K}_{n n} u^{2}+\mathrm{F}_{3}^{n} .
\end{align*}
$$

We substitute $u=z^{n}-i H$ in the fourth and sixth terms of (3.13), and $u=\bar{z}^{n}+i \mathrm{H}$ in the fifth. Defining

$$
\begin{equation*}
\mathrm{E}_{2}^{n}=-\mathbf{K}_{\alpha j} z^{\alpha} z^{j}-\mathbf{K}_{n n}\left(z^{n}\right)^{2}+\overline{\mathbf{K}}_{\bar{\alpha} \bar{j}} z^{\alpha} z^{j} \tag{3.14}
\end{equation*}
$$

then gives

$$
\begin{align*}
\mathrm{F}_{2}^{n}+\mathrm{E}_{2}^{n}=\mathrm{K}_{a \bar{\beta}} z^{\alpha} z^{\bar{\beta}}+ & 2 \operatorname{Re}\left(\mathrm{~K}_{\bar{\alpha} \bar{j}} z^{\bar{\alpha}} z^{\bar{j}}\right) \\
& +\mathrm{F}_{3}-i\left(\mathrm{~K}_{\alpha n} z^{\alpha}-\mathrm{K}_{\bar{\alpha} \bar{n}} z^{\bar{\alpha}}\right) \mathrm{H}-\mathrm{K}_{n n}\left(2 i z^{n} \mathrm{H}+\mathrm{H}^{2}\right) \tag{3.15}
\end{align*}
$$

Equation (3.12) may now be written as

$$
\begin{gather*}
v^{\prime}=b_{\alpha \bar{\beta}}^{1} z^{\alpha} z^{\bar{\beta}}+h(z, u)+h^{*}(z, u),  \tag{3.16}\\
b_{\alpha \bar{\beta}}^{1}=b_{\alpha \bar{\beta}}+\hat{b}_{\alpha \bar{\beta} \bar{\beta}}, \quad \hat{b}_{\alpha \bar{\beta}}=-\mathrm{K}_{n}^{\prime n} b_{\alpha \bar{\beta}}+\frac{1}{2 i}\left(\mathrm{~K}_{\alpha \bar{\beta}}-\overline{\mathrm{K}}_{\beta \bar{\alpha}}\right),  \tag{3.17}\\
h^{*}=-\mathrm{K}_{n}^{\prime n} h+\operatorname{Im}\left\{\mathrm{F}_{3}^{n}-i\left(\mathrm{~K}_{\alpha n} z^{\alpha}-\mathrm{K}_{\bar{\alpha} \bar{n}} z^{\bar{\alpha}}\right) \mathrm{H}-\mathrm{K}_{n n}\left(2 i u \mathrm{H}-\mathrm{H}^{2}\right)\right\} . \tag{3.18}
\end{gather*}
$$

We denote by $f: \mathrm{D}_{\rho} \rightarrow \mathbb{R}^{2 n-1}$ the mapping given by (3.9), (3.11) and by $g$ its (local) inverse:

$$
f: \quad \begin{align*}
\mathrm{z}^{\prime \alpha} & =z^{\alpha}+f_{2}^{\alpha}(z, u) ;  \tag{3.19}\\
u^{\prime} & =u+f_{2}^{n}(z, u)
\end{aligned} \quad g: \quad \begin{aligned}
\mathrm{z}^{\alpha} & =z^{\prime \alpha}+g_{2}^{\alpha}\left(z^{\prime}, u^{\prime}\right) \\
u & =u^{\prime}+g_{2}^{n}\left(z^{\prime}, u^{\prime}\right)
\end{align*}
$$

Then $Z_{1}=Z^{\prime} \circ g$, defined on an appropriate neighborhood of 0 in $\mathbb{R}^{2 n-1}$, is the new embedding. By (3.16) and (3.19) the new hypersurface $\mathrm{M}^{1}$ has the equation

$$
\left.\begin{array}{c}
r_{1} \equiv-v+\mathrm{H}_{1}=0, \quad \mathrm{H}_{1}=b^{1}(z)+h_{1}(z, u)  \tag{3.20}\\
h_{1}=\left(h+h^{*}\right) \circ g+b_{\alpha \overline{\mathrm{\beta}}}^{1}\left(z^{\alpha} g_{2}^{\bar{\beta}}+g_{2}^{\alpha} z^{\bar{\beta}}+g_{2}^{\alpha} g_{2}^{\bar{\beta}}\right) .
\end{array}\right\}
$$

We also form the new fields $Y_{\alpha}^{1}, T^{1}$ as in (1.10).
To define the new frame $X_{\alpha}^{1}$ adapted to $Z_{1}$, we first put

$$
\left.\begin{array}{rr}
\mathrm{X}_{\alpha}^{\prime}=\mathrm{C}_{\alpha}^{\beta} \mathrm{X}_{\beta}, & \mathrm{C}=\mathrm{I}+\hat{\mathrm{C}}  \tag{3.21}\\
\delta_{\alpha}^{\beta}=\mathrm{X}_{\alpha}^{\prime} z^{\prime \beta}=\mathrm{C}_{\alpha}^{\gamma}\left(\delta_{\gamma}^{\beta}+\mathrm{X}_{\gamma} f_{2}^{\beta}\right) .
\end{array}\right\}
$$

Then we set

$$
\begin{equation*}
\mathrm{X}_{\alpha}^{1}=d f\left(\mathrm{X}_{\alpha}^{\prime}\right) \tag{3.22}
\end{equation*}
$$

We may regard the $X_{\alpha}^{1}$ as defining a new $C R$ structure equivalent to the original via the map $f$. We have

$$
\begin{gather*}
\mathbf{X}_{\alpha}^{1} z_{1}^{\beta}=X_{\alpha}^{\prime}\left[z_{1}^{\beta} \circ f\right] \circ g=\left(X_{\alpha}^{\prime} z^{\prime \beta}\right) \circ g=\delta_{\alpha}^{\beta}, \\
\mathbf{X}_{\bar{\alpha}}^{1} Z_{1}=\mathbf{X}_{\bar{\alpha}}^{\prime}\left[Z_{1} \circ f\right] \circ g=\left(X_{\bar{\alpha}}^{\prime} Z^{\prime}\right) \circ g,  \tag{3.23}\\
\partial_{\mathbf{x}^{1}} Z_{1}(0)=0 .
\end{gather*}
$$

By repeating this construction we shall define a sequence of embeddings $\mathrm{Z}_{j}$ and diffeomorphisms $f_{j}, j=0,1,2, \ldots$ We must then show that the sequences $\mathrm{Z}_{j}$ and $\mathrm{F}_{j}=f_{j} \circ f_{j-1}{ }^{\circ} \ldots{ }^{\circ} f_{0}$ converge in a neighborhood of 0 in $\mathbb{R}^{2 n-1}$ to $Z_{\infty}$ and $F_{\infty}$. Then $Z_{\infty}$ will give a holomorphic embedding of the structure $d \mathrm{~F}_{\infty}\left(\mathrm{X}_{\alpha}\right)$, or equivalently, $\mathrm{Z}_{\infty} \circ \mathrm{F}_{\infty}$ will embed our original structure.

## 4. GEOMETRIC PROPERTIES OF THE EMBEDDING

The approximating real hypersurface $M$ given by (1.8) has an essentially nonlinear character, which is fully gauged only via the Cartan-ChernMoser theory. We shall not require this theory, but we shall have to control the function $h$, which remains non-zero throughout our argument. Otherwise the domain $\mathrm{D}_{\mathrm{p}}$ would tend to shrink too rapidly during iteration. We should point out that the domains $\left|z^{n}(z, u)\right|<\rho$ are in some ways more natural but need not be convex, a point which causes considerable difficulty. Thus, we have chosen $D_{\rho}$ as in (1.9).

We set

$$
\begin{equation*}
x=(z, u), \quad|x|^{2}=|z|^{2}+u^{2}, \quad|z|^{2}=z \cdot \bar{z} \tag{4.1}
\end{equation*}
$$

and assume that the domain D of $h$ is initially a ball,

$$
\mathbf{D}=\mathbf{B}(\mathbf{R})=\{x:|x|<\mathbf{R}\}, \quad 0<\mathbf{R} \leqq 1 .
$$

Also, we define

$$
\left.\begin{array}{l}
\psi(x)=u^{2}+\mathrm{H}(x)-\mathrm{H}(x)^{2}=\operatorname{Re}\left(z^{n}(x)^{2}-i z^{n}(x)\right)  \tag{4.2}\\
\psi_{0}(x)=u^{2}+b(z)-b(z)^{2}=\operatorname{Re}\left(z_{0}^{n}(x)^{2}-i z_{0}^{n}(x)\right)
\end{array}\right\}
$$

where $z_{0}^{n}=u+i b$. Since $\frac{1}{2}|z|^{2} \leqq b(z) \leqq \frac{3}{2}|z|^{2}$, taking $\mathrm{R}^{2}<1 / 3$ gives

$$
\psi_{0}(x) \leqq u^{2}+b(z) \leqq \frac{3}{2}|x|^{2}
$$

$$
\psi_{0}(x) \geqq u^{2}+b(z)\left(1-\frac{3}{2} \mathrm{R}^{2}\right) \geqq \frac{1}{4}|x|^{2}
$$

If $|h(x)| \leqq c_{h}|x|^{3}$ on D , then $(\mathrm{H}=b+h)$

$$
\begin{aligned}
& \left|\psi-\psi_{0}\right|=|h(1-2 b-h)| \leqq c_{h}|x|^{3}\left(1+3|x|^{2}+c_{h}|x|^{3}\right) \\
& \leqq c_{h}|x|^{3} \leqq c_{h} \mathrm{R}|x|^{2} \leqq c_{h} \mathrm{R} \psi_{0} \leqq \frac{1}{2} \psi_{0}
\end{aligned}
$$

after changing $c_{h}$ and shrinking $R$. Thus on D

$$
\begin{gather*}
\frac{1}{8}|x|^{2} \leqq \psi(x) \leqq \frac{9}{4}|x|^{2} \\
B\left(\frac{2}{3} \sqrt{\rho}\right) \subset D_{\rho} \subset B(\sqrt{8 \rho}) \tag{4.3}
\end{gather*}
$$

We set

$$
\begin{equation*}
\mathrm{N} \equiv \mathbf{N}(h, \rho)=1+\sup \left\{\left|\partial^{\mathrm{I}} h(x)\right|:|\mathrm{I}| \leqq 3, x \in \mathrm{D}_{\rho}\right\} \tag{4.4}
\end{equation*}
$$

In what follows we shall frequently have to assume $\rho$ is so small that

$$
\begin{equation*}
c_{0} N \sqrt{\rho}<1 \tag{4.5}
\end{equation*}
$$

for a sufficiently large absolute constant $c_{0} \geqq 1$, in order that $D_{\rho}$ have certain desired properties. In particular, we have just shown that $D_{p}$ has compact closure and (4.3) holds on $D_{p}$ if $c_{0}$ is sufficiently large. (Absolute constants, denoted by $c, c^{\prime}, c_{j}$, etc., are those independent of particular functions and of the number of derivatives taken in our argument.)

We claim that $D_{p}$ is a smoothly bounded strictly convex neighborhood of 0 if $c_{0}$ in (4.5) is sufficiently large. Since $d \psi(0)=0$, it will suffice to show that the hessian of $\psi$ is positive definite on $D_{\rho}$. But

$$
\begin{aligned}
& \Psi_{i j}(x) \dot{x}^{i} \dot{x}^{j}=2 \dot{u}^{2}+ 2 b(\dot{z})+h_{i j} \dot{x}^{i} \dot{x}^{j} \\
&-2\left(\mathrm{H}_{i} \dot{x}^{i}\right)^{2}-2 \mathrm{HH}_{i j} \dot{x}^{i} \dot{x}^{j} \\
& \geqq|\dot{x}|^{2}-c|\dot{x}|^{2}\left\{\mathrm{~N}|x|+(1+\mathrm{N}|x|)^{2}|x|^{2}\right\} \\
& \geqq|\dot{x}|^{2}-c\left|\dot{x}^{2}\right|^{2}\left\{c_{0}^{-1}+\left(1+c_{0}^{-1}\right)^{2} c_{0}^{-2}\right\} .
\end{aligned}
$$

Hence, if $c_{0}$ is chosen large (relative to $c$ ), this is positive definite.
Next we estimate the distance between $\partial \mathrm{D}_{\mathrm{\rho}}$ and $\partial \mathrm{D}_{\mathrm{\rho}(1-\sigma)}$ for $0<\sigma<1$. For this take $x_{0} \in \partial \mathrm{D}_{\rho(1-\sigma)}, x_{1} \in \partial \mathrm{D}_{\rho}$. By the mean value theorem on $x_{0} x_{1}$,

$$
\begin{gather*}
\rho \sigma=\psi\left(x_{1}\right)-\psi\left(x_{0}\right) \leqq c\left|\partial_{x} \psi\right|_{\rho}\left|x_{1}-x_{0}\right|, \\
\left|\partial_{x} \psi\right| \leqq c\left\{|x|+\mathrm{N}|x|^{2}+|x|^{3}(1+\mathrm{N}|x|)^{2}\right\}  \tag{4.6}\\
\leqq c \sqrt{\rho}\left\{\left(1+c_{0}^{-1}\right)+\rho\left(1+c_{0}^{-1}\right)^{2}\right\}
\end{gather*}
$$

Thus for an absolute constant $c_{1} \geqq 1$,

$$
\begin{equation*}
\operatorname{dist}\left(\partial \mathrm{D}_{\rho}, \partial \mathrm{D}_{\rho(1-\sigma)}\right) \geqq c_{1}^{-1} \rho \sigma \tag{4.7}
\end{equation*}
$$

Next we consider some properties of $D_{p}$ which are necessary for the estimates of [10]. For this let $\mathrm{W}=\left(w, w^{n}\right)$ be a second coordinate vector for $\mathbb{C}^{n}$ and put

$$
\begin{gathered}
\mathrm{S}=i\left(r_{\overline{\mathrm{Z}}} \cdot \partial_{\mathrm{z}}-r_{\mathrm{Z}} \cdot \partial_{\overline{\mathrm{Z}}}\right), \\
p(\mathrm{Z}, \mathrm{~W})=r_{\mathrm{Z}} \cdot(\mathrm{Z}-\mathrm{W}), \quad q(\mathrm{Z}, \mathrm{~W})=-p(\mathrm{~W}, \mathrm{Z}) .
\end{gathered}
$$

In [10] we required

$$
\begin{align*}
& |\mathrm{S} p| \geqq c_{2}^{-1}, \quad|\mathrm{~S} q| \geqq c_{2}^{-1}  \tag{4.8}\\
& |p(\mathrm{Z}, \mathrm{~W})| \geqq c_{2}^{-1}|\mathrm{Z}-\mathrm{W}|^{2} \tag{4.9}
\end{align*}
$$

for a constant $c_{2}>0$, and all $\mathbf{W}, \mathbf{Z}$ in $\mathbf{M}_{\rho}$.
First consider $\mathrm{S} p$ ( $\mathrm{S} q$ is similar and simpler),

$$
\mathrm{S} p=\mathrm{S} r_{\mathrm{Z}} \cdot(\mathrm{Z}-\mathrm{W})+i r_{\mathrm{Z}} \cdot r_{\overline{\mathbf{Z}}}=i\left(r_{\overline{\mathrm{Z}}} \cdot \mathrm{H}_{\mathrm{ZZ}}-r_{\mathrm{Z}} \cdot \mathrm{H}_{\mathrm{Z} \overline{\mathrm{Z}}}\right) \cdot(\mathrm{Z}-\mathrm{W})+\frac{i}{4}+i \mathrm{H}_{\mathrm{Z}} \cdot \mathrm{H}_{\overline{\mathbf{Z}}} .
$$

As

$$
\begin{gathered}
\left|r_{\overline{\mathbf{Z}}} \cdot \mathbf{H}_{\mathrm{ZZ}}\right|+\left|r_{\mathbf{Z}} \cdot \mathbf{H}_{\mathrm{Z} \overline{\mathbf{Z}}}\right| \leqq c\left\{\mathrm{~N}|x|+|x|(1+\mathrm{N}|x|)^{2}\right\} \\
\mathbf{H}_{\mathbf{Z}} \cdot \mathbf{H}_{\overline{\mathbf{Z}}} \leqq c|x|^{2}(1+\mathrm{N}|x|)^{2} \\
|\mathbf{W}-\mathrm{Z}| \leqq c \sqrt{\rho}
\end{gathered}
$$

it's clear that (4.8) will hold for an absolute constant $c_{2}$ if $c_{0}$ is sufficiently large.

For (4.9) we take a second order Taylor expansion of $c$ about $Z$,
$r(\mathrm{~W})-r(\mathrm{Z})=2 \operatorname{Re} p(\mathrm{Z}, \mathrm{W})+2 b(z-w)+2 \int_{0}^{1}(1-t) \partial^{2} h(\mathrm{~W}-\mathrm{Z}, \mathrm{W}-\mathrm{Z}) d t$.
Here the hessian $\partial^{2} h$ is evaluated at $Z_{t}=\mathbf{Z}+t(\mathbf{W}-\mathbf{Z}) \in \mathrm{D}_{\rho} \times \mathbb{R}$. Since $r(\mathrm{Z})=r(\mathrm{~W})=0$,

$$
|z-w|^{2} \leqq 2 b(z-w) \leqq 2|p(\mathrm{Z}, \mathrm{~W})|+c \mathrm{~N}|x| \cdot|\mathbf{Z}-\mathbf{W}|^{2}
$$

In the last term we substitute

$$
z^{n}-w^{n}=\left(r_{z^{n}}\right)^{-1}\left(p(\mathrm{Z}, \mathrm{~W})-\mathrm{H}_{z} .(z-w)\right)
$$

to get

$$
\left.\left.\begin{array}{rl}
|z-w|^{2} \leqq|p|(2+ & \left.c \mathrm{~N}|x| \cdot\left|r_{n}\right|^{2}|p|\right) \\
& +|z-w|^{2} c \mathrm{~N}|x|\left(1+\left|r_{n}\right|^{2}(|x|+\right.
\end{array} \quad \mathrm{N}|x|^{2}\right)\right) .
$$

Here we have used a lower bound on $\left|r_{n}\right|$ and an upper bound on $|p|$ on $\mathrm{D}_{\rho}$. By increasing $c_{0}$, (4.9) is attained with an absolute constant $c_{2}$.

## 5. $\mathbf{C}^{k}$-NORM ESTIMATES

Using multi-index notation for derivatives, we define the usual $C^{k}$ norm of a function $f$ on a domain $D$ in $\mathbb{R}^{n}$ :

$$
\left.\begin{array}{c}
\left|\partial^{j} f\right|_{\mathbf{D}}=\sup \left\{\left|\partial^{J} f(x)\right|:|\mathrm{J}|=j, x \in \mathrm{D}\right\},  \tag{5.1}\\
|f|_{\mathrm{D}, k}=\max \left\{\left|\partial^{j} f\right|_{\mathrm{D}}: 0 \leqq j \leqq k\right\} .
\end{array}\right\}
$$

For vector-valued functions the max is taken over all components. We begin with some elementary properties of these norms relating to the product rule, chain rule, and inverse mapping lemma. With a fixed $k$ in mind we shall consider differently derivatives of order $b$ which are low $b<k$, intermediate $k<b<2 k$, and high $b>2 k$. High derivatives and the usual methods [7] for dealing with them will not enter into our arguments.

From the product rule we clearly have

$$
\begin{equation*}
|f g|_{\mathrm{D}, k} \leqq c(k)|f|_{\mathrm{D}, k}|g|_{\mathrm{D}, k} \tag{5.2}
\end{equation*}
$$

with a constant $c(k)$ depending on $k$. If we take at most $2 k$ derivatives, more than $k$ can fall on only one function; thus

$$
\left.\begin{array}{c}
|f g|_{\mathrm{D}, k+\mu} \leqq c(k+\mu)\left\{|f|_{\mathrm{D}, k}|g|_{\mathbf{D}, k+\mu}+|f|_{\mathbf{D}, k+\mu}|g|_{\mathbf{D}, k}\right\}  \tag{5.3}\\
0<\mu \leqq k
\end{array}\right\}
$$

The estimate is linear in $(k+\mu)$-norms over $k$-norms.
For the chain rule let $G: U \rightarrow V, \varphi: V \rightarrow \mathbb{R}^{p}$, where $U \subset \mathbb{R}^{n}, V \subset \mathbb{R}^{m}$. For $j \geqq 1,|\mathrm{~J}|=j$, we have

$$
\partial^{\mathrm{J}}(\varphi \circ \mathrm{G})=\sum_{s=1}^{j} \varphi_{\alpha_{1} \ldots \alpha_{s}} \sum_{\mathrm{I}_{1} \cup \ldots \cup \mathrm{I}_{s}=\mathrm{J}} \partial^{\mathrm{I}_{1}} \mathrm{G}^{\alpha_{1}} \ldots \partial^{\mathrm{I}_{s}} \mathrm{G}^{\alpha_{s}}
$$

where the subscripts on $\varphi$ denote single derivatives. Hence, for $1 \leqq j \leqq k$

$$
\left|\partial^{j}(\varphi \circ \mathbf{G})\right|_{\mathbf{U}} \leqq c(j)|\varphi|_{\mathbf{v}, j}\left(|\mathbf{G}|_{\mathbf{U}, j}\right)^{j}
$$

and by (5.1)

$$
\begin{equation*}
|\varphi \circ \mathrm{G}|_{\mathrm{U}, k} \leqq c(k)|\varphi|_{\mathrm{v}, k}\left(1+|\mathrm{G}|_{\mathrm{U}, k}\right)^{k} . \tag{5.4}
\end{equation*}
$$

For $|\mathrm{J}|=k+\mu, 0<\mu \leqq k$, we write $\partial^{\mathrm{J}}(\varphi \circ \mathrm{G})=\mathrm{S}_{1}+\mathrm{S}_{2}$, where $\mathrm{S}_{1}$ is the sum of the terms with $s \leqq k$, and $\mathrm{S}_{2}$ the rest. If in $\mathrm{S}_{1}$ one $\left|\mathrm{I}_{v}\right|>k$, then the rest are $<k$. Thus,

$$
\left|\mathrm{S}_{1}\right|_{\mathrm{U}} \leqq c(k+\mu)|\varphi|_{\mathrm{V}, k}\left\{|\mathrm{G}|_{\mathrm{U}, k+\mu}\left(|\mathrm{G}|_{\mathrm{U}, k}\right)^{k-1}+\left(|\mathrm{G}|_{\mathrm{U}, k}\right)^{k}\right\} .
$$

Since $s \geqq k+1$ and each $\left|I_{v}\right| \geqq 1$, each $\left|I_{v}\right| \leqq \mu \leqq k$ in $\mathrm{S}_{2}$, so that

$$
\left|\mathrm{S}_{2}\right|_{\mathrm{U}} \leqq c(k+\mu)|\varphi|_{\mathrm{V}, k+\mu}\left(1+|\mathrm{G}|_{\mathrm{U}, k}\right)^{k+\mu}
$$

Combining these we get

$$
\begin{equation*}
|\varphi \cdot \mathrm{G}|_{\mathrm{U}, \mathrm{k}+\mu} \leqq c(k+\mu)\left(1+|\mathrm{G}|_{\mathrm{U}, \mathrm{k}}\right)^{k+\mu}\left\{|\varphi|_{\mathrm{V}, k+\mu}+|\varphi|_{\mathrm{V}, \boldsymbol{k}}|\mathrm{G}|_{\mathrm{U}, k+\mu}\right\} \tag{5.5}
\end{equation*}
$$

which again is linear in the $(k+\mu)$-norms.
For the inverse mapping lemma let $\mathbf{D}$ and $\mathrm{D}^{\prime}$ be domains containing 0 in $\mathbb{R}^{n}$. Suppose $f$ maps D onto $\mathrm{D}^{\prime}, f(0)=0, f=\mathrm{I}+f_{2}, f_{2}=O\left(|x|^{2}\right)$, and $f$ has inverse $g=\mathrm{I}+g_{2}$ mapping $\mathrm{D}^{\prime}$ into $\mathrm{D}, f \circ g\left(x^{\prime}\right) \equiv x^{\prime}$. With $d$ denoting the Jacobian matrix and $k \geqq 1$, we also assume

$$
\begin{equation*}
\left|d f_{2}\right|_{\mathrm{D}}<\frac{1}{2}, \quad\left|f_{2}\right|_{\mathrm{D}, k}<1 \tag{5.6}
\end{equation*}
$$

Then

$$
\begin{equation*}
d g_{2}=\varphi \circ g, \quad \varphi=\Phi\left(d f_{2}\right) d f_{2}, \quad \Phi(\mathbf{X})=-(\mathbf{I}+\mathbf{X})^{-1} \tag{5.7}
\end{equation*}
$$

By (5.4), (5.2), (5.6) and $V=\left\{|X|<\frac{1}{2}\right\}$,

$$
\begin{align*}
&\left|d g_{2}\right|_{\mathbf{D}^{\prime}, k-1} \leqq c(k-1)|\varphi|_{\mathrm{D}, k-1}\left(1+\left|g_{2}\right|_{\mathbf{D}^{\prime}, k-1}\right)^{k-1} \\
&|\varphi|_{\mathrm{D}, k-1} \leqq c(k-1)\left|d f_{2}\right|_{\mathrm{D}, k-1}|\Phi|_{\mathrm{V}, k-1}(1\left.+\left|d f_{2}\right|_{\mathrm{D}, k-1}\right)^{k-1}  \tag{5.8}\\
& \leqq c(k-1)\left|d f_{2}\right|_{\mathrm{D}, k-1}
\end{align*}
$$

Inductively, we see that

$$
\left.\begin{array}{c}
\left|d g_{2}\right|_{\mathrm{D}^{\prime}, k-1} \leqq c(k-1)\left|d f_{2}\right|_{\mathrm{D}, k-1},  \tag{5.9}\\
\left|g_{2}\right|_{\mathrm{D}^{\prime}, k} \leqq c(k)\left|f_{2}\right|_{\mathrm{D}, k} .
\end{array}\right\}
$$

For the intermediate derivatives we assume $\mu \leqq k-1$. In (5.7) we use (5.5), (5.9), (5.6), (5.8) to get

$$
\left|d g_{2}\right|_{\mathrm{D}^{\prime}, k-1+\mu} \leqq c(k-1+\mu)\left\{|\varphi|_{\mathrm{D}, k-1+\mu}+|g|_{\mathrm{D}^{\prime}, k-1+\mu}\right\} .
$$

From (5.7), (5.5), (5.6)

$$
\begin{aligned}
&|\varphi|_{\mathrm{D}, k-1+\mu} \leqq c(k-1+\mu)\left\{|\Phi \mathbf{X}|_{\mathbf{V}, k-1+\mu}\right. \\
&\left.+|\Phi \mathbf{X}|_{\mathbf{V}, k-1}\left|d f_{2}\right|_{\mathbf{D}, k-1+\mu}\right\} \\
& \leqq c(k+\mu)\left(1+\left|f_{2}\right|_{\mathbf{D}, k+\mu}\right) .
\end{aligned}
$$

It follows by induction that

$$
\begin{equation*}
\left|g_{2}\right|_{\mathrm{D}^{\prime}, k+\mu} \leqq c(k+\mu)\left(1+\left|f_{2}\right|_{\mathrm{D}, k+\mu}\right), \quad 0 \leqq \mu \leqq k-1 . \tag{5.10}
\end{equation*}
$$

Next we recall the estimates for $P$ (and $Q$ ) from [10] and make some adaptations. It's easy to see that inequalities (4.5) of [10] hold with $\delta=c \rho \sigma$ for $M_{\rho}$ as in (1.9). Since $D_{\rho} \times \mathbb{R}$ is convex, the norms over $M_{\hat{\rho}}$ in (4.12) of [10] are bounded by the norms over $D_{p}$.

$$
\begin{equation*}
\left|\partial^{j} f\right|_{\rho}=\left|\partial^{j} f\right|_{D_{\rho}}, \quad|f|_{\rho, k}=|f|_{D_{\rho}, k} \tag{5.11}
\end{equation*}
$$

For $0<\rho, \sigma<1$ formula (4.9) of [10] gives, for a ( $0, q$ )-form $\varphi$ on $D_{\rho}$,

$$
\begin{equation*}
\left|\hat{\partial}^{b} \mathrm{P} \varphi\right|_{\rho(1-\sigma)} \leqq c(b)(\rho \sigma)^{-2(n+b-1)} \sum_{\mathrm{j}=0}^{b}\left|\partial^{j} \varphi\right|_{\rho}\left|w_{0}^{j} w_{1}^{b-j} \mathrm{~A}\right|_{\rho} \tag{5.12}
\end{equation*}
$$

where as explained there

$$
\begin{gather*}
w_{0}^{j} w_{1}^{b-j} \mathrm{~A}=\sum_{\alpha_{1}+\ldots+\alpha_{s}=b-j} \mathrm{~F}_{\alpha_{1}} \ldots \alpha_{s}(v) \partial^{\alpha_{1}} v \ldots \partial^{\alpha_{s}} v  \tag{5.12a}\\
v=\partial^{i} r, \quad 1 \leqq i \leqq 3
\end{gather*}
$$

The F's are certain rational functions whose denominators (4.8) are bounded away from 0 on $D_{p}$. With $b=k-3$, (5.2) and (1.8) readily give

$$
\begin{align*}
& |\mathrm{P} \varphi|_{\rho(1-\sigma), k-3} \leqq \mathrm{~K}|\varphi|_{\rho, k-3}  \tag{5.13}\\
& \mathrm{~K}=c(k) \mathrm{N}^{\gamma(k)}(\rho \sigma)^{-2(n+k)+8} \tag{5.14}
\end{align*}
$$

Here $\gamma(k)$ is a polynomial in $k$, and we have redefined N in (4.4) as

$$
\begin{equation*}
\mathrm{N}=1+|h|_{\mathrm{p}, k}, \tag{5.15}
\end{equation*}
$$

since we shall assume $k \geqq 3$.
For the intermediate derivative estimate we take $b=k-3+\mu, \mu \leqq k-3$, and write $\left|\partial^{b} P \varphi\right|_{\rho(1-\sigma)} \leqq S_{1}+S_{2}$, where $S_{1}$ is the sum on $j$ from 0 to $k-3$. In (5.12a) for $S_{1}, \alpha_{1}+\ldots+\alpha_{s}=k-3+\mu-j$, so that at most one $\alpha_{i}>k-3$. Thus

$$
\mathrm{S}_{1} \leqq c(b) \mathrm{N}^{\gamma(b)} \mathrm{N}(\mu)(\rho \sigma)^{-2(n+k+\mu)+8}|\varphi|_{\rho, k-3}
$$

where

$$
\begin{equation*}
\mathrm{N}(\mu)=1+|h|_{\rho, k+\mu} \tag{5.16}
\end{equation*}
$$

For $S_{2}$ we have $b-j<\mu-1<k-3$ in (5.12), (5.12a) so that

$$
\mathrm{S}_{2} \leqq c(b) \mathrm{N}^{\gamma(b)}(\rho \sigma)^{-2(n+k+\mu)+8}|\varphi|_{\rho, k-3+\mu}
$$

Hence, for $0<\mu \leqq k-3$

$$
\begin{gather*}
|\mathrm{P} \varphi|_{\rho(1-\sigma), k-3+\mu} \leqq \mathrm{K}^{*}\left\{|\varphi|_{\rho, k-3+\mu}+\mathrm{N}(\mu)|\varphi|_{\rho, k-3}\right\}  \tag{5.17}\\
\mathrm{K}^{*}=c(k-3+\mu) \mathrm{N}^{\gamma(k+\mu)}(\rho \sigma)^{-2(n+k+\mu)+8} \tag{5.18}
\end{gather*}
$$

Finally, we recall some standard results from [6] on smoothing operators. For each integer $k>0$ and for $0<t<\infty$ one may construct $\mathrm{S}_{t}$, which is convolution with a smooth function supported in a ball of radius $t$ in $\mathbb{R}^{2 n-1}$. Thus, by (4.7)

$$
\begin{equation*}
\mathrm{S}_{t}: \quad \mathrm{C}^{0}\left(\mathrm{D}_{\rho}\right) \rightarrow \mathrm{C}^{\infty}\left(\mathrm{D}_{\rho(1-\sigma)}\right), \quad 0<t<c_{1}^{-1} \rho \sigma \tag{5.19}
\end{equation*}
$$

For $0 \leqq a \leqq b \leqq k$, the following two properties also hold,

$$
\begin{gather*}
\left|\mathrm{S}_{t} f\right|_{\mathrm{p}(1-\sigma), b} \leqq c(k) t^{a-b}|f|_{\mathrm{p}, a}  \tag{5.20}\\
\left|\left(\mathrm{I}-\mathrm{S}_{t}\right) f\right|_{\mathrm{p}(1-\sigma), a} \leqq c(k) t^{b-a}|f|_{\mathrm{p}, b} \tag{5.21}
\end{gather*}
$$

## 6. ESTIMATES FOR THE NEW EMBEDDING

We now assume that our vector fields $\mathrm{X}_{\alpha}$ are of class $\mathrm{C}^{m}, m \geqq k \geqq 3$, and that we have an approximate holomorphic embedding $\mathbf{Z}$ on $\mathrm{D}_{\mathrm{p}}$. Also, we assume inequality (4.5) with $c_{0}$ so large that all the results of section 4 hold. With $0<t<c_{1}^{-1} \rho \sigma, 0<\sigma<1 / 2$, we make the perturbation (3.1), (3.5). Then for $\tilde{\rho}=\rho(1-2 \sigma),(5.20)$ with $a=k-3, b=k$ and (5.13) give

$$
\begin{aligned}
|\mathrm{F}|_{\tilde{\rho}, k} & \leqq c(k) t^{-3}\left|\mathrm{P}_{\mathrm{M}} \bar{\partial}_{\mathrm{X}} \mathrm{Z}\right|_{\rho(1-\sigma), k-3} \\
& \leqq c(k) t^{-3} \mathrm{~K}\left|\bar{\partial}_{\mathrm{X}} \mathrm{Z}\right|_{\rho, k-3} .
\end{aligned}
$$

Absorbing $c(k)$ into K , we get

$$
\begin{equation*}
|\mathbf{F}|_{\tilde{\mathrm{p}}, k} \leqq \theta, \quad \theta \equiv t^{-3} \mathrm{~K} \delta, \quad \delta \equiv\left|\bar{\partial}_{\mathbf{X}} Z\right|_{\mathrm{\rho}, k-3} \tag{6.1}
\end{equation*}
$$

With $0<\mu \leqq k-3$, all the constants $c(k+\mu)$ in section 5 may be denoted $c(k)$. Then (5.20) $(a=k-3, b=k+\mu)$ and (5.13) give

$$
\begin{equation*}
|\mathbf{F}|_{\tilde{\rho}, k+\mu} \leqq t^{-\mu} \theta, \quad 0<\mu \leqq k-3 \tag{6.2}
\end{equation*}
$$

where another $c(k)$ has been absorbed into $K$. With $f_{2}$ as in (3.19), (3.9), (3.11), we have, using (5.15), (3.14),

$$
\begin{equation*}
\left|f_{2}\right|_{\tilde{\rho}, k} \leqq c\left\{|\mathbf{F}|_{\tilde{\rho}, k}+\theta \mathbf{N}\right\} \leqq c \theta \mathrm{~N}, \tag{6.3}
\end{equation*}
$$

whereas (6.2), (5.16) give

$$
\begin{equation*}
\left|f_{2}\right|_{\tilde{\mathrm{p}}, k+\mu} \leqq c \theta\left(t^{-\mu}+\mathrm{N}(\mu)\right), \quad 0<\mu \leqq k-3 \tag{6.4}
\end{equation*}
$$

By section 1 and (4.3) we may assume initially that for $\rho_{0}>0$

$$
\begin{equation*}
\delta_{0} \leqq c_{\mathrm{X}} \rho_{0}^{(1 / 2)(m-k+3)} \tag{6.5}
\end{equation*}
$$

where the constant $c_{\mathrm{X}}$ depends on the $\mathrm{X}_{\alpha}$. Then with $t_{0}=c_{1}^{-1} \rho_{0} \sigma_{0}, \sigma_{0}$ fixed, (6.1) and (5.14) give

$$
\begin{equation*}
\theta_{0} \leqq c_{\mathrm{X}}(k) \rho_{0}^{(1 / 2)(m-5 k-4 n+13)} \tag{6.6}
\end{equation*}
$$

Thus, if we require [to also get (7.21-22) with $s=1, j=0$ ]

$$
\begin{equation*}
m \geqq 5 k+4 n-10 \tag{6.7}
\end{equation*}
$$

we can make $\delta_{0}$ and $\theta_{0}$ as small as we like by shrinking $\rho_{0}$. We therefore assume $0<t<1,0 \leqq \delta \leqq \theta<1$ in the following.

Next we analyse the properties of the map $f(3.19)$. We want to show that $f$ maps $\mathrm{D}_{\tilde{\rho}}$ onto $\mathrm{D}_{\tilde{\rho}(1-\sigma)}$ and has inverse $g$ mapping $\mathrm{D}_{\tilde{\rho}(1-\sigma)}$ into $\mathrm{D}_{\tilde{\rho}}$. For this we fix $x^{\prime}$ in $\mathrm{D}_{\tilde{\boldsymbol{\rho}}(1-\sigma)}$ and must show that the transformation $\mathrm{W}(x)=x^{\prime}-f_{2}(x)$ has a unique fixed point $x$ in $\mathrm{D}_{\tilde{p}}$. Since $f_{2}=O(2)$, (6.3) and (4.3) give, for $x \in \mathrm{D}_{\tilde{\boldsymbol{p}}}$,

$$
\left|\mathrm{W}(x)-x^{\prime}\right| \leqq c\left|f_{2}\right|_{\tilde{\tilde{p}}, 2}|x|^{2} \leqq c \theta \mathrm{~N} \tilde{\rho} .
$$

Thus $W$ will map $D_{\tilde{\rho}}$ into itself by (4.7), if

$$
\begin{equation*}
c_{3} \theta \mathrm{~N}<\sigma \tag{6.8}
\end{equation*}
$$

for an absolute constant $c_{3}>0$. If $x_{1}, x_{2}$ are in $D_{\tilde{\rho}}$, we apply the mean value theorem on the segment $x_{1} x_{2} \subset \mathrm{D}_{\tilde{\rho}}$ and use (6.3) to get

$$
\begin{aligned}
& \left|\mathrm{W}\left(x_{2}\right)-\mathrm{W}\left(x_{1}\right)\right|=\left|f_{2}\left(x_{1}\right)-f_{2}\left(x_{2}\right)\right| \\
& \leqq c\left|f_{2}\right|_{\tilde{p}, 1}\left|x_{2}-x_{1}\right| \\
& \leqq c \theta \mathrm{~N}\left|x_{2}-x_{1}\right| \text {. }
\end{aligned}
$$

Thus by (6.8), W is a contraction. We denote by $x=g\left(x^{\prime}\right) \in \mathrm{D}_{\tilde{\rho}}$ the unique fixed point, $f\left(g\left(x^{\prime}\right)\right) \equiv x^{\prime}$. It follows that $g=\mathrm{I}+g_{2}$. By (6.3) and (6.8), (5.6) holds, so by (5.9) and (5.10)

$$
\begin{gather*}
\left|g_{2}\right|_{\tilde{\mathcal{p}}(1-\sigma), k} \leqq c(k)\left|f_{2}\right|_{\tilde{\mathfrak{p}}, k},  \tag{6.9}\\
\left|g_{2}\right|_{\tilde{\mathfrak{p}}(1-\sigma), k+\mu} \leqq c(k)\left(1+\left|f_{2}\right|_{\tilde{\rho}, k+\mu}, \quad 0<\mu \leqq k-3 .\right. \tag{6.10}
\end{gather*}
$$

It furthermore $c_{3}$ in (6.8) is sufficiently large (i. e. $\theta$ sufficiently small), then $|f(x)-x|$ will be less than $c_{1}^{-1} \rho^{*} \sigma(4.7)$ on $D_{\rho^{*}(1-\sigma)}$, so that

$$
f\left(\mathrm{D}_{\rho^{*}(1-\sigma)}\right) \subset \mathrm{D}_{\rho^{*}}, \quad \text { for } \quad \rho^{*} \leqq \tilde{\rho}
$$

We have now established the new embedding $Z_{1}=Z^{\prime} \circ g$ on $D_{\tilde{\rho}(1-\sigma)}$; however, its proper domain will be [see (1.9), (4.2)]

$$
\begin{equation*}
\mathrm{D}_{\rho_{1}}^{1}=\left\{x: \psi_{1}(x)<\rho_{1}\right\}, \quad \psi_{1}=\operatorname{Re}\left\{\left(z_{1}^{n}\right)^{2}-i z_{1}^{n}\right\} \tag{6.11}
\end{equation*}
$$

for a suitable $\rho_{1}$. We also set $\psi^{\prime}=\operatorname{Re}\left\{\left(z^{\prime n}\right)^{2}-i z^{\prime n}\right\}$, so that $\psi_{1}=\psi^{\prime} \circ g$. For $x \in \mathrm{D}_{\tilde{\rho}}$, (3.10) and (4.5) give

$$
\begin{gathered}
\left|\psi^{\prime}-\psi\right| \leqq\left|z^{\prime n}-z^{n}\right|\left(1+2\left|z^{n}\right|+\left|z^{\prime n}-z^{n}\right|\right) \\
\left|z^{\prime n}-z^{n}\right| \leqq c\left(\theta|x|^{2}+\theta|\mathbf{H}|\right) \leqq c \theta|x|^{2}(1+\mathrm{N}|x|) \leqq c^{\prime} \theta \psi(x) \\
\left|z^{n}\right| \leqq c|x|\left(1+|x|+\mathrm{N}|x|^{2}\right) \leqq c^{\prime \prime}
\end{gathered}
$$

Thus on $\mathrm{D}_{\tilde{\rho}}$,

$$
\left.\begin{array}{c}
\left|\psi^{\prime}-\psi\right| \leqq \hat{c} \theta \psi  \tag{6.12}\\
\psi(1-\hat{c} \theta) \leqq \psi^{\prime} \leqq \psi(1+\hat{c} \theta),
\end{array}\right\}
$$

so that $D_{p^{*}(1-\hat{c} \theta)}^{\prime} \subset D_{\rho^{*}}$, if $\rho^{*} \leqq \tilde{\rho}$. Since

$$
\mathrm{D}_{\rho_{1}}^{1}=f\left(\mathrm{D}_{\rho_{1}}^{\prime}\right) \subset f\left(\mathrm{D}_{\rho_{1}(1-\hat{c} \theta)}-1\right) \subset \mathrm{D}_{\rho_{1}(1-\hat{c} \theta)}-1_{(1-\sigma)}-1,
$$

$\hat{c} \theta<\sigma$ by (6.8), and we need $D_{\rho_{1}}^{1} \subset D_{\tilde{\rho}(1-\sigma)}$, we require $\rho_{1} \leqq \tilde{\rho}(1-\sigma)^{3}$. Since $\rho(1-2 \sigma)(1-\sigma)^{2} \geqq \rho(1-5 \sigma)$, we take

$$
\begin{equation*}
\rho_{1}=\rho(1-5 \sigma), \quad \sigma<\frac{1}{5} \tag{6.13}
\end{equation*}
$$

On $D_{\rho_{1}}^{1}$ we may compare $Z_{1}$ and $Z$. By (1.7)

$$
\begin{gathered}
\mathrm{Z}_{1}(z, u)-\mathrm{Z}(z, u)=\left(0, i\left(\mathrm{H}_{1}(z, u)-\mathbf{H}(z, u)\right)\right) \\
\mathrm{H}_{1}-\mathrm{H}=\hat{b}(z)+h_{1}(z, u)-h(z, u)
\end{gathered}
$$

By (3.17), since $\left|b_{\alpha \beta}\right|<2$,

$$
\begin{equation*}
|\hat{b}(z)|_{\rho_{1}, k} \leqq c \theta \tag{6.14}
\end{equation*}
$$

From (3.20), (6.9), (6.3), (6.8), (5.2), (5.4)

$$
\begin{gather*}
\left|h_{1}-h\right|_{\rho_{1}, k} \leqq \mathrm{~J}_{1}+\mathrm{J}_{2}+\mathrm{J}_{3}, \\
\mathrm{~J}_{3} \leqq c(k)\left|g_{2}\right|_{\rho_{1}, k}\left(1+\left|g_{2}\right|_{\rho_{1}, k} \leqq c(k) \theta \mathrm{N},\right. \\
\mathrm{J}_{2}=\left|h^{*} \circ g\right|_{\rho_{1}, k} \leqq c(k)\left|h^{*}\right|_{\tilde{\rho}, k},  \tag{6.15}\\
\mathrm{~J}_{1}=|h \circ g-h|_{\rho_{1}, k} .
\end{gather*}
$$

From (3.18)

$$
\begin{equation*}
\left|h^{*}\right|_{\rho_{\rho}, k} \leqq c(k)\left\{\theta|h|_{\rho, k}+\theta+\theta|\mathbf{H}|_{\rho, k}+\theta\left|\mathbf{H}^{2}\right|_{\rho, k}\right\} \leqq c(k) \theta \mathbf{N}^{2} . \tag{6.16}
\end{equation*}
$$

For $\mathrm{J}_{1}$ we set $g_{t}(x)=x+\operatorname{tg}_{2}(x) \in \mathrm{D}_{\tilde{p}}$, then

$$
\begin{gathered}
h(g(x))-h(x)=\int_{0}^{1} \nabla_{x} h\left(g_{t}(x)\right) \cdot g_{2}(x) d t, \\
|h \circ g-h|_{\rho_{1}, k} \leqq c(k)|h|_{\tilde{\rho}, k+1}\left(1+\left.\left|g_{2}\right|_{\rho_{1}, k}\right|^{k}\left|g_{2}\right|_{\rho_{1}, k} \leqq c(k) \theta \mathrm{NN}(1),\right.
\end{gathered}
$$

where $\mathrm{N}(1)$ is given by (5.16) with $\mu=1$. It follows that

$$
\begin{equation*}
\left|h_{1}-h\right|_{\rho_{1}, k} \leqq c(k) \theta \mathrm{NN}(1), \tag{6.17}
\end{equation*}
$$

and with $\mathrm{N}_{1}=1+\left|h_{1}\right|_{\mathrm{p}_{1}, k}$,

$$
\begin{equation*}
\mathrm{N}_{1} \leqq \mathrm{~N}(1+c(k) \theta \mathrm{N}(1)) . \tag{6.18}
\end{equation*}
$$

Next we consider the $(k-3)$-norm of $\bar{\partial}_{\mathbf{x}^{1}} \mathrm{Z}_{1}$, which we denote by $\delta_{1}$. By (3.23), (3.22), (3.21), and (3.8)

$$
\begin{align*}
\left|\mathrm{X}_{\bar{\alpha}}^{1} \mathrm{Z}_{1}\right|_{\rho_{1}, k-3} & =\left|\mathrm{C}_{\bar{\alpha}}^{\bar{\beta}} \mathrm{X}_{\bar{\beta}}\left[\mathrm{Z}_{*}+\mathrm{E}\right] \circ g\right|_{\rho_{1}, k-3} \\
& \leqq\left. c(k)\left(1+\left|g_{2}\right|_{\rho_{1}, k-3}\right)^{k-3}\left|\mathrm{C}_{\tilde{\rho}_{\tilde{p}}, k-3}\right| \bar{\partial}_{\mathrm{x}} \mathrm{Z}_{*}\right|_{\tilde{\tilde{p}}, k-3} \\
& \leqq c(k)|\mathrm{C}|_{\tilde{\rho}, k-3}\left|\bar{\partial}_{\mathrm{X}} Z_{*}\right|_{\tilde{\rho}, k-3} . \tag{6.19}
\end{align*}
$$

By (3.21) $\hat{\mathrm{C}}=\left(\mathbf{I}+\mathbf{X} f_{2}\right)^{-1} \mathbf{X} f_{2}$. Hence, if $\left|\mathbf{X} f_{2}\right|_{\tilde{\mathrm{p}}, 0}<\frac{1}{2}$, then as in (5.7)

$$
|\hat{\mathrm{C}}|_{\tilde{\rho}, k-3} \leqq c(k-3)\left|\mathbf{X} f_{2}\right|_{\tilde{\rho}, k-3}\left(1+\left|\mathbf{X} f_{2}\right|_{\tilde{\rho}, k-3}\right)^{k-3} .
$$

For any function $\varphi$, (1.11) gives
$\left|\mathrm{X}_{\alpha} \varphi\right|_{\tilde{\rho}, k-3} \leqq c(k-3)|\varphi|_{\tilde{\tilde{p}, k-2}}\left\{\left(1+\left|r_{\alpha} / r_{n}\right|_{\tilde{\rho}, k-3}\right)\left(1+|\mathbf{A}|_{\tilde{\rho}, k-3}\right)+|\mathbf{B}|_{\tilde{\rho}, k-3}\right\}$.
Since

$$
\left|r_{\alpha} / r_{n}\right|_{\rho, k-3} \leqq c(k) \mathrm{N}^{k-2}, \quad\left|\left(1+i h_{u}\right)^{-1}\right|_{\rho, k-3} \leqq c(k) \mathrm{N}^{k-3},
$$

we get from (1.12)

$$
\begin{equation*}
|\mathbf{B}|_{\tilde{\rho}, k-3} \leqq c(k) \mathrm{N}^{2 k-5}\left|\bar{\partial}_{\mathrm{X}} \mathrm{Z}\right|_{\tilde{\rho}, k-3} . \tag{6.20}
\end{equation*}
$$

Thus, as $\delta \leqq 1$,

$$
\left|\mathbf{X}_{\alpha} \varphi\right|_{\tilde{\rho}, k-3} \leqq c(k) N^{2 k-5}|\varphi|_{\tilde{\rho}, k-3},
$$

and from (6.3)

$$
\begin{gathered}
\left|\mathrm{X}_{\alpha} f_{2}\right|_{\tilde{\mathrm{p}}, k-3} \leqq c(k) \theta \mathrm{N}^{2 k-4}, \\
\left|\mathrm{X}_{\alpha} f_{2}\right|_{\tilde{\mathrm{p}}, 0} \leqq c \theta \mathrm{~N}^{2} .
\end{gathered}
$$

If we replace (6.8) by

$$
\begin{equation*}
c_{3} \theta \mathrm{~N}^{2 k}<\sigma \tag{6.21}
\end{equation*}
$$

(i.e. decrease $\theta$ ) we have $\mid \mathrm{C}_{\tilde{\mathrm{p}}, k-3} \leqq c(k)$, so that (6.19) gives

$$
\begin{equation*}
\delta_{1} \equiv\left|\bar{\partial}_{\mathbf{x}^{1}} \mathbf{Z}_{1}\right|_{\rho_{1}, k-3} \leqq c(k)\left|\bar{\partial}_{\mathbf{x}} \mathbf{Z}_{*}\right|_{\tilde{\mathrm{p}}, k-3} \tag{6.22}
\end{equation*}
$$

From (3.6), (3.3), (5.20), (5.13) we get
$\left|\mathrm{I}_{2}\right|_{\tilde{\mathrm{p}}, \boldsymbol{k}-3} \leqq c(k)\left\{|\mathrm{A}|\left(1+\left|r_{\alpha} / r_{n}\right|\right)+|\mathrm{B}|\right\}\left|\mathrm{S}_{t} \mathrm{P}_{\mathrm{M}} \bar{\partial}_{\mathrm{X}} \mathrm{Z}\right|_{\tilde{\mathrm{p}}, k-2}$

$$
\leqq c(k)\left|\partial_{\mathrm{X}} \mathrm{Z}\right|_{\rho, k-3} \mathrm{~N}^{2 k-5} t^{-1} \mathrm{~K}\left|\partial_{\mathrm{X}} \mathrm{Z}\right|_{\mathrm{p}, k-3}
$$

By absorbing $c(k) \mathrm{N}^{\mathbf{2 k - 5}}$ into K , we have

$$
\begin{equation*}
\left|\mathrm{I}_{2}\right|_{\tilde{p}, k-3} \leqq t^{-1} \mathrm{~K} \delta^{2} \tag{6.23}
\end{equation*}
$$

From (3.6), (5.21) with $a=k-2, b=k-3+\mu$, we get

$$
\begin{aligned}
&\left|\mathrm{I}_{1}\right|_{\tilde{\mathrm{p}}, k-3} \leqq c(k-3)\left(1+\mathrm{N}^{k-2}\right)\left|\left(\mathrm{I}-\mathrm{S}_{t}\right) \mathrm{P}_{\mathrm{M}} \bar{\partial}_{\mathrm{X}} \mathrm{Z}\right|_{\tilde{\mathrm{p}}, k-2} \\
& \leqq c(k) \mathrm{N}^{k-2} t^{-1+\mu}\left|\dot{\mathrm{P}}_{\mathrm{M}} \bar{\partial}_{\mathrm{X}} \mathrm{Z}\right|_{\rho(1-\sigma), k-3+\mu}
\end{aligned}
$$

With $0<\mu \leqq k-3$, (5.17) gives

$$
\left|\mathrm{I}_{1}\right|_{\tilde{\mathrm{p}}, k-3} \leqq t^{-1+\mu} \mathrm{K}^{*}\left\{\left|\bar{\partial}_{\mathrm{X}} \mathrm{Z}\right|_{\mathrm{p}, k-3+\mu}+\mathrm{N}(\mu) \delta\right\},
$$

where we have absorbed $c(k) \mathrm{N}^{k-2}$ into $\mathrm{K}^{*}$. Since $\delta \leqq 1$,

$$
\begin{gather*}
\left|\mathrm{I}_{1}\right|_{\tilde{\mathrm{p}}, k-3} \leqq t^{-1+\mu} \mathrm{K}^{*} \mathrm{~L}(\mu)  \tag{6.24}\\
\mathrm{L}(\mu)=\left|+\left|\bar{\partial}_{\mathrm{X}} \mathrm{Z}\right|_{\mathrm{p}, k-3+\mu}+|h|_{\mathrm{\rho}, k+\mu}\right. \tag{6.25}
\end{gather*}
$$

Finally, we consider $I_{3}$ in (3.6). By (5.13) for $Q$, (3.4), and arguing as for $I_{2}$, we have

$$
\left\lvert\, \begin{aligned}
\left|\mathrm{I}_{3}\right|_{\tilde{\rho}, k-3} \leqq \mathrm{~K}\left|\bar{\partial}_{\mathrm{M}} \bar{\partial}_{\mathrm{X}} \mathrm{Z}\right|_{\rho(1-\sigma), k-3} \\
\leqq c(k) \mathrm{K}\left|\bar{\partial}_{\mathrm{X}} \mathrm{Z}\right|_{\rho, k-3} \mathrm{~N}^{2 k-5}\left|\bar{\partial}_{\mathrm{X}} \mathrm{Z}\right|_{\rho(1-\sigma), k-2} .
\end{aligned}\right.
$$

For the last factor we use

$$
\begin{gathered}
\bar{\partial}_{\mathrm{X}} \mathrm{Z}=\mathrm{S}_{t} \bar{\partial}_{\mathrm{x}} \mathrm{Z}+\left(\mathrm{I}-\mathrm{S}_{t}\right) \bar{\partial}_{\mathrm{X}} \mathrm{Z} \\
\left|\mathrm{~S}_{t} \bar{\partial}_{\mathrm{X}} \mathrm{Z}\right|_{\rho(1-\sigma), k-2} \leqq c(k) t^{-1}\left|\bar{\partial}_{\mathrm{X}} \mathrm{Z}\right|_{\rho, k-3} \\
\left|\left(\mathrm{I}-\mathrm{S}_{t}\right) \bar{\partial}_{\mathrm{X}} \mathrm{Z}\right|_{\rho(1-\sigma), k-2} \leqq c(k) t^{-1+\mu}\left|\bar{\partial}_{\mathrm{X}} \mathrm{Z}\right|_{\rho, k-3+\mu}
\end{gathered}
$$

Hence, we may write

$$
\begin{equation*}
\left|\mathrm{I}_{3}\right|_{\tilde{p}, k-3} \leqq t^{-1} \mathrm{~K} \delta^{2}+t^{-1+\mu} \mathrm{K} \delta\left|\bar{\partial}_{\mathbf{X}} \mathbf{Z}\right|_{\rho, k-3+\mu} \tag{6.26}
\end{equation*}
$$

Combining the foregoing, we have, as $\mathrm{K} \leqq \mathrm{K}^{*}$,

$$
\begin{equation*}
\delta_{1} \leqq t^{-1} \mathrm{~K} \delta^{2}+t^{-1+\mu} \mathrm{K}^{*} \mathrm{~L}(\mu) . \tag{6.27}
\end{equation*}
$$

The next step is to control the growth of $N(\mu)$ and $L(\mu)$. Arguing as in (6.15) but using the corresponding intermediate derivative estimates, we have

$$
\begin{aligned}
&\left|h_{1}\right|_{\rho_{1}, k+\mu} \leqq \mathrm{J}_{1}(\mu)+\mathrm{J}_{2}(\mu) \\
& \mathrm{J}_{2}(\mu) \leqq c(k)\left(1+\left|g_{2}\right|_{\rho_{1}, k}\right)\left.g_{2}\right|_{\rho_{1}, k+\mu} \\
& \leqq c(k)\left(1+\left|f_{2}\right|_{\tilde{p}, k+\mu}\right) \leqq c(k)\left(1+\theta t^{-\mu}+\theta \mathrm{N}(\mu)\right) \\
& \mathrm{J}_{1}(\mu)=\left|\left(h+h^{*}\right) \circ g\right|_{\rho_{1}, k+\mu} \leqq c(k) \\
& \times\left\{\left|h+h^{*}\right|_{\tilde{\rho}, k+\mu}+\left|h+h^{*}\right|_{\tilde{\rho}, k}\left(1+\left|g_{2}\right|_{\rho_{1}, k+\mu}\right)\right\} .
\end{aligned}
$$

By (6.16), (6.21), (3.18), (6.4), we get

$$
\begin{aligned}
&\left|h+h^{*}\right|_{\tilde{\mathfrak{p}}, k} \leqq c(k) \mathrm{N} \\
&\left|h^{*}\right|_{\tilde{p}, k+\mu} \leqq c(k)\left\{\theta \mathrm{N}+|\mathbf{F}|_{\tilde{\mathfrak{p}}, k+\mu}+\theta \mathrm{NN}(\mu)\right\} \\
& \leqq c(k) \theta\left\{t^{-\mu}+\mathrm{NN}(\mu)\right\} .
\end{aligned}
$$

Thus, by (6.21), (6.10), and (6.4)

$$
\begin{equation*}
\mathrm{N}_{1}(\mu) \leqq c(k)\left(\mathrm{N}(\mu)+\theta t^{-\mu}\right) \tag{6.28}
\end{equation*}
$$

For the term $\bar{\partial}_{\mathbf{x}^{1}} Z_{1}$ we first note that by (3.21-23)

$$
\begin{aligned}
\mathrm{X}_{\bar{\alpha}}^{1} \mathrm{Z}_{1} & =\varphi \circ g, \quad \varphi=\Phi(\mathrm{G}), \\
\mathrm{G} & =\left(\bar{\partial}_{\mathrm{x}} \mathrm{Z}, \partial^{1} \mathrm{~F}, \partial^{1} h\right),
\end{aligned}
$$

and apply (5.5) twice. This gives

$$
\begin{gathered}
\left|\partial_{\mathrm{x}^{1}} \mathrm{Z}_{1}\right|_{\mathfrak{p} 1, k-3+\mu} \leqq c(k)\left\{|\varphi|_{\tilde{p}, k-3+\mu}+|\varphi|_{\tilde{\mathrm{p}}, k-3}\left(1+\left|g_{2}\right|_{\rho_{1}, k-3+\mu}\right)\right\}, \\
|\varphi|_{\tilde{p}, k-3} \leqq c(k)\left(1+|\mathbf{G}|_{\tilde{p}, k-3}\right)^{k-3} \leqq c(k) \mathrm{N}^{k-3}, \\
|\varphi|_{\tilde{\mathrm{p}}, k-3+\mu} \leqq c(k) \mathrm{N}^{k-3+\mu}\left(1+|\mathrm{G}|_{\tilde{\mathrm{p}}, k-3+\mu}\right), \\
|\mathrm{G}|_{\tilde{\mathrm{p}}, k-3+\mu} \leqq\left|\bar{\partial}_{\mathrm{X}} \mathrm{Z}\right|_{\tilde{\mathrm{p}}, k-3+\mu}+|h|_{\tilde{\mathrm{p}}, k-2+\mu}+t^{2-\mu} \theta .
\end{gathered}
$$

By (6.10) and (6.4)

$$
1+\left|g_{2}\right|_{p_{1}, k-3+\mu} \leqq c(k)\left(1+\theta t^{3-\mu}+\theta \mathrm{N}(\mu)\right)
$$

Combining gives

$$
\left|\bar{\partial}_{\mathbf{x}^{1}} Z_{1}\right|_{\rho_{1}, k-3+\mu} \leqq c(k) \mathrm{N}^{k-3+\mu}\left\{\mathrm{L}(\mu)+\theta t^{2-\mu}\right\}
$$

Thus, using (6.28) we get

$$
\begin{equation*}
\mathrm{L}_{1}(\mu) \leqq c(k) \mathrm{N}^{k-3+\mu}\left\{\mathrm{L}(\mu)+\theta t^{-\mu}\right\} \tag{6.29}
\end{equation*}
$$

## 7. THE SEQUENCE OF EMBEDDINGS

We must show that the foregoing process can be repeated an infinity of times and leads to a sequence of embeddings which converges to a
holomorphic embedding. In view of (6.13) we first define for $j \geqq 0$

$$
\begin{equation*}
\rho_{j+1}=\rho_{j}\left(1-5 \sigma_{j}\right), \quad \sigma_{j}=5^{-j-2} \tag{7.1}
\end{equation*}
$$

where $\rho_{0}>0$ is yet to be determined. By taking the logarithm of the infinite product, one sees that $\rho_{*}=\lim \rho_{j}>0$,

$$
\begin{equation*}
\rho_{*}=\rho_{0} \prod_{j=0}^{\infty}\left(1-5 \sigma_{j}\right) . \tag{7.2}
\end{equation*}
$$

We first assume that we have constructed the $\mathrm{Z}_{j}=\left(z, u+i \mathrm{H}_{j}\right)$ on $\mathrm{D}\left(\rho_{j}\right)$ with

$$
\begin{gather*}
c_{0} \mathrm{~N}_{j} \sqrt{\rho_{j}}<1,  \tag{7.3}\\
0<t_{j} \leqq c_{1}^{-1} \rho_{j} \sigma_{j},  \tag{7.4}\\
c(k) \theta_{j} \mathrm{~N}_{j}^{2 k}<\sigma_{j}, \quad \theta_{j}=t_{j}^{-3} \mathrm{~K}_{j} \delta_{j}, \tag{7.5}
\end{gather*}
$$

where $\mathrm{N}_{j}=1+\left|h_{j}\right|_{\rho_{j}, k}$, etc, and determine what further conditions are need for convergence.

We define

$$
\begin{gather*}
\mathrm{F}_{l}=f_{l} \circ f_{l-1} \circ \ldots \circ f_{0}=f_{l} \circ \mathrm{~F}_{l-1}: \mathrm{U}_{l} \rightarrow \mathbb{R}^{2 n-1}  \tag{7.6}\\
\mathrm{G}_{l}=g_{0} \circ g_{1} \circ \ldots \circ g_{l}=\mathrm{G}_{l-1} \circ g_{l}: \mathrm{D}\left(\rho_{l+1}\right) \hookrightarrow \mathrm{D}\left(\rho_{0}\right) . \tag{7.7}
\end{gather*}
$$

Here, $\mathrm{U}_{l}=\mathrm{G}_{l}\left(\mathrm{D}\left(\rho_{l+1}\right)\right) \subset \mathrm{D}_{\mathrm{p}_{\mathrm{o}}}$ contains 0 but is otherwise difficult to specify. However, since the $D\left(\rho_{j}\right) \supset B\left(2 \sqrt{\rho_{j}} / 3\right)(4.3)$, and they decrease, all the $\mathrm{G}_{l}$ are defined on

$$
\mathbf{D}^{*} \equiv \bigcap\left\{\mathbf{D}\left(\rho_{j}\right): 0 \leqq j<\infty\right\}
$$

which contains $B\left(2 \sqrt{\rho^{*}} / 3\right)$. By (6.14) and (6.17)

$$
\begin{gather*}
\left|b_{j+1}-b_{j}\right|_{\rho_{j+1}, k} \leqq c \theta_{j}  \tag{7.8}\\
\left|h_{j+1}-h_{j}\right|_{\mathrm{p}_{j+1}, k} \leqq c(k) \theta_{j} \mathrm{~N}_{j} \mathrm{~N}_{j}(1) \tag{7.9}
\end{gather*}
$$

Thus, if we can show

$$
\begin{equation*}
\sum \theta_{j} N_{j} N_{j}(1)<\infty \tag{7.10}
\end{equation*}
$$

then the $\mathrm{Z}_{j}$ will converge in $\mathrm{C}^{k}$ norm to $\mathrm{Z}_{*}=\left(z, u+i \mathrm{H}_{*}(z, u)\right)$, an embedding defined on a neighborhood of 0 in $\mathbb{R}^{2 n-1}$. The corresponding vector fields (1.10), $\mathrm{Y}_{\alpha}^{j}$, will converge in $\mathrm{C}^{k-1}$ to $\mathrm{Y}_{\alpha}^{*}$. Since, (1.11),

$$
\begin{equation*}
\mathbf{X}_{\alpha}^{j}=\mathrm{Y}_{\alpha}^{j}+\mathrm{A}_{\alpha}^{j \bar{\gamma}} \mathrm{Y}_{\bar{\gamma}}+\mathrm{B}_{\alpha}^{j} \partial_{u} \tag{7.11}
\end{equation*}
$$

and $\left(\mathrm{A}^{j}, \mathrm{~B}^{j}\right) \rightarrow 0$ in $\mathrm{C}^{k-3}$ by (1.12), (6.20), and (7.5), $\mathrm{X}_{\alpha}^{j}$ will converge to $\mathrm{X}_{\alpha}^{*} \equiv \mathrm{Y}_{\alpha}^{*}$ in $\mathrm{C}^{k-3}$. Thus, we shall have produced an embedded real hypersurface of class $C^{k}$ which is, in fact, equivalent to our original structure. However, we shall not yet have produced a $C^{k}$ solution to $\bar{\partial}_{\mathrm{X}} \mathrm{Z}=0$, for the original $\mathrm{X}_{\alpha}=\mathrm{X}_{\alpha}^{0}$. We must still analyse the maps $\mathrm{F}_{l}, \mathrm{G}_{\mathrm{l}}$, as in [9].

By (6.3), (6.9)

$$
\begin{equation*}
\left|f_{2, j}\right|_{\tilde{p}_{j}, k} \leqq c \theta_{j} \mathbf{N}_{j,} \quad\left|g_{2, j}\right|_{\tilde{p}_{j}\left(1-\sigma_{j}\right), k} \leqq c(k) \theta_{j} \mathbf{N}_{j} . \tag{7.12}
\end{equation*}
$$

Letting $d$ denote Jacobian matrix, and $\|$.$\| the operator norm on matrices,$ we have

$$
\begin{gather*}
d \mathrm{G}_{l}=\left(d \mathrm{G}_{l-1}{ }^{\circ} g_{l}\right) d g_{l}=d g_{0} \ldots d g_{l} \\
\left\|d \mathrm{G}_{l}\right\|_{\rho_{l+1}}
\end{gather*} \begin{aligned}
& j=0  \tag{7.13}\\
& \\
& \leqq \prod_{j=0}^{\infty}\left(1+\left\|d g_{2, j}\right\|_{\rho_{j+1}}\right)
\end{aligned}
$$

By (7.10) this is finite giving $\left|G_{l}\right|_{\rho_{l+1}, 1} \leqq b^{\prime \prime}$ for all $l$. From this follows the uniform convergence of the $G_{l}$ to a continuous $G_{*}: D^{*} \rightarrow D\left(\rho_{0}\right)$. For, arguing as for $\mathrm{J}_{1}$ in (6.15),

$$
\left|\mathrm{G}_{l}-\mathrm{G}_{l-1}\right|_{\rho_{l+1}, 0} \leqq\left\|d \mathrm{G}_{l-1} \circ g_{l}\right\|_{\rho_{l+1}}\left\|g_{2, l}\right\|_{\rho_{l+1}, 0} \leqq b^{\prime \prime} c(0) \theta_{j} \mathrm{~N}_{j}
$$

and we have convergence by (7.10). Similarly, the $F_{l}$ are defined on the set $\mathrm{U}_{*}=\mathrm{G}_{*}\left(\mathrm{D}^{*}\right)$, and we have a common bound

$$
c^{-1}\left|\mathrm{~F}_{l}\right|_{\mathrm{U}_{l}, 1} \leqq\left\|d \mathrm{~F}_{l}\right\|_{\mathrm{U}_{l}} \leqq b^{\prime} .
$$

Now,

$$
\left.\begin{array}{c}
d \mathrm{~F}_{j}-d \mathrm{~F}_{l-1}=\left(d f_{2, l} \circ \mathrm{~F}_{l-1}\right) d \mathrm{~F}_{l-1}  \tag{7.14}\\
\left\|d \mathrm{~F}_{l}-d \mathrm{~F}_{l-1}\right\|_{\mathrm{U}_{l}} \leqq\left\|d f_{2, l}\right\|_{\tilde{\mathfrak{p}} l}\left\|d \mathrm{~F}_{l-1}\right\|_{\mathrm{U}_{l-1}} \leqq c \theta_{l} \mathrm{~N}_{l} b^{\prime} .
\end{array}\right\}
$$

Thus, by (7.10), (7.14), and (7.15) with $s=3$, the $\mathrm{DG}_{l}=\left(d \mathrm{~F}_{l}\right)^{-1} \circ \mathrm{G}_{l}$ converge uniformly on $\mathrm{D}^{*}$. Since the $d \mathrm{G}_{l}$ now have the positive lower bound $1 / b^{\prime}$ in norm, it follows that the $\mathrm{G}_{l}$ converge in $\mathrm{C}^{1}$-norm to $\mathrm{G}_{*}$, which is a diffeomorphism (after shrinking $D_{*}$ ). Thus, $U_{*}$ is an open set and the $\mathrm{F}_{j}$ converge in $\mathrm{C}^{1}$ to a diffeomorphism $\mathrm{F}_{*} . \mathrm{F}_{*}$ is a CR equivalence between the structures $\mathrm{X}_{\alpha}^{0}, \mathrm{X}_{\alpha}^{*}$.
To show that the $\mathrm{F}_{j}$ converge in $\mathrm{C}^{k}$ to $\mathrm{F}_{*}$, we fix $s, 2 \leqq s \leqq k$ and assume that we have a bound

$$
\begin{equation*}
\left|\mathrm{F}_{l}\right|_{\mathrm{U}_{l}, s-1} \leqq b_{s-1}^{\prime} \tag{7.15}
\end{equation*}
$$

where $b_{s-1}^{\prime}$ is a constant depending on our given structure. Then,

$$
\begin{gather*}
d \mathrm{~F}_{l}=\left(d f_{l} \circ \mathrm{~F}_{l-1}\right)\left(d \mathrm{~F}_{l-1}\right), \\
\left|d \mathrm{~F}_{l}\right|_{\mathrm{U}_{l}, s-1} \leqq c(s)\left|d f_{l}\right|_{\tilde{p} l}, s-1  \tag{7.16}\\
\left(1+\left|\mathrm{F}_{l-1}\right|_{\mathrm{U}_{l}, s-1}\right)^{s-1}\left|d \mathrm{~F}_{l-1}\right|_{\mathrm{U}_{l-1}, s-1} \\
\leqq c(s)\left(1+b_{s-1}^{\prime}\right)^{s-1}\left|d \mathrm{~F}_{l-1}\right|_{\mathrm{U}_{l-1}, s-1} .
\end{gather*}
$$

Thus,

$$
\begin{equation*}
\left|\mathrm{F}_{l}\right|_{\mathrm{U}_{l}, s} \leqq b_{s}\left|\mathrm{~F}_{l-1}\right|_{\mathrm{U}_{l-1}, s} \leqq \ldots \leqq b_{s}^{l} a_{s} \tag{7.17}
\end{equation*}
$$

where $b_{s}$ and $a_{s}$ are constants depending on $s \leqq k$ and the $\mathrm{X}_{\alpha}^{0}$. By (7.14)

$$
\begin{align*}
&\left|d \mathrm{~F}_{l}-d \mathrm{~F}_{l-1}\right|_{\mathrm{U}_{l}, s-1} \leqq c(s)\left|d f_{2, l}\right|_{\tilde{\tilde{l}_{l}}, s-1}\left(1+b_{s-1}^{\prime}\right)^{s-1} b_{s}^{l} a_{s} \\
& \leqq c(s) \theta_{l} \mathrm{~N}_{l}\left(b_{s}^{2}\right)^{l} a_{s} \tag{7.18}
\end{align*}
$$

By induction there is a constant $b=b\left(k, \mathrm{X}_{\alpha}^{0}\right)$, such that, if

$$
\begin{equation*}
\sum_{j=0}^{\infty} \theta_{j} N_{j} N_{j}(1) b^{j}<\infty \tag{7.19}
\end{equation*}
$$

then (7.15) holds for $s=k+1$, and $\mathrm{F}_{l} \rightarrow \mathrm{~F}_{*}$ in $\mathrm{C}^{k}$.
If we can establish an a priori upper bound $\mathrm{N}^{*} \geqq \mathrm{~N}_{j}$, and then shrink $\rho_{0}$ so that $c_{0} N^{*} \sqrt{\rho_{0}}<1$, (7.3) will hold for all $j$. From (6.18)

$$
\begin{align*}
\mathbf{N}_{j+1} & \leqq \mathrm{~N}_{j}\left(1+c(k) \theta_{j} \mathrm{~N}_{j}(1)\right) \\
& \leqq \mathrm{N}_{\mathrm{O}} \prod_{i=0}^{j}\left(1+c(k) \theta_{i} \mathbf{N}_{i}(1)\right), \tag{7.20}
\end{align*}
$$

so that (7.19) will guarantee such a bound $\mathrm{N}^{*}$. For a suitable $s>0$ to be determined, we define

$$
\begin{equation*}
\zeta_{j}=t_{j}^{-s} \theta_{j} \mathrm{~N}_{j} \mathrm{~N}_{j}(1) b^{j}=t_{j}^{-s-3} \delta_{j} \mathrm{~K}_{j} \mathrm{~N}_{j}(1) b^{j}, \tag{7.21}
\end{equation*}
$$

where for convenience we have absorbed another factor of $\mathrm{N}_{j}$ into $\mathrm{K}_{j}$. If we can show $\zeta_{j} \leqq 1$, then by (7.4), (7.1)

$$
\theta_{j} \mathbf{N}_{j} \mathbf{N}_{j}(1) b^{j} \leqq c_{1}^{-s}\left(\rho_{j} \sigma_{j}\right)^{s} \leqq \sigma_{j}^{s} .
$$

We shall then have an upper bound for all partial sums in (7.19).
Now we assume that we have constructed $Z_{i}, i \leqq j$, satisfying (7.4), (7.5) and for a constant $M$

$$
\begin{equation*}
\zeta_{j} \leqq \mathrm{M}, \quad \mathrm{M}<\frac{1}{2} . \tag{7.22}
\end{equation*}
$$

Then with any $t_{j+1} \leqq c_{1}^{-1} \rho_{j+1} \sigma_{j+1}$, we may construct $\mathrm{Z}_{j+1}$ on $\mathrm{D}_{\rho_{j+1}}$. By (6.18), (6.27), (6.28) and (6.29) we have (7.20) and

$$
\begin{gather*}
\delta_{j+1} \leqq t_{j}^{-1} \mathrm{~K}_{j} \delta_{j}^{2}+t_{j}^{-1+\mu} \mathrm{K}_{j}^{*} \mathrm{~L}_{j}(\mu),  \tag{7.23}\\
\mathrm{N}_{j+1}(\mu) \leqq c(k)\left(\mathrm{N}_{j}(\mu)+\theta_{j} t_{j}^{-\mu}\right),  \tag{7.24}\\
\mathrm{L}_{j+1}(\mu) \leqq c(k) \mathrm{N}_{j}^{k-3+\mu}\left(\mathrm{L}_{j}(\mu)+\theta_{j} t_{j}^{-\mu}\right) . \tag{7.25}
\end{gather*}
$$

We investigate the growth of $\mathbf{N}_{j}, \mathbf{K}_{j}, \mathbf{K}_{j}^{*}$, and $\mathbf{N}_{j}(1)$ with $j$. By (7.20), (7.21), (7.22), (5.14), (5.18) we have, successively increasing $\hat{c}(k) \geqq 1$
each time,

$$
\begin{gather*}
\frac{\mathrm{N}_{j+1}}{\mathrm{~N}_{j}} \leqq c(k)\left(1+t_{j}^{s} \zeta_{j}\right) \leqq \hat{c}(k), \\
\frac{\mathrm{K}_{j+1}}{\mathrm{~K}_{j}} \leqq c(k) \hat{c}(k)^{\gamma(k)}\left(\frac{5}{1-5 \sigma_{j}}\right)^{2(n+k)-8} \leqq \hat{c}(k),  \tag{7.26}\\
\frac{\mathrm{K}_{j+1}^{*}}{\mathrm{~K}_{j}^{*}} \leqq c(k) \hat{c}(k)^{\gamma(k+\mu)}\left(\frac{5}{1-5 \sigma_{j}}\right)^{2(n+k+\mu)-8} \leqq \hat{c}(k) .
\end{gather*}
$$

From (7.24) and (7.21) and $\mu=1$

$$
\begin{equation*}
\frac{\mathrm{N}_{j+1}(1)}{\mathrm{N}_{j}(1)} \leqq c(k)\left(1+\zeta_{j} t_{j}^{s-1}\right) \leqq \hat{c}(k), \tag{7.27}
\end{equation*}
$$

if we choose $s \geqq 1$.
Next we consider $\zeta_{j+1}$. By (7.21), (7.23), (7.26), (7.27), we have

$$
\begin{gather*}
\zeta_{j+1} \leqq t_{j+1}^{-s-3} \mathrm{~K}_{j+1} \mathrm{~N}_{j+1}(1) b^{j+1}\left\{t_{j}^{-1} \mathrm{~K}_{j} \delta_{j}^{2}+t_{j}^{-1+\mu} \mathrm{K}_{j}^{*} \mathrm{~L}_{j}(\mu)\right\} \\
\leqq t_{j+1}^{-s-3} c_{b}\left\{t_{j}^{2 s+5} \zeta_{j}^{2}+t_{j}^{-1+\mu} b^{j} \mathrm{~K}_{j} \mathrm{~K}_{j}^{*} \mathrm{~N}_{j}(1) \mathrm{L}_{j}(\mu)\right\} \\
c_{b} \equiv \hat{c}(k)^{2} b . \tag{7.28}
\end{gather*}
$$

If we choose $t_{j+1}$ so that the coefficient of $\zeta_{j}^{2}$ is 1 ;

$$
\begin{equation*}
t_{j+1}=c_{b}^{2-x} t_{j}^{x}, \quad \kappa=\frac{2 s+5}{s+3}, \quad s=\frac{3 \kappa-5}{2-\kappa} \tag{7.29}
\end{equation*}
$$

then

$$
\begin{equation*}
\zeta_{j+1} \leqq \zeta_{j}^{2}+t_{j}^{\mu-2 s-6} b^{j} \mathbf{K}_{j} \mathbf{K}_{j}^{*} \mathbf{N}_{j}(1) \mathbf{L}_{j}(\mu) \tag{7.30}
\end{equation*}
$$

Therefore, we define the $t_{j}$ inductively by $t_{0} \leqq c_{1}^{-1} \rho_{0} \sigma_{0}$ and (7.29). If (7.4) holds for $j$ then
$t_{j+1} \leqq c_{b}^{2-x}\left(\frac{\rho_{j} \sigma_{j}}{\rho_{j+1} \sigma_{j+1}}\right)^{x}\left(\rho_{j+1} \sigma_{j+1}\right)^{x-1} \rho_{j+1} \sigma_{j+1} \leqq c\left(\rho_{0} \sigma_{0}\right)^{x-1} \rho_{j+1} \sigma_{j+1}$.
Once we have determined $\kappa>1$, we shrink $\rho_{0}$ so that $c\left(\rho_{0} \sigma_{0}\right)^{x-1} \leqq C_{1}^{-1}$, then (7.4) ${ }_{j+1}$ holds.

To (7.5) and (7.22) we add the inductive hypothesis

$$
\begin{equation*}
t_{j}^{\mu-2 s-6} b^{j} \mathrm{~K}_{j} \mathrm{~K}_{j}^{*} \mathrm{~N}_{j}(1) \mathrm{L}_{j}(\mu)<\frac{1}{2} \mathbf{M} \tag{7.31}
\end{equation*}
$$

Clearly (7.22) ${ }_{j}$ and (7.31) imply (7.22) ${ }_{j+1}$. By (7.21) and (7.22), $\theta_{j} \mathrm{~N}_{j}^{2 k} \leqq t_{j}^{s} \mathrm{~N}_{j}^{2 k-1}$, so that (7.5) may be replaced by the stronger

$$
\begin{equation*}
c(k) t_{j}^{s} \mathbf{N}_{j}^{2 k-1} \sigma_{j}^{-1}<1 \tag{7.32}
\end{equation*}
$$

But the ratio of the left hand side of (7.32) at $j+1$ to that at $j$ does not exceed $\left(c_{b}^{2-x} t_{j}^{x-1}\right)^{s} \hat{c}(k)^{2 k-1} 5$. This will be less than one if $\rho_{0}$ is chosen sufficiently small, and (7.32) ${ }_{j}$ will imply (7.32) ${ }_{j+1}$.

Finally, consider (7.31) $\boldsymbol{j}_{j+1}$. If we set

$$
\mathrm{Q}_{j} \equiv b^{j} \mathrm{~K}_{j} \mathrm{~K}_{j}^{*} \mathrm{~N}_{j}(1) \mathrm{N}_{j}^{k-3+\mu},
$$

then by (7.28), (7.24-27), and $\theta_{j} \leqq \mathrm{M} t_{j}^{s}$,

$$
\begin{align*}
t_{j+1}^{\mu-2 s-6} b^{j+1} \mathrm{~K}_{j+1} \mathrm{~K}_{j+1}^{*} & \mathrm{~N}_{j+1}(1) \mathrm{L}_{j+1}(\mu) \\
& \leqq c(k, b) \mathrm{MQ}_{j} x_{j}^{(\mu-2 s-6)}\left(t_{j}^{-\mu+2 s+6}+t_{j}^{s-\mu}\right) . \tag{7.33}
\end{align*}
$$

We require that both exponents of $t_{j}$ be positive [note (7.29)]

$$
\mu>2 s+6, \quad(\kappa-1) \mu>3 s+10 .
$$

Then (7.33) is bounded by $c(k, b) \mathrm{MQ}_{j} t_{j}^{2 \alpha}, \alpha>0$. But

$$
\frac{\mathrm{Q}_{j+1} t_{j+1}^{\alpha}}{\mathrm{Q}_{j} t_{j}^{\alpha}} \leqq c(k, b) t_{j}^{(\mathrm{x}-1) \alpha}
$$

which can be made arbitrarily small uniformly in $j$ by shrinking $\rho_{0}$. Hence, if $\rho_{0}$ is taken sufficiently small $(7.31)_{j+1}$ will hold. This completes the induction step.

It remains to verify the above conditions for $j=0$. For this we take $t_{0}=\left(c_{1}^{-1} \rho_{0} \sigma_{0}\right)^{\beta}, \beta \geqq 1$. Then (7.4) ${ }_{0}$ holds and (7.32) $)_{0}$ will hold if $\rho_{0}$ is small enough. We can achieve (7.31) $)_{o}$ by shrinking $\rho_{0}$ if [see (5.14) and (5.18)]

$$
\begin{equation*}
\beta(\mu-2 s-6)>2 \mu+4(n+k)-16 . \tag{7.34}
\end{equation*}
$$

Likewise, by (6.5) and (7.21) we can get (7.22) if

$$
\begin{equation*}
m-k+3>2 \beta(s+3)+4(n+k)-16 . \tag{7.35}
\end{equation*}
$$

Thus, if all parameters are chosen as indicated and $\rho_{0}>0$ is sufficiently small, the construction is possible for all $j$ and yields a sequence of embeddings $Z_{j} \circ{ }^{\circ}{ }_{j}$ of the original CR structure which converges in $\mathrm{C}^{k}$ norm on a neighborhood of 0 in $\mathbb{R}^{2 n-1}$ to a holomorphic embedding. We may take $s=1, \kappa=\frac{7}{4}$, and $\mu=18$. Since $\mu \leqq k-3$, we need $k \geqq 21$. By the construction of section 1 we may take $\mathbf{Z}_{0}$ to be a polynomial. For $j \geqq 1$ the $Z_{j}$ as constructed in section 3 are $\mathbf{C}^{\infty}$. Hence $\mathrm{N}_{j}(\mu)$ is finite, and by (1.11) $\bar{\partial}_{\mathrm{x}^{j}} \mathrm{Z}_{j}$ is as smooth as $\left(\mathrm{A}^{j}, \mathrm{~B}^{j}\right)$, i.e. of class $\mathrm{C}^{m}$. By (6.25) $\mathrm{L}_{j}(\mu)$ will be finite if $m \geqq k-3+\mu=k+15$. If we choose $\beta=2+n+k$, then (7.34) holds and (7.35) becomes $m \geqq 6 k+5 n-2$. This implies $m \geqq k+15$ since $k \geqq 21$, and the theorem is proved.

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