## WICHITA STATE UNIVERSITY

A LOCAL EXISTENCE THEOREM FOR QUASI-LINEAR
PARTIAL DIFFERENTIAL EQUATIONS*
by
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## MATHEMATICS DEPARTMENT

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## I. Introduction.

We investigate here conditions on a quasi-linear partial differential equation of second order sufficient to guarantee existence of a smooth solution. More precisely we consider an equation

$$
\begin{equation*}
A_{i j} u_{x_{i} x_{j}}+D=0 \tag{1}
\end{equation*}
$$

in a bounded region $R$ in $E^{n}$ subject to the condition

$$
u=0 \quad \text { on } \quad \partial R .
$$

We use the subscript $x_{i}$ to denote partial differentiation with respect to $\mathrm{X}_{\boldsymbol{i}}$ and make free use of the summation convention on repeated indices (two terms with the same index are summed over the common index). The coefficients $A_{i j}, D$ are functions of $2 n$ variables, the first $n$ being $x=\left(x_{1}, \ldots, x_{n}\right)$, and the second $n$ being $u_{x}=\left(u_{x_{1}}, \ldots, u_{x_{n}}\right)$.

We will seek solutions which together with their derivatives through second order, are uniformly Hölder continuous, that is which belong to the Banach Space $C_{2+\alpha}(R)$, where $\alpha$ is a positive number less than one. Therefore we assume immediately that $\partial R$ is sufficiently smooth, that is of class $C_{2+\alpha}$ (see [1]), and that the functions $A_{i j}$ äre "uniformly elliptic" with respect to elements of $C_{2+\alpha}(R)$, that is

$$
A_{i j}\left(x, u_{x}(x)\right) \lambda_{i} \lambda_{j} \geq \text { const. }|\lambda|^{2}, \lambda=\left(\lambda_{1}, \ldots, \lambda_{n}\right) \varepsilon E^{n}
$$

where the constant may depend on the $u \in C_{2+\alpha}(R)$ appearing but does not depend on $x$.

The approach taken will be to give sufficient conditions for the differential operator involved to map one Banach Space of Hölder continuous functions into another in a way sufficiently nice for the theorem of Kantorovich on converaance of Newtons method in Banach Spaces to be applied.

Before proceeding to the proof of our main theorem we must. make a survey of some prerequisite material.
II. Preliminaries.

The work done below will be carried out in two Sanach Spaces of Holder continuous functions, $C_{2+\alpha}(R)$ and $C_{\alpha}(R)$. The latter consists of functions uniformly Holder continuous in $p$, and is normed by

$$
|\phi|_{\alpha}=|\phi|_{0}+H(\phi)
$$

where $1 l_{0}$ denotes the maximum norm over ? and

$$
H(\phi)=\sup _{x, y \in R!} \frac{|\phi(x)-\phi(y)|}{|x-y|^{\alpha}} .
$$

The former consists of functions which, together with their derivatives through second order, belong to $C_{\alpha}($ 吅 and is normed by

$$
|\phi|_{2+\alpha}=|\phi|_{\alpha}+\text { sum of } C_{\alpha} \text { norms of derivatives through 2nd order. }
$$ The reader is referred to [1] for relevant facts about these spaces. We need two results on linear elliptic equations, which follow from theorems aiven in [1], the first an a priori estimate, the second an existence theorem. He consider an equation

$$
a_{i j} u_{x_{i} x_{j}}+b_{i} u_{x_{i}}=f
$$

in $R$ together with prescription of zero boundary values. The functions
$a_{i j}, b_{i}, f$ are assumed to be elements of $C_{\alpha}(R)$, and it is assumed that

$$
a_{i j} \lambda_{i} \lambda_{j} \geq m|\lambda|^{2} \quad \lambda=\left(\lambda_{1}, \ldots, \lambda_{n}\right) \in \dot{E}^{n}
$$

where $m$ is a positive number independent of $x$. The needed results are
(1) There is a constant $C$ depending only on $\alpha, R, m$ and an upper bound on the $c_{\alpha}$ norms of $a_{i j}, b_{i}$ such that for any solution $u$ of the above linear problem that is an element of $C_{2+\alpha}(R)$ the inequality

$$
|u|_{2+\alpha} \leq c|f|_{\alpha}
$$

holds.
and
(2) Under the above assumptions for any $f$ in $C_{\alpha}(R)$ there is a solution $u$ in $C_{2+\alpha}(R)$ of the above linear problem. It is worth noting here that a solution of the above problem is unique by the maximum principle. (See [2].)

We assume that the reader is familiar with elementary facts about Frechét derivatives, which we will call F-derivatives. (These are available in [3] or [4] for example.) If $X$ and $Y$ are Banach Spaces, we denote by $B(X, Y)$ the space of continuous linear overtors from $X$ into $Y$, and by $B(X, X ; Y)$ the space of continuous bilinear operators from $X \times X$ into $Y$.

The theorem of Kantorovich on convergance of Newtons method is given below. A proof can be found in [4].

Theorem. (Kantorovich) Suppose $P$ is a continuous mapping of the sphere $A=\left\{x \mid\left\|x-x_{0}\right\|<R\right\}$ in the Banach space. $X$ into the
Banach space $Y$, and that $P$ is twice continuously differentiable
on the closed sphere
$A_{0}=\left\{x \mid\left\|x-x_{0}\right\| \leq r\right\}$, where $r<R$. Further, suppose that
(1) $\left[P^{\prime}\left(x_{0}\right)\right]^{-1}$ exists,
(2) $\left\|\left[P^{\prime}\left(x_{0}\right)\right]^{-1}\right\| \leq B$,
(3) $\left\|P\left(x_{0}\right)\right\| \leq n$,
(4) $\left\|P^{\prime \prime}(x)\right\| \leq K$ for $x$ in $A$.

Then, if $h=K B^{2} n<1 / 2$ and $r_{0}=(1-\sqrt{1-2 h}) B_{n} / h \leq r$, the equation

$$
P(x)=0
$$

has a solution $x^{*}$ in the sphere $\left\{x \mid\left\|x_{0}-x\right\| \leq r_{0}\right\}$ to which the
sequence

$$
x_{n+1}^{\prime}=x_{n}^{\prime}-\left[P^{\prime}\left(x_{0}\right)\right]^{-1} P\left(x_{n}^{\prime}\right)
$$

converges. In addition, the sequence

$$
x_{n+1}=x_{n}-\left[P^{\prime}\left(x_{n}\right)\right]^{-1} P\left(x_{n}\right)
$$

is well defined and converges to $x^{*}$. The following error bounds hold:

$$
\begin{aligned}
& \left\|x^{*}-x_{n}\right\| \leq(2 h)^{2 n} n / h 2^{n}, \frac{\text { and }}{} \\
& \left\|x^{*}-x_{n}^{\prime}\right\| \leq(1-\sqrt{1-2 h})^{n+1} n / h .
\end{aligned}
$$

If

$$
\therefore \quad r<r_{1}=(1+\sqrt{1-2 h}) B n / h .
$$

the solution $x^{*}$ is unique in $A_{0}$.

The space $Y$ which we will utilize is just $C_{\alpha}(R)$. A moments reflection shows that the functions in $C_{2+\alpha}(R)$ which vanish on $\partial R_{\text {, }}$ form a Banach Space and this is the space $X$ which we use. The operator. $P$ is defined by letting $P(u)$ be the left hand side of (1) with $u$ substituted. Therefore we have translated our problem into that of existence of a zero in $X$ for the operator $P$.

## III. The Existence Theorem.

The main task in what follows is to show that under appropriate conditions the operator $P$, and the spaces $X$ and $Y$ defined above, satisfy the conditions needed to apply the Kantorovich theorem. The first step is to see that $P$ does in fact map $X$ into $Y$. We must have

$$
\begin{equation*}
P(u)=A_{i j} u_{x_{i} x_{j}}+D \varepsilon Y \tag{2}
\end{equation*}
$$

for $u$ in $X$. Since $Y$ is closed under (pointwise) multiplication, it suffices to show that $A_{i j}\left(x, u_{x}(x)\right), D\left(x, u_{x}(x)\right)$ are uniformly Hölder continuous for $u$ in $x$. Letting $F$ denote anyone of $A_{i j}$, or $D$, we propose the following sufficient condition on $F$ in order that this be true:
E)

$$
|F(x, \xi)-F(y, n)| \leq K\left[|x-y|^{\alpha}+|\xi-n|\right]
$$

for $x, y \in \bar{R} \quad|\xi|,|\eta| \leq c$ where $K$ depends. only on $c$.

Lemma 1. If each of $A_{i j}, D$ satisfy $E, P$ maps $X$ into ${ }^{*} Y$.
PROOF: Using $F$ as above, and applying $E$, we have

$$
\left|F\left(x, u_{x}(x)\right)-F\left(y, u_{x}(y)\right)\right| \leq K\left[|x-y|^{\alpha}+\left|u_{x}(x)-u_{x}(y)\right|\right]
$$

where $K$ is determined by $U$. Then, since $u \varepsilon \cdot X$,

$$
\mid F\left(x, u_{x}(x)\right)-F\left(y, u_{x}(y)\left|\leq k\left[|x-y|^{\alpha}+k_{1}|x-y|^{\alpha}\right]=k\left(1+k_{1}\right)\right| x-\left.y\right|^{\alpha},\right.
$$

and the lemma is proven.

In order to investigate continuity of $P$, suppose that we have chosen $u^{0}$ in $X$ and consider

$$
\begin{equation*}
P(u)-P\left(u^{0}\right)=A_{i j} u_{x_{i} x_{j}}-A_{i j}^{0} u_{x_{i} x_{j}}^{0}+D-D^{0} \tag{3}
\end{equation*}
$$

where we have denoted the coefficient with $u^{0}$ substituted by the same superscript. We let

$$
\begin{aligned}
& A_{i j}=A_{i j}^{0}+\Delta A_{i j}, \\
& D=D^{0}+\Delta D, \\
& u=u^{0}+h
\end{aligned}
$$

and (3) becomes

$$
P(u)-P\left(u^{0}\right)=\Delta A_{i j} u_{x_{i} x_{j}}^{0}+A_{i j}^{0} h_{x_{i} x_{j}}+\Delta A_{i j} h_{x_{\mathbf{i}} x_{j}}
$$

so that

$$
\left|P(u)-P\left(u^{0}\right)\right|_{\alpha} \leq\left|\Delta A_{i j}\right|_{\alpha}\left|u_{x_{i} x_{j}}^{0}\right|_{\alpha}+\left|A_{i j}^{0}\right|_{\alpha}\left|h_{x_{i} x_{j}}\right|_{\alpha}+\left|\Delta A_{i j}\right|_{\alpha}\left|h_{x_{i} x_{j}}\right|_{\alpha}
$$

Clearly then, a condition which guarantees that

$$
\left|\Delta A_{i j}\right|_{\alpha} \text { and }|\Delta D|_{\alpha} \text { go to zero as }|h|_{2+\alpha}
$$

does will imply continuity of $P$ at $u^{0}$. We assert that it suffices to assume that the first and second derivatives of $A_{i j}$ and $D$ with respect to their "gradient" variables satisfy the condition $E$ given above. For brevity we introduce the class of functions $G_{m}$ : a realvalued function with domain $\bar{R} \times E^{n}$ which has continuous derivatives
through the moth order with respect to its last $n$ variables and which is such that it and all of these derivatives satisfy condition $E$ is in $G_{m}$.

Lemma 2. If $A_{i j}$ and $D$ are in $G_{2}, P$ is continuous on $X$.
PROOF: Retaining the above notation, and letting $F$ represent any of $A_{i j}$, $D$ as before, it suffices to show that

$$
|\Delta F|_{\alpha} \rightarrow 0 \text { as }|h|_{2+\alpha} \rightarrow 0
$$

The assumptions on $F$ clearly allow application of Taylors' theorem and imply the validity of

$$
\begin{equation*}
\Delta F=F_{u_{x_{i}}}\left(x, u_{x}^{0}\right) h_{x_{i}}+F_{u_{x_{i}}} u_{x_{j}}\left(x, u_{x}^{0}+\theta h_{x}\right) h_{x_{i}} h_{x_{j}}, \tag{4}
\end{equation*}
$$

where $\theta$ is a positive number less than one.
Since functions which satisfy $E$ yield elements of $Y$ when elements of $X$ are substituted in them, we have

$$
|\Delta F|_{\alpha} \leq\left|F_{u_{x_{i}}}^{0}\right|_{\alpha}\left|h_{x_{i}}\right|_{\alpha}+\left|F_{u_{x_{i}}}^{\theta} u_{x_{j}}\right|_{\alpha}\left|h_{x_{i}}\right|_{\alpha}\left|h_{x_{i}}\right|_{\alpha}
$$

where $F_{u_{x_{i}}}^{\theta}{ }_{u_{j}}$ denotes the appropriate function from (4). It remains to show that there is an $r>0$ such that $\left|F_{u_{x_{i}}}^{\theta}{ }_{u_{x_{j}}}\right|_{\alpha}$ is bounded for $|h|_{2+\alpha}<r$. Let $G$ denote the function $F_{u_{x_{i}}} u_{x_{j}}(\cdot, \cdot)$. We have, since $F \in \dot{\dot{G}_{2}}$,

$$
\begin{aligned}
\left|G\left(x, u_{x}^{0}(x)+\theta h_{x}(x)\right)-G\left(y, u_{x}^{0}(y)+\theta h_{x}(y)\right)\right| & \leq \\
\left|G\left(x, u_{x}^{0}(x)+\theta h_{x}(x)\right)-G\left(x, u_{x}^{0}(x)+\theta h_{x}(y)\right)\right| & +\mid G\left(x, u_{x}^{0}(x)+\theta h_{x}(y)\right) \\
& -G\left(y, u_{x}^{0}(y)+\theta h_{x}(y)\right) \mid
\end{aligned}
$$

which is not greater than

$$
0\left|h_{x}(x)-h_{x}(y)\right|+k|x-y|^{\alpha}+k\left|u_{x}^{0}(x)-u_{x}^{0}(y)\right|
$$

for $x, y$ in $\bar{R}$. Dividing by $|x-y|^{\alpha}$ and taking the supremum for $x, y$ in $R$ we get the desired result.

We are now ready to investigate differentiability properties of the operator $P$. To this end we note that assuming each of $A_{i j}$, $D$ belong to $G_{3}$ guarantees the validity of

$$
\begin{aligned}
F\left(x, u_{x}+\right. & \left.h_{x}\right)=F\left(x, u_{x}\right)+F_{u_{x_{i}}}\left(x, u_{x}\right) h_{x_{i}}+F_{u_{x_{i}} u_{x_{j}}}\left(x, u_{x}\right) h_{x_{i}} h_{x_{j}} \\
& +F_{u_{x_{i}}} u_{x_{j}}{ }^{u_{x_{k}}}\left(x, u_{x}+\theta h_{x}\right) h_{x_{i}} h_{x_{j}}{ }^{h} x_{k}, 0<\theta<1
\end{aligned}
$$

where $F$ represents any of $A_{i j}, D$, and $u, h \varepsilon X$. Then we have

$$
\begin{aligned}
& P(u+h)=A_{i j}\left(x, u_{x}+h_{x}\right) u_{x_{i} x_{j}}+A_{i j}\left(x, u_{x}+h_{x}\right) h_{x_{i} x_{j}}+D\left(x, u_{x}+h_{x}\right) \\
& =A_{i j}\left(x, u_{x}\right) u_{x_{i} x_{j}}+\left(A_{i j}\right) u_{x_{k}} u_{x_{i} x_{j}}{ }^{h} x_{k}+\left(A_{i j}\right) u_{x_{k}}{ }^{u_{x}}{ }{ }^{u_{x_{i}} x_{j}}{ }^{h} x_{k}{ }^{h} x_{1} \\
& +\left(A_{i j}^{\theta}\right)_{u_{x_{k}}} u_{x_{1}}{ }^{u} x_{m}{ }_{u_{x_{i}} x_{j}}{ }^{h} x_{k}{ }^{h} x_{1}{ }^{h} x_{m} \\
& +A_{i j}\left(x, u_{x}\right) h_{x_{i} x_{j}}+\left(A_{i j}\right)_{u_{x_{k}}} h_{x_{i} x_{j}}{ }^{h} x_{k}+\left(A_{i j}\right)_{u_{x_{k}}} u_{x_{1}}{ }^{h} x_{i} x_{j}{ }^{h} x_{k}{ }^{h} x_{1} \\
& +\left(A_{i j}^{\theta}\right)_{u_{x_{k}}} u_{x_{1}} u_{x_{m}}{ }^{h} x_{i} x_{j}{ }^{h} x_{k}{ }^{h} x_{1}{ }^{h} x_{m} \\
& +D\left(x, u_{x}\right)+D_{u_{x_{i}}} h_{x_{i}}+D_{u_{x_{i} x_{j}}} h_{x_{i}}{ }^{h}{x_{j}}+D_{u_{x_{i}}}^{\theta} u_{x_{j}}{ }^{u_{x_{k}}}{ }^{h} x_{i}{ }^{h} x_{j}{ }^{h} x_{k}
\end{aligned}
$$

so that

$$
P(u+h)=P(u)+L_{u}(h)+B_{u}(h, h)+R_{u}(h)
$$

where

$$
\begin{align*}
& L_{u}(h)=A_{i j} h_{x_{i} x_{j}}+\left(A_{i j}\right)_{u_{x_{k}}} u_{x_{i} x_{j}}{ }^{h} x_{k}+D_{u_{x_{i}}} h_{x_{i}},  \tag{5}\\
& B_{u}(h, h)=\left(A_{i j}\right)_{u_{x_{k}} u_{x_{1}}} u_{x_{i} x_{j}}{ }^{h} x_{k}{ }^{h} x_{j}+\left(A_{i j}\right){ }_{u_{x_{k}}}{ }_{x_{i} x_{j}}{ }^{h} x_{k}+D_{u_{x_{i}}} u_{x_{j}}{ }^{h} x_{i} h_{x_{j}}, \tag{6}
\end{align*}
$$

and

$$
\begin{align*}
R_{u}(h) & =\left(A_{i j}^{\theta}\right)_{u_{x_{k}}} u_{x_{1}} u_{x_{m}}\left[u_{x_{i} x_{j}}+h_{x_{i} x_{j}}\right]{ }^{h} x_{k}{ }^{h} x_{1}{ }^{h} x_{m}  \tag{7}\\
& +D_{u_{x_{i}}{ }^{u} x_{j}}{ }^{u}{ }_{x_{k}}{ }^{h} x_{x_{i}}{ }^{h} x_{j}{ }^{h} x_{k}+\left(A_{i j}\right){ }_{u_{x_{k}}} u_{x_{1}}{ }^{h} x_{x_{i}} x_{j}{ }^{h} x_{k}{ }^{h} x_{1}
\end{align*}
$$

A quick check shows that, for fixed $u$ in $x, L_{u}$ and $B_{u}$ are elements of $B(X, Y), B(X, X ; Y)$, respectively. Therefore, application of the reasoning used at the end of the proof of lemma 2 to the terms of $R_{u}(h)$ with the superscript $\theta$ proves.

Lemma 3. If $A_{i j}, D$ are in $G_{3}, P$ is $F$-differentiable at each $u$ in $X$, and $P^{\prime}(u)$ is given by (6).

To proceed we must guarantee continuity of the mapping

$$
u \rightarrow p^{\prime}(u)
$$

## Consider

$$
\begin{aligned}
{\left[P^{\prime}(u)-P^{\prime}(v)\right](h) } & =\left(A_{i j}^{u}-A_{i j}^{v}\right) h_{x_{i} x_{j}}+\left[\left(A_{i j}^{u}\right)_{u_{x_{k}}} u_{x_{i} x_{j}}-\left(A_{i j}^{v}\right)_{u_{x_{k}}} v_{x_{i} x_{j}}\right] h_{x_{k}} \\
& +\left(D_{u_{x_{k}}}^{u}-D_{u_{x_{k}}}^{v}\right) h_{x_{k}},
\end{aligned}
$$

where a function $F$ with say $u$ substituted is denoted $F^{u}$. It
follows from the above that

$$
\begin{aligned}
\left\|P^{\prime}(u)-P^{\prime}(v)\right\| \leq C \max \left\{\left|A_{i j}^{u}-A_{i j}^{v}\right|_{\alpha},\right. & \left|\left(A_{i j}^{u}\right)_{u_{x_{k}}} u_{x_{i} x_{j}}-\left(A_{i j}^{v}\right)_{u_{x_{k}}} v_{x_{i} x_{j}}\right|_{\alpha}, \\
& \left.\left|0_{u_{x_{k}}^{u}}^{u}-D_{u_{x_{k}}}^{v}\right|_{\alpha}\right\},
\end{aligned}
$$

where $C$ is independent of $u, v$. We must insure that this quantity go to zero as $v$ approaches $u$ in the norm of $X$. That this is true for the first and third terms in brackets follows from the assumption that $A_{i j}, D$ belong to $G_{3}$ and the content of lemma 2. In order to deal with the other term we denote $\left(A_{i j}\right)_{u_{x_{k}}}$ by $F$, observe that

$$
\left|F^{u} u_{x_{i} x_{j}}-F^{v} v_{x_{i} x_{j}}\right|_{\alpha} \leq\left|u_{x_{i} x_{j}}\right|_{\alpha}\left|F^{u}-F^{v}\right|+\left|F^{v}\right|_{\alpha}|v-u|_{2+\alpha},
$$

and again use the reasoning applied at the end of the proof of lemma 2. Thus the hypothesis of lerma 3 yields continuous differentiability as well. Since $\left|P_{u}(h)\right|_{\alpha}$ is clearly of order 3 with respect to the norm of $X$, it now follows that $P$ has a second Frechét derivative at each $u$, which is identified in the usual way (see either [3] or [4]) with the bilinear operator $B_{u}$. In fact, an unexciting return to the ground just covered shows that with no further assumptions $P$ is twice continuously differentiable.

Corollary (to lemma 3). If $A_{i j}, D \in G_{3}, P$ is twice continuously F-differentiable.

We now observe that the assumptions we made in the Introduction guarantee that the operator $\mathrm{P}^{\prime}(\mathrm{u})$ is invertible for each $u$ in $X$. In fact, the regularity assumptions, and the "uniform ellipticity" assumption made on the $A_{i j}$ in that section allow application of the
existence theorem for linear equations given in section 2, so that we know $P^{\prime}(u)$ maps $X$ onto $Y$. Then the maximum principle implies that $P^{\prime}(u)$ is one-to-one, and the closed graph theorem guarantees continuity of the inverse operator. We have

Lemma 4. The uniform ellipticity assumption implies that $P^{\prime}(u)$ is invertible for each $u$ in $X$.

To complete the preparation for use of the Kantorovich theorem, we give the following lemma which says, roughly, that the norm of $P^{\prime \prime}(u)$ is bounded on spheres.

Lemma 5. For each $u$ in $x, r>0$ there is a number $k=k(u, r)$ such that

$$
\left\|P^{\prime \prime}(u)\right\| \leq K
$$

for $u$ in $S(u, r)$.

PROOF: Looking back at the expression giving $P^{\prime \prime}(u)$ in (6), we see that one further application of the reasoning given at the end of the proof of lemma 2 suffices.

By simply looking applying the Kantorovich Theorem to the operator $P$ we can now deduce the following theorem.

Theorem. If $u \varepsilon \times$ and $r>0$ are such that

$$
\therefore K(u, r)\left\|\left[P^{\prime}(u)\right]^{-1}\right\|^{2}|P(u)|_{\alpha}<1 / 2 \text {, }
$$

and

$$
1 / K(u, r)\left\|\left[P^{\prime}(u)\right]^{-1}\right\| \leq r \text {, }
$$

then there is an element $u^{*}$ in $S(u, r)$ such that
(8)

$$
P\left(u^{\star}\right)=0
$$

to which both the Newton sequence and the modified Newton sequence (Begun at u) converge.

Of course, the uniformly ellipticity condition, and the assumption that the coefficients in $P$ are of class $G_{3}$ have been implicitly assumed in order to obtain this theorem.

We now deduce the following
Corollary. For each $u$ in $X$ there is a number $H$ such that

$$
|P(u)|_{\alpha}<H
$$

implies that there is a solution of the equation (8) in a sphere around $u$ in $x$.

PROOF: We first need to observe that, for fixed $u$, the second of the inequalities needed in the above theorem can always be fulfilled by taking, if necessary, a larger upper bound for $\left\|P^{\prime \prime}(u)\right\|$. More precisely, denoting $\left\|\left[P^{\prime}(u)\right]^{-1}\right\|$ by $B$, for any given $r$ we need only guarantee that

$$
K \geq 1 / B r
$$

which is certainly possible since $K$ is an upper bound. Then observing the first inequality needed in the theorem we deduce that

$$
\therefore \quad|P(u)|_{\alpha}<1 / 2 K B^{2}
$$

implies that there is a solution of ( 8 ) in $S(u, r)$, and the corollary is proven.

Remark 1. The conclusions of the Kantorovich theorem could also be used to deduce a "local uniqueness" result. However, for the partial differential equation under consideration a solution can be
shown to be unique using the maximum principle (see for example [2]), so this is of limited interest here.

Remark 2. The theorem and corollary proven above say something essentially about retention of existence under a perturbation. To see this more clearly, suppose we assume for the moment that the function $u$ in the above is the function identically zero. Also, note that when

$$
D(x, 0) \equiv 0
$$

the function identically zero is a solution of the equation being studied, Therefore, the corollary says that under a sufficiently small perturbation of the function

$$
D(x, 0)=P(0)
$$

from zero (in the norm of $X$ ) there still exists a solution (in $X$ ) of the equation.

Remark 3. He give here a more concrete class of coefficients which is contained in the class $G_{m}$ utilized above. This class consists of functions with domain $\bar{R} \times E^{n}$ which are polynomials in the last $n$ variables with coefficients that are uniformly Holder continuous functions of the first $n$. More precisely we consider functions

$$
F(x, \xi)=a_{i \ldots j}(x) \xi_{1}^{v_{i}} \ldots \xi_{n}^{v} j
$$

(recall the summation convention) where $v_{i}, \ldots, v_{j}$ vary, over the lowest through the highest powers appearing on $\xi_{1}, \ldots, \xi_{n}$, respectively, and each function $a_{i} \ldots j$ is an element of $Y$.

To see that such a function is in $G_{m}$ we note that this new class is closed under partial differentiations with respect to any of the last $n$ variables, and that

$$
\begin{aligned}
F(x, \xi)-F(y, n) & =\left[a_{i} \ldots j(x)-a_{i} \ldots j(y)\right] \xi_{1}^{v_{i}} \ldots \xi_{n}^{v_{j}} \\
& +a_{i \ldots j}(y)\left[\xi_{1}^{v_{i}} \ldots \xi_{n}^{v_{j}}-\eta_{1}^{v_{i}} \ldots n_{n}^{v_{j}}\right]
\end{aligned}
$$

The result follows from Holder continuity of $a_{i} \ldots j$ and the fact that a continuously differentiable function of $n$ variables satisfies a uniform Lipschitz condition on any compact set (the Lipschitz constant depending, of course, on the compact set.).

We will deal with an equation whose coefficients are of this type in the appendix.

Appendix. An example: the Minimal Surface Equation.
We apply the results obtained above to the $n$-dimensional "minimal surface" equation. We seek a function $u$ satisfying

$$
\begin{equation*}
M(u) \equiv\left(1+u_{i} u_{i}\right) \Delta u-u_{i} u_{j} u_{i j}=0 \tag{9}
\end{equation*}
$$

in $R$, and

$$
u=\phi
$$

on $\partial R$. We have used here and will use henceforth the notation

$$
u_{i}=\partial u / \partial x_{i}, u_{i j}=\partial^{2} u / \partial x_{i} \partial x_{j},
$$

that is partial derivatives are indicated simply by subscripts.
In order to pose this problem in the setting of the previous paragraphs we assume immediately that the given boundary function is the restriction to $\partial R$ of some function in $C_{2+\alpha}(R)$ which we will also call $\phi$. Then, setting

$$
v=u-\phi
$$

we obtain from (9) the equation

$$
P(v) \equiv M(v+\phi)=0
$$

for $v$ subject to a homogeneous Dirichlet boundary condition. Clearly, we may think of $P$ as an operator mapping $X$ into $Y$, and a solution of the operator equation

$$
P(v)=0
$$

in $X$ is a solution of our problem. In fact, the last remark of paragraph 3 together with the lemmas of that paragraph show that $P$ is twice continuously F-differentiable on $X$, and that the norm of its second derivative is bounded on spheres. In order to make use of the theorem of paraqraph 3 we need only show that the F-derivative of $P$ is a (uniformly) elliptic operator for each $v$ in $X$. We have

$$
P(v)=N(v)+L(v)+M(\phi)
$$

where the operator $: 1$ is defined in (9), and

$$
\begin{aligned}
& N(v)=\left(v_{i} v_{i}+2 \phi_{i} v_{i}\right) \Delta v-v_{i} v_{j} v_{i j}-2 \phi_{i} v_{j} v_{i j}+\Delta \phi v_{i} v_{i}-\phi_{i j} v_{i} v_{j}, \\
& L(v)=\left(1+\phi_{i} \phi_{i}\right) \Delta v-\phi_{i} \phi_{j} v_{i j}-2 \phi_{i} \phi_{i j} v_{j}+2 \Delta \phi \phi_{i} v_{i} .
\end{aligned}
$$

It is clear that $M(\phi)$ is an element of $Y$ and that the linear operator $L$ defined above is an element of $B(X, Y)$. Thus, the F-derivative of $P$ is that of $N$ added to $L$. (The derivative of a linear operator is the operator itself and the derivative of a. is the sum of the derivatives.). Routine manipulations show that

$$
M(v+h)=H(v)+Q_{v}(h)+R_{v}(h)
$$

where

$$
\begin{aligned}
Q_{v}(h) & =\left(v_{i} v_{i}+2 \phi_{i} v_{i}\right) \Delta h-\left(v_{i} v_{j}+2 \phi_{i} v_{j}\right) h_{i j}+2 \Delta v\left(v_{i} h_{i}+\phi_{i} h_{i}\right) \\
& +2 v_{i j}\left(v_{i} h_{j}+\phi_{i} h_{j}\right)+2 \Delta \phi v_{i} h_{i}-2 \phi_{i j} v_{i} h_{j}
\end{aligned}
$$

and $R_{v}(h)$, for fixed $v$, vanishes quadratically in $|h|_{2+\alpha}$ as $|h|_{2+\alpha}$ goes to zero. Since the operator $\eta_{v}$ so defined is an. element of $B(X, Y)$ we have

$$
\eta_{v}=N^{\prime}(v) .
$$

Therefore the principle part (the part involving only second partials of $h$ ) of $H^{\prime}(v)$ is given by

$$
\operatorname{pr} H^{\prime}(v)(h)=\left(v_{i} v_{i}+2 \phi_{i} v_{i}\right) \Delta h-\left(v_{i} v_{j}+2 \phi_{i} v_{j}\right) h_{i j}
$$

and that of $L$ by

$$
\operatorname{pr} L(h)=\left(1+\phi_{i} \phi_{i}\right) \Delta h-\phi_{i} \phi_{j} h_{i j}
$$

so that

$$
p r p^{\prime}(v)(h)=\left(1+\left(\phi_{\mathbf{i}}+\dot{v}_{i}\right)\left(\phi_{\mathbf{i}}+v_{i}\right)\right) \Delta h-\left(\phi_{\mathbf{i}}+v_{i}\right)\left(\phi_{\mathbf{j}}+v_{j}\right) h_{i j}
$$

Setting $u=\phi+v$, the relevant quadratic form is given by

$$
g(\lambda)=\left(1+u_{i} u_{i}\right)|\lambda|^{2}-u_{i} u_{j} i_{i}{ }_{j}
$$

Since, by the Schwartz inequality,

$$
u_{i} u_{j}^{\lambda} i_{i}=\left(u_{i} \lambda_{i}\right)^{2} \leq\left(u_{i} u_{i}\right)|\lambda|^{2}
$$

we have

$$
q(\lambda) \geq|\lambda|^{2}
$$

and the required uniform ellipticity of the $F$-derivative is demonstrated.

We are now able to apply the theorem of the previous paragraph to the operator $P$. We restrict ourselves to asking whether a solution exists in a neighborhood in $X$ of 0 for the operator equation

$$
\begin{equation*}
P(v)=0 . \tag{10}
\end{equation*}
$$

Let $K=K(r)$ be such that

$$
\left\|P^{\prime \prime}(v)\right\| \leq k, v \in S(0, r)
$$

Then the above mentioned theorem quarantees that if

$$
K B^{2}|P(0)|_{\alpha}<1 / 2,
$$

and

$$
1 / K B \leq r,
$$

where

$$
B=\left\|\left[P^{\prime}(0)\right]^{-1}\right\|
$$

there exists a solution of (10) in $S(0, r)$. We force the second of the above inequalities to hold simply by choosing

$$
K \geq 1 / B r,
$$

that is by choosing $K$ to be the maximum of $1 / B r$ and some upper bound for $\left\|P^{\prime \prime}(v)\right\|$ on $S(0, r)$. Then, noting that

$$
P(0)=: 1(\phi)
$$

we know there is a solution of (10) in $S(0, r)$ whenever

$$
|M(\phi)|_{\alpha}<1 / 2 K B^{2}
$$

We can now state

Theorem. For each $r>0$ there is a number $H$ depending only on $r$ and a constant bounding the $C_{\alpha}$ norms of the first and second de-
rivatives of $\phi$ such that

$$
|M(\phi)|_{\alpha}<H
$$

implies existence of a solution $u$ of the equation (9) such that $u$ is in $C_{2+\alpha}(R)$, and

$$
|u-\phi|_{2+\alpha}<r .
$$

PROOF: We will, in fact, give an $H$ explicitly in terms of $r, B$, and $C_{\alpha}$ norms of the first and second derivatives of $\phi$. The result then follows from the fact that $B$ is the smallest constant possible in the a priori estimate

$$
|h|_{2+\alpha} \leq \text { const. }\left|P^{\prime}(0)(h)\right|_{\alpha}
$$

which holds for the linear elliptic boundary value problem

$$
P^{\prime}(0)(h)=f \quad(f \text { in } Y)
$$

since these constants depend only on an upper bound for the $C_{\alpha}$ norms of the coefficients of the differential operator $P^{\prime}(0)$ (which are quadratic expressions in the first and 2nd derivative of $\phi)$ and the ellipticity constant for $P^{\prime}(0)$ which is one. It must be noted here that we have suppressed the fact that such a constant also depends on $\alpha, n$, and $R$ so that $H$ does as well.

We must now give an upper bound for $\left\|P^{\prime \prime}(u)\right\|$ on $S(0, r)$.
We have

$$
\begin{aligned}
P^{\prime \prime}(v)(h, h) & =(\Delta v+\Delta \phi) h_{i} h_{i}-\left(v_{i j}+\phi_{i j}\right) h_{i} h_{j} \\
& +2\left(v_{i}+\phi_{i}\right) h_{i} \Delta h-2\left(v_{i}+\phi_{i}\right) h_{j} h_{i j},
\end{aligned}
$$

so that

$$
\left|P^{\prime \prime}(v)(h, h)\right|_{\alpha} \leq\left[4|v|_{2+\alpha}+|\Delta \phi|_{\alpha}+\max _{i, j}\left|\phi_{i j}\right|_{\alpha}+2 \max _{i}\left|\phi_{i}\right|_{\alpha}\right]\left(|h|_{2+\alpha}\right)^{2}
$$

1
and

$$
\left\|P^{\prime \prime}(v)\right\| \leq 4|v|_{2+a}+\lambda(\phi)
$$

where

$$
\lambda(\phi)=|\Delta \phi|_{\alpha}+\max _{i, j}\left|\phi_{i j}\right|_{\alpha}+\underset{\mathbf{i}}{\max }\left|\phi_{\mathbf{i}}\right|_{\alpha}
$$

Thus, for $v$ in $s(0, r)$, we have

$$
\left\|P^{\prime \prime}(v)\right\| \leq 4 r+\lambda(\phi) .
$$

We may then take

$$
K=\max \left\{\frac{1}{B r}, 4 r+\lambda(\phi)\right\}
$$

in the expression

$$
1 / 2 B^{2} K
$$

for $H$.
Some light can be shed on the expression for $K$ by observing that for $r \leq r_{0}$, where

$$
r_{0}=\frac{-\lambda B+\sqrt{\lambda^{2} B^{2}+16 B}}{8 B}
$$

we have

$$
\frac{1}{B r} \geq 4 r+\lambda
$$

and :

$$
H=\frac{r}{2 B},
$$

and for $r \geq r_{0}$

$$
\frac{1}{\mathrm{Br}} \leq 4 r+\lambda,
$$

and

$$
H=\frac{1}{2 B^{2}(4 r+\lambda)},
$$

where $r_{0}$ is the unique value of $r$ for which

$$
\frac{1}{B r}=4 r+\lambda .
$$

We also note that $r_{0}$ is the optional choice for $r$ in the sense that $H$ has an absolute maximum for this choice of $r$, and that $H$ goes to zero as $r$ either goes to zero or to infinity.

The above theorem is stated so that the function $\phi$, defined on all of $R$, is somehow given to us. We could turn the situation around somewhat.

Suppose we are given a boundary function $\phi_{0}$ which can be extended to a function defined on $R$ which is an element of $C_{2+\alpha}(R)$ and which is such that the $C_{\alpha}$ norms of the first and second derivatives of at least one of the "extensions" are bounded by some specified constant. He call such extensions admissable. Note then, that the numbers $B$ and $\lambda$ used above are determined by the constant used to define "admissable". Further, suppose we choose the value $r_{0}$ of $r$ determined by this constant. The above theorem then guarantees existence of a solution of the boundary value problem (in $\left.C_{2+\alpha}(R)\right)$ if $\phi_{0}$ has an admissable extension $\phi$ which is such that:

$$
|M(\phi)|_{\alpha}<H=r_{0} / 2 B,
$$

that is, the problem has a solution if there exists a function which is sufficiently close to being a solution in the above sense.

Remark. In connection with the above we mention a recent result of Jenkins an Serrin [5]. They show that for a $C_{2}$ region there is a number $B$, depending on the region and the uniform norms of the first and second derivatives of the boundary data, such that if the oscillation of the boundary data is less than $\boldsymbol{B}$ there exists a solution. The number $B$ is shown to depend in a precise way on the mean curvature of the boundary (mean curvature everywhere non-negative implies $b=\infty$ ), and their result is also far better than ours in other ways. We would, however, like to point out that in the case of a region whose boundary does not have an everywhere non-negative mean curvature ( $B$ finite) our result gives a sufficient condition for existence that involves relations between norms of the first and second derivatives of an extension of the boundary function, without explicit restriction on its oscillation on the boundary.

We conclude by noting that the theorems we have presented are constructive in that by using the Kantorovich theorem they produce a sequence of functions converging to a solution, the elements of this sequence being solutions of linear problems.

## References

1
[1] A. Friedman, Partial Differential Equation of Parabolic Type, Prentice Ha11, Inc., Englewood Cliffs, N.J., 1964.
[2] L. Bers, F. John, and M. Schecter, Partial Differential Equations, Wiley (Interscience), 1964.
[3] J. Dieudonne, Foundations of Modern Analysis, Academic Press, New York, 1960.
[4] L. V. Kantorovich, and G. P. Akilov, Functional Analysis in Normed Spaces (translated by D. E. Brown), The MacMillan Co., New York, 1964.
[5] H. Jenkins, and J. Serrin, The Dirichlet problem for the minimal surface equation in higher dimensions, Journal für die Reine und Angervandte Mathematik, 229 (1968), 170-187.


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