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## Layering methods for nonlinear partial differential equations of first order

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$\mathcal{N u m b a m}^{\prime}$

# LAYERING METHODS <br> FOR NONLINEAR <br> PARTIAL DIFFERENTIAL EQUATIONS OF FIRST ORDER* 

by Avron DOUGLIS

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## Introduction.

Global solutions of nonlinear first order partial differential equations have been an object of inquiry since 1950-51, when for the particular equation

$$
\begin{equation*}
u_{t}+u u_{x}=0 \tag{1}
\end{equation*}
$$

E. Hopf [16] and J. Cole [2] first succeeded in constructing "weak solutions" -solutions in a distribution sense- in the halfplane $t \geqslant 0$ with bounded, measurable initial values. Their methods, which were similar but independently arrived at, after some years were accomodated to other equations than (1) and additional approaches developed. At present, several ways are known to construct global weak solutions with prescribed initial values of fairly general types of scalar first order equations, including equations of HamiltonJacobi form

$$
\begin{equation*}
u_{t}+f(x, t, \operatorname{grad} u)=0 \tag{2}
\end{equation*}
$$

and quasilinear "conservation laws"

$$
\begin{equation*}
u_{t}=\sum_{i=1}^{n} \frac{\partial}{\partial x_{i}} \mathrm{~A}_{i}(x, t, u)=\mathrm{C}(x, t, u) \tag{3}
\end{equation*}
$$

Here, $x=\left(x_{1}, x_{2}, \ldots, x_{n}\right) \in \mathrm{E}^{n}, u(x, t)$ denotes a scalar function, grad $u=\left(u_{x_{1}}, u_{x_{2}}, \ldots, u_{x_{n}}\right), u_{x_{i}}=\partial u / \partial x_{i}, u_{t}=\partial u / \partial t$. A beginning has been made with scalar equations in two independant variables of the form $(f(u))_{t}+(g(u))_{x}=0$ (D.P. Ballou [1]), and weak solutions of certain hyperbolic systems of quasilinear equations also have been obtained (see P.D. Lax [28], J. Glimm [14], Glimm and Lax [15], J.A. Smoller [32], J.L. Johnson and J.A. Smoller [19], ( ${ }^{\mathbf{1}}$ ) and the bibliographies in these papers). Nevertheless, global theory today is incomplete, and often it is complex and delicate. Hence, it seems worthwhile to continue to explore the scalar case, to which this paper is devoted.
( ${ }^{1}$ ) See also J.A. Smoller and C.C. Conley [1-6] and the included references.

Our methods are based on the idea of approximating weak solutions by functions that are built up out of strict $C^{\mathbf{1}}$ solutions in appropriately narrow domains. This kind of construction seems to have been first suggested in principle by B.L. Rozhdestvenskii [31] (1961) and was first carried out by N.N. Kuznetsov [27] (1967). I hit upon it later and outlined its application to equations (2) and (3) in an A.M.S. lecture $\left({ }^{1}\right)$ in 1968, not knowing that Kuznetsov had already introduced it. In this paper, the original methods are greatly broadened and developed.

The principal estimates needed in this approach are, as Kuznetsov also stresses, suitable a priori inequalities for strict solutions - solutions of class $C^{1}-$ of the given partial differential equations. Besides these inequalities, a smoothing transformation $S$ of $C^{1}\left(\mathrm{E}^{n}\right)$ into itself is required under which the inequalities are preserved. Supposing we have $S$, we can describe the idea as follows. Let $u_{0}(x)$ denote the prescribed initial data. The first step is to find, in a suitably thin layer $0 \leqslant t \leqslant h$, a strict solution $u_{1}(x, t)$ satisfying the initial condition $u_{1}(x, 0)=\operatorname{Su}_{0}(x)$. The second step is to find in the layer $h \leqslant t \leqslant 2 h$ a strict solution $u_{2}(x, t)$ such that $u_{2}(x, h)=\operatorname{Su}_{1}(x, h)$. The third step is to find in the layer $2 h \leqslant t \leqslant 3 h$ a strict solution $u_{3}(x, t)$ such that $u_{3}(x, 2 h)=\operatorname{Su}_{2}(x, 2 h)$, and so forth. The sectionally continuous, sectionally smooth function

$$
\begin{aligned}
u(x, t) & =u_{1}(x, t) \quad \text { for } \quad 0 \leqslant t<h \\
& =u_{2}(x, t) \quad \text { for } \quad h \leqslant t<2 h \\
& =u_{3}(x, t) \quad \text { for } \quad 2 h \leqslant t<3 h
\end{aligned}
$$

we call a layered, or stratified, solution. It obeys inequalities depending on $u_{0}$ that derive from the estimates originally established for strict solutions. It is of course not a weak solution of the original problem, but under appropriate conditions a sequence of layered solutions will approximate a weak solution if the smoothing performed is made finer and finer with h accordingly tending toward zero.
$\left.{ }^{1}{ }^{1}\right)$ The existence of weak solutions of first order partial differential equations, an invited address before the Eastern Sectional Meeting of the American Mathematical Society at Johns Hopkins University, October 26, 1968.

A large class of smoothing operators, including arithmetic and other types of averaging as well as quite different kinds of transformations, exists with which layering works comparatively simply for equations of forms (2) and (3).

## CHAPTER 2

## GENERALIZED SOLUTIONS OF QUASILINEAR CONSERVATION Laws WITH BOUNDED, MEASURABLE INITIAL DATA

Global weak solutions of multi-dimensional quasilinear conservation laws were first constructed by E.D. Conway and J.A. Smoller [4], who confined their efforts to equations of the form

$$
u_{t}+\sum_{i=1}^{n} \frac{\partial}{\partial x_{i}} \mathrm{~A}_{i}(u)=0
$$

For equations

$$
\begin{equation*}
u_{t}+\sum_{i=1}^{n} \frac{\partial}{\partial x_{i}} \mathrm{~A}_{i}(x, t, u)=\mathrm{C}(x, t, u) \tag{A}
\end{equation*}
$$

that depend explicitly upon $x$ and $t$, similar results were then obtained by K. Kojima [20], extending Conway's and Smoller's finite difference methods, and by N.N. Kuznetsov [27], who used a layering procedure. Finite difference schemes in relation to the two equations $u_{t}+\nabla \cdot \mathrm{F}(u)=0$ and $u_{t}+\nabla \cdot \mathrm{F}(u)=\nu \Delta u$ have been further elucidated by D.B. Kotlow [21].

These writers required the initial data to be of bounded variation in a sense of Tonelli and Cesari, but recently new estimates in conjunction with the "viscosity" method have made it possible for S.N. Kruzhkov [25, 26] to dispense with all demands upon initial data other than boundedness and measurability. In the same papers, Kruzhkov also has presented a new definition of "generalized solution" under which generalized solutions of $(A)_{0}$ are weak solutions but, almost regardless of the nature of the $A_{i}$, depend uniquely and continuously (in $L^{1}$ ) upon their initial data $\left({ }^{1}\right)$. In the one-dimensional case, Kruzhkov's generalized solutions satisfy the "entropy" or analogous conditions that hitherto have served as the basis of uniqueness theorems.

Ideas closely related to Kruzhkov's were presented by E. Hopf [18]. For sectionally continuous solutions, B.K. Quinn [1-5] has provided an elegant alternative to Kruzhkov's methods.
${ }^{(1)}$ P.A. Andreyanov [1-2] recently has extended Kruzhkov's methods and ideas to problems with just locally bounded initial data.

Viscosity methods like Kruzhkov's are of interest in themselves, as Kruyhkov stresses, but procedures not involving parabolic equations doubtless also are worthwhile. Such a procedure, which leads relatively easily to Kruzhkov's existence theorem, is to approximate the given initial data by functions of locally bounded variation (in the sense of Tonelli and Cesari) and then to pass from solutions that have these data to the solution with the data desired by means of a limit process based on Kruzhkov's theorem on continuous dependence alluded to earlier. (Kuznetsov [27] used a method somewhat like this for equations of the form $\left.u_{t}+\sum_{i}\left(\mathrm{~A}_{i}(u)\right)_{x_{i}}=0\right)$. In this approach, incidentally, second instead of third order differentiability of the $\mathrm{A}_{\boldsymbol{i}}$ and $C$ suffices.

This chapter is devoted to a direct approach, by means of layering, to the problem of bounded, measurable initial data and also provides a pattern for our discussion of the case of initial data of bounded variation (Chapter 4) and our treatment of Hamilton-Jacobi equations (Chapter 5). The lack of a needed estimate (an appropriate generalization of Theorem 3) restricts present considerations to conservation laws of the form

$$
\begin{equation*}
u_{t}+\sum_{i=1}^{n} \frac{\partial}{\partial x_{i}} \mathrm{~A}_{i}(t, u)=\mathrm{C}(t, u) \tag{A}
\end{equation*}
$$

without explicit dependence upon $x$. Initial conditions

$$
\begin{equation*}
u(x, 0)=u_{0}(x) \tag{B}
\end{equation*}
$$

with bounded, measurable $u_{0}$ are imposed.

## 1. Notation, assumptions, aims.

For an arbitrary point $x=\left(x_{1}, \ldots, x_{n}\right)$ of real Euclidean space $\mathrm{E}^{n}$, we write

$$
|x|=\left[\sum_{i=1}^{n} x_{i}^{2}\right]^{1 / 2}
$$

If $v$ is a bounded function of $x$ on $E^{n}$, we write

$$
M_{v}=\sup _{x \in \mathbb{E}}|v(x)| ;
$$

if $v$ is a function of $x, t$ for $x \in \mathrm{E}^{n}$ and $t \geqslant 0$, and $v(\bullet, t)$ is bounded for fixed $t$, we write

$$
\mathrm{M}_{v}(t)=\sup _{x \in \mathrm{E}^{n}}|v(x, t)|
$$

Similar notation is used for vector functions.
Let $C\left(E^{n}\right)$ denote the class of functions that are bounded and continuous on $\mathrm{E}^{n}$, and for $v \in \mathrm{C}\left(\mathrm{E}^{n}\right)$ write

$$
|v|_{0}=\sup _{x \in \mathbf{E}^{n}}|v(x)|
$$

For a vector $f=\left(f_{1}, \ldots, f_{n}\right)$ with $f_{i} \subset \mathrm{C}\left(\mathrm{E}^{n}\right)(i=1, \ldots, n)$, we write

$$
|f|_{0}=\left|\left[\sum_{i=1}^{n} f_{i}^{2}\right]^{1 / 2}\right|_{0}
$$

For $\mathrm{T}>0, \mathrm{M}>0, a<b, h>0$, and positive integers $m$, define

$$
\begin{aligned}
\mathrm{Z}(\mathrm{~T}) & =\left\{(x, t): x \in \mathrm{E}^{n}, 0 \leqslant t<\mathrm{T}\right\}, \\
\mathrm{Z}(a, b) & =\left\{(x, t): x \in \mathrm{E}^{n}, a \leqslant t \leqslant b\right\}, \\
\mathrm{Z}_{m} & =\mathrm{Z}_{m, h}=\mathrm{Z}((m-1) h, m h), \\
\mathrm{Z}_{m}^{\prime} & =\left\{(x, t): x \in \mathrm{E}^{n},(m-1) h \leqslant t<m h\right\} ; \\
\mathrm{Z}(\mathrm{~T} ; \mathrm{M}) & =\left\{(x, t, u): x \in \mathrm{E}^{n}, 0 \leqslant t<\mathrm{T},|u| \leqslant \mathrm{M}\right\}, \\
\mathrm{Z}(a, b ; \mathrm{M}) & =\left\{(x, t, u): x \in \mathrm{E}^{n}, a \leqslant t \leqslant b,|u| \leqslant \mathrm{M}\right\}, \\
\widetilde{\mathrm{Z}}(\mathrm{~T} ; \mathrm{M}) & =\{(x, t): 0 \leqslant t \leqslant \mathrm{~T},|u| \leqslant \mathrm{M}\} .
\end{aligned}
$$

If $Z$ represents any domain, by $C^{k}(Z)$ we shall mean the class of functions that are bounded and continuous and have bounded, continuous partial derivatives on Z of orders up to $k\left({ }^{1}\right)$. By $\mathrm{C}_{0}^{k}(\mathrm{Z})$ we mean the subclass of $C^{k}(Z)$ consisting of the functions that have compact support in Z . Here, $k$ denotes any nonnegative integer. If
(1) The conditions of boundedness will be dropped in Section 3, Chapter 3, and thereafter.
$f \in \mathrm{C}^{k}(\mathrm{Z})$, we write $|f|_{\mathrm{C}^{k}(\mathrm{Z})}$ to mean, as is usual, the sum of the least upper bounds in Z of the absolute values of $f$ and of its partial derivatives of orders up to $k$. By $L^{1}(Z)$ we mean the class of Lebesgue summable functions on $Z$, and by $|f|_{L^{1}(Z)}$ the usual norm of $f$ in that space.

For all positive constants $\mathrm{M}, \mathrm{T}$, we assume

$$
\mathrm{C} \in \mathrm{C}^{1}(\widetilde{\mathrm{Z}}(\mathrm{~T} ; \mathrm{M})), \mathrm{A}_{i} \in \mathrm{C}^{2}(\widetilde{\mathrm{Z}}(\mathrm{~T} ; \mathrm{M})), i=1,2, \ldots, n
$$

We also assume $u_{0}$ to be bounded and measurable in $\mathrm{E}^{n}$. Following Kruzhkov [26], we shall say that a bounded, measurable function $u$ on the set $Z(T)$ is a generalized solution of $(A)_{0}$, (B) if for all real constants $k$ and all nonnegative functions $f \in \mathrm{C}_{0}^{1}(\mathrm{Z}(\mathrm{T}))$ with $f(x, 0)=0$, we have
$\iint_{\mathrm{Z}(\mathrm{T})}\left\{|u-k| f_{t}+\operatorname{sign}(u-k) \cdot \sum_{i=1}^{n}\left[\mathrm{~A}_{i}(x, t, u)-\mathrm{A}_{i}(x, t, k)\right]\right.$

$$
\left.+\operatorname{sign}(u-k) \cdot\left[\mathrm{C}(x, t, u)-\sum_{i=1}^{n} \mathrm{~A}_{i, x_{i}}(x, t, k)\right] f\right\} d x d t
$$

$$
\begin{equation*}
\geqslant 0 \tag{3a}
\end{equation*}
$$

here, $f_{t}=\partial f / \partial t, f_{x_{i}}=\partial f / \partial x_{i}$. In addition, it is required that a subset $\mathcal{E}$ of $[0, \mathrm{~T}]$ of measure zero exist such that for every ball

$$
\begin{gather*}
\mathrm{S}(r)=\left\{x \in \mathrm{E}^{n}:|x| \leqslant r\right\} \\
\lim _{t \rightarrow 0} \int_{\mathrm{S}(r)}\left|u(t, x)-u_{0}(x)\right| d x=0 \tag{3b}
\end{gather*}
$$

where during the limit process $0<t<\mathrm{T}$ and $t \notin \mathbb{E}$. As Kruzhkov proves, generalized solutions are weak solutions. Furthermore, all generalized solutions of an equation for which the $A_{i}$ and $C$ are of class $\mathrm{C}^{1}$, and the first derivatives of the $\mathrm{A}_{i}$ with respect to $x$ and $t$ are Lipschitz-continuous, for $(x, t, u) \in \mathrm{Z}(\mathrm{T}, \mathrm{M})$ with arbitrary M , depend uniquely and continuously (in $\mathrm{L}^{1}$ ) upon their initial data. (See also B.K. Quinn [1-5]).

For our purposes, we add to the previous assumptions the following growth condition upon $\mathrm{C}:$ for $x \in E^{n}$ and $t \geqslant 0$

$$
\sup _{|v| \leqslant w}|\mathrm{C}(t, v)| \leqslant \mathrm{E}(w)
$$

where E is a nonnegative, nondecreasing function defined for $w>0$ such that

$$
\int_{a}^{\infty} \frac{d w}{\mathrm{E}(w)}=\infty \text { for } a>0
$$

Our principal aim in this chapter is to prove the following result :

Thiorim A. - Under the assumptions stated, a generalized solution of (A) and (B) exists.

To employ stratified solutions in constructing the solution demanded, we must first establish certain properties of strict solutions and also suitable smoothing operators. The latter will be discussed in detail in the next chapter, from which we shall borrow as required. We now turn to the former.

## 2. Strict solutions of equation ( $\mathrm{A}_{0}$ ).

By a strict solution of a partial differential equation we mean a continuously differentiable solution in the ordinary sense. In this section, we shall discuss strict solutions of first order conservation laws of the form

$$
\begin{equation*}
u_{t}+\sum_{i=1}^{n} \frac{\partial}{\partial x_{i}} \mathrm{~A}_{i}(x, t, u)=\mathrm{C}(x, t, u) \tag{A}
\end{equation*}
$$

where $C \in C^{1}(Z(T ; M)), \quad A_{i} \in C^{2}(Z(T ; M))$ for $i=1, \ldots, n$. We suppose a generalized form of the previous growth condition to hold : for $x \in \mathrm{E}^{n}$ and $t \geqslant 0$,

$$
\sup _{|v| \leqslant w}\left|\sum_{i=1}^{n} \mathrm{~A}_{i, x_{i}}(x, t, v)-\mathrm{C}(x, t, v)\right| \leqslant \mathrm{E}(w)
$$

where $\mathrm{A}_{i, x_{i}}=\frac{\partial \mathrm{A}_{i}}{\partial x_{i}}$ and E , as before, is a nonnegative, nondecreasing function defined for $w>0$ such that $\int_{a}^{\infty} \frac{d w}{\mathrm{E}(w)}=\infty$ for $a>0$. An initial condition of the form

$$
\begin{equation*}
u\left(x, t_{0}\right)=w(x) \tag{1}
\end{equation*}
$$

is assumed, where $w$ is an arbitrary bounded, continuously differentiable function with bounded gradient :

$$
|w|_{0} \leqslant w_{0},|\operatorname{grad} w|_{0} \leqslant w_{1}
$$

$w_{0}$ and $w_{1}$ being constants. Standard theory will provide positive constants A and $a$ independent of $w_{1}$ such that $u$ exists in the layer

$$
\mathrm{Z}\left(t_{0}, t_{0}+1 /\left(\mathrm{A} w_{1}+a\right)\right)
$$

and will furnish an upper bound pertaining to this layer for $\mathrm{M}_{u}(t)$. Well known methods will also enable us to show that $u$-a strict solution of $(A)_{0}$ and (1)- is a generalized solution in an appropriate sense. Finally, we shall prove $u$ to be mean continuous (in the $\mathrm{L}_{1}$ sense) with respect to $x$, uniformly with respect to $w_{1}$. This information is the main content of the following three theorems.

Thiorlim 1. A strict solution $u$ of $(\mathrm{A})_{0}$, (1) exists in the layer $\mathrm{Z}\left(t_{0}, t_{0}+h\right), h=1 /\left(\mathrm{A} w_{1}+a\right)$, where the constants A and a may depend upon $w_{0}$, the bounds assumed for the first derivatives of C , and those for the second derivatives of the $\mathrm{A}_{i}$. This solution $u$ is subject to a bound described as follows. Define the function

$$
\phi(t) \equiv \phi\left(t ; t_{0}, w_{0}\right)
$$

by the condition

$$
\int_{w_{0}}^{\phi(t)} \frac{d v}{\mathrm{E}(v)}=t-t_{0}
$$

(Under our growth assumption, $\phi$ exists and is an increasing, absolutely continuous function for $t \geqslant t_{0}$ ). Then

$$
\begin{equation*}
\mathrm{M}_{u}(t) \leqslant \phi\left(t ; t_{0}, w_{0}\right) \text { for } t_{0} \leqslant t \leqslant t_{0}+h \tag{2}
\end{equation*}
$$

Theorlm 2. - Let $u$ be a strict solution of $(\mathrm{A})_{0}$ in $\mathrm{Z}\left(t_{0}, t_{1}\right)$, where $t_{0}<t_{1}$, and suppose $f \in \mathrm{C}^{1}\left(\mathrm{Z}\left(t_{0}, t_{1}\right)\right)$ to vanish identically for $|x|>r, r$ being any positive constant. Then for any constant $k$

$$
\begin{aligned}
\int_{\mathrm{z}_{\left(t_{0}, t_{1}\right)}}\left\{|u-k| f_{t}\right. & +\operatorname{sign}(u-k) \cdot \sum_{i=1}^{n}\left[\mathrm{~A}_{i}(x, t, u)\right. \\
& \left.-\mathrm{A}_{i}(x, t, k)\right] f_{x_{i}}+\operatorname{sign}(u-k) \cdot[\mathrm{C}(x, t, u) \\
& \left.\left.-\sum_{i=1}^{n} \mathrm{~A}_{i, x_{i}}(x, t, k)\right] f\right\} d x d t \\
& =\left.\int_{\mathrm{E}^{n}}|u-k| f d x\right|_{t=t_{0}} ^{t_{1}}
\end{aligned}
$$

Theorlm 3( ${ }^{1}$ ). - Let $u$ be a strict solution of (A) in a layer $\mathrm{Z}\left(t_{0}, t_{1}\right)\left(t_{0}<t_{1}\right)$, in. which $|u(x, t)| \leqslant \mathrm{M}$. Let

$$
\mathrm{N}=\sup _{Z\left(t_{0}, t_{1} ; \mathrm{M}\right)}\left[\sum_{r=1}^{n}\left(\partial \mathrm{~A}_{r} / \partial u\right)^{2}\right]^{1 / 2}
$$

With $\xi \in \mathrm{E}^{n}$ and $\mathrm{T} \geqslant t_{1}$, denote by $\mathrm{D}(t ; \xi, \mathrm{T})$ the horizontal disk

$$
\mathrm{D}(t ; \xi, \mathrm{T})=\{x:|x-\xi|<\mathrm{N}(\mathrm{~T}-t)\}
$$

cut out of the plane with ordinate $t$ by the back horizontal cone with vertex $(\xi, \mathrm{T})$ and "slope" N . Then for any $z \in \mathrm{E}^{n}$ and for $t_{0} \leqslant t \leqslant t_{1}$, we have :

$$
\begin{aligned}
\int_{\mathrm{D}(t ; \xi, \mathrm{T})} & |u(x+z, t)-u(x, t)| d x \\
& \leqslant \exp \left[\mathrm{C}_{1}\left(t-t_{0}\right)\right] \cdot \int_{\mathrm{D}\left(t_{0} ; \xi, \mathrm{T}\right)}\left|u\left(x+z, t_{0}\right)-u\left(x, t_{0}\right)\right| d x
\end{aligned}
$$

Here, $C_{1}=\sup _{Z\left(t_{0}, t_{1} ; M\right)}\left|C_{u}\right|$.
Proof of Theorem 1. - For strict solutions, equation $(A)_{0}$ is equivalent to
(1) Kuznetsov [27] gives such a result in the case in which $\mathrm{C}=0$.

$$
\begin{equation*}
u_{t}+\sum_{i=1}^{n} a_{i}(x, t, u) u_{x_{i}}=c(x, t, u) \tag{3}
\end{equation*}
$$

where $a_{i}=\partial \mathrm{A}_{i} / \partial u$ and $c=\mathrm{C}-\sum_{i=1}^{n}\left(\partial \mathrm{~A}_{i} / \partial x_{i}\right)$. The coefficients in (3) satisfy the following conditions :
i) For all positive constants $\mathrm{M}, \mathrm{T}$,

$$
c, a_{i} \in \mathrm{C}^{1}(\mathrm{Z}(\mathrm{~T} ; \mathrm{M})), i=1, \ldots, n ;
$$

ii) in terms of the function $E$ previously defined, $c$ satisfies the condition

$$
\sup _{|v| \leqslant w}|c(x, t, v)| \leqslant \mathrm{E}(w)
$$

for $x \in \mathrm{E}^{n}$ and $t \geqslant 0$.
We use the method of characteristics as follows. For an arbitrary point $\xi$ of the initial plane $t_{0}=0$, let the $n+1$ functions

$$
\begin{equation*}
\mathrm{U}(t ; \xi), x_{j}(t ; \xi), j=1,2, \ldots, n \tag{4}
\end{equation*}
$$

denote solutions of the differential equations

$$
\begin{equation*}
d \mathrm{U} / d t=c(x, t, \mathrm{U}), d x_{j} / d t=a_{j}(x, t, \mathrm{U}), \quad j=1,2, \ldots, n \tag{5}
\end{equation*}
$$

satisfying the initial conditions

$$
\begin{equation*}
\mathrm{U}\left(t_{0} ; \xi\right)=w(\xi), x_{j}\left(t_{0} ; \xi\right)=\xi_{j}, \quad j=1,2, \ldots, n \tag{6}
\end{equation*}
$$

We shall find an a priori bound for U in any interval

$$
\mathrm{I}: t_{0} \leqslant t \leqslant t_{1} \quad\left(t_{0}<t_{1}\right)
$$

within which, for fixed $\xi$, the solution (4) exists and is of class $C^{1}$. By (5),

$$
\mathrm{U}(t ; \xi)=w(\xi)+\int_{0}^{t} c(x(s ; \xi), s, \mathrm{U}(s ; \xi)) d x \quad \text { for } t \in \mathrm{I}
$$

Assumption (ii) implies that

$$
\begin{aligned}
\mathrm{U}(t ; \xi) \leqslant|w|_{0} & +\int_{0}^{t} \mathrm{E}(|\mathrm{U}(s ; \xi)|) d s \\
& \leqslant w_{0}+\int_{0}^{t} \mathrm{E}(|\mathrm{U}(s ; \xi)|) d s
\end{aligned}
$$

where, as required, $w_{0}$ is a constant such that

$$
|w|_{0} \leqslant w_{0} .
$$

Since $\phi$ satisfies the condition $\phi^{\prime}(t)=\mathrm{E}(\phi(t))$ almost everywhere and $\phi\left(t_{0}\right)=w_{0}, \phi$ satisfies the condition

$$
\phi(t)=w_{0}+\int_{t_{0}}^{t} \mathrm{E}(\phi(s)) d s
$$

and reasoning such as is used for Gronwall's inequality shows that

$$
\begin{equation*}
|\mathrm{U}(t ; \xi)| \leqslant \phi(t) \quad \text { for } \quad t_{0} \leqslant t \leqslant t_{1} \tag{7}
\end{equation*}
$$

Therefore, in particular, $\left|\mathrm{U}\left(t_{1}, \xi\right)\right| \leqslant \phi\left(t_{1}\right)$, and standard theorems show that the values $\tau$ such that the solution of (5), (6) exists for $t_{0} \leqslant t \leqslant \tau$ have no finite least upper bound. Hence, this solution (4) will exist for $t \geqslant t_{0}$. By well known theorems, this solution also is of class $\mathrm{C}^{1}$ with respect to $t, \xi$ for $\xi \in \mathrm{E}^{n}, t \geqslant t_{0}$.

For fixed $t$, consider the transformation

$$
\begin{equation*}
x=x(t ; \xi) . \tag{8}
\end{equation*}
$$

We can invert (8), say as

$$
\begin{equation*}
\xi=\xi(x, t) \tag{9}
\end{equation*}
$$

in any layer

$$
\begin{equation*}
x \in \mathrm{E}^{n}, \quad t_{0} \leqslant t \leqslant t_{0}+\delta \tag{10}
\end{equation*}
$$

in which the functional determinant

$$
\mathrm{J}(t) \equiv \mathrm{J}(t ; \xi)=\partial(x) / \partial(\xi)
$$

exceeds a positive constant. In such a layer, Cauchy's theory says that the function

$$
u(x, t) \equiv \mathrm{U}(t ; \xi(x, t))
$$

in a solution of (3) and the only solution that satisfies the initial conditions (1). Furthermore, the functions

$$
p_{j}(t ; \xi) \equiv u_{x_{j}}(x(t ; \xi), t), \quad j=1,2, \ldots, n
$$

which are the derivatives of $u$ regarded as functions of $t, \boldsymbol{\xi}$, satisfy the differential equations

$$
\begin{equation*}
d p_{j} / d t=-\sum_{i=1}^{n}\left(a_{i, x_{j}} p_{i}+a_{i, u} p_{i} p_{j}\right)+c_{x_{j}}+c_{u} p_{j} \tag{11}
\end{equation*}
$$

These show that positive constants $\alpha$ and $\beta$ depending on the first derivatives of the $a_{i}$ and $c$ exist such that

$$
\left|d p_{j} / d t\right| \leqslant \alpha(p(t)+\beta)^{2}
$$

where

$$
p(t)=p(t ; \xi)=\max _{i}\left|p_{i}(t ; \xi)\right|
$$

By reasoning similar to that which led to (7) we have

$$
p(t ; \xi) \leqslant q(t)
$$

where $q(t)$ is the function that satisfies the differential equation $d q / d t=\alpha(q+\beta)^{2}$ and the initial condition $q\left(t_{0}\right)=w_{1} \geqslant p\left(t_{0}\right)$. Consequently,

$$
p(t ; \xi) \leqslant \frac{w_{1}+\alpha \beta\left(w_{1}+\beta\right)\left(t-t_{0}\right)}{1-\alpha\left(w_{1}+\beta\right)\left(t-t_{0}\right)}
$$

the right hand side in fact being $q(t)$. This shows, in particular, that if the transformation $x=x(t ; \xi)$ can be inverted as $\xi=\xi(x, t)$ in a layer

$$
\begin{equation*}
\xi \in \mathrm{E}^{n}, \quad t_{0} \leqslant t \leqslant t_{0}+\delta \tag{12}
\end{equation*}
$$

with

$$
\delta \leqslant 1 /\left[2 \alpha\left(w_{1}+\beta\right)\right] \equiv h
$$

then in that layer

$$
\begin{equation*}
p(t) \leqslant 2 w_{1}+\beta \tag{13}
\end{equation*}
$$

We shall now exhibit a positive lower bound for $\mathrm{J}(t) \equiv \mathrm{J}(t ; \xi)$ in this layer. It is well known -and easily proved from (5) and (11)that

$$
\begin{equation*}
d \mathrm{~J} / d t=\sum_{i=1}^{n}\left(a_{i, x_{i}}+a_{i, u} p_{i}\right) \mathrm{J} . \tag{14}
\end{equation*}
$$

It follows from this that

$$
\begin{equation*}
|d \mathrm{~J} / d t| \leqslant(\alpha p+\gamma) \mathrm{J} \quad \text { for } \quad t_{0} \leqslant t \leqslant t_{0}+\delta \tag{15}
\end{equation*}
$$

with the same $\alpha$ as before and $\gamma$ a bound for $\left|\sum_{i} a_{i, x_{i}}\right|$. In view of (6), we have $\mathrm{J}\left(t_{0}\right)=1$, and the homogeneous differential equation (14) implies that $\mathrm{J}>0$. By (15) and (13),

$$
\begin{aligned}
\log \mathrm{J} \geqslant & -\left[\alpha\left(2 w_{1}+\beta\right)+\gamma\right] h \\
& =-\frac{\alpha\left(2 w_{1}+\beta\right)+\gamma}{2 \alpha\left(w_{1}+\beta\right)} \\
& =-1-\frac{\gamma-\alpha \beta}{2 \alpha w_{1}+2 \alpha \beta} \\
& \geqslant-j_{0}
\end{aligned}
$$

where

$$
j_{0}=\left\{\begin{array}{l}
1 \quad \text { if } \quad \gamma \leqslant \alpha \beta \\
\frac{\gamma+\alpha \beta}{2 \alpha \beta} \quad \text { if } \quad \gamma>\alpha \beta
\end{array}\right.
$$

Thus,

$$
\begin{equation*}
\mathrm{J}(t) \geqslant e^{-j_{0}} \quad \text { for } \quad t_{0} \leqslant t \leqslant t_{0}+\delta \tag{16}
\end{equation*}
$$

But $\delta$ is any positive number not exceeding $h$ such that the equations $x=x(t ; \xi)$ can be solved for $\xi$ in the layer (12). Hence, in view of this estimate for the Jacobian of the transformation, the least upper bound of such $\delta$ must be $h$. This implies that the inversion (9) can be performed in the layer $\mathrm{Z}\left(t_{0}, t_{0}+h\right)$, and Cauchy's theory now guarantees that $\mathrm{U}(t ; \xi(x, t))$ is a solution of (1), (3) in this layer. In view of (7), the absolute value of this solution is not greater than $\phi\left(t ; t_{0}, w_{0}\right)$, as the theorem states. Thus, this theorem is wholly proved.

Proof of Theorem 2. - Since $k$ is a constant, we have

$$
\begin{align*}
(u-k)_{t} & +\sum_{i=1}^{n}\left(\mathrm{~A}_{i}(x, t, u)-\mathrm{A}_{i}(x, t, k)\right)_{x_{i}} \\
& =\mathrm{C}(x, t, u)-\sum_{i=1}^{n} \mathrm{~A}_{i, x_{i}}(x, t, k) \equiv \mathrm{C}^{*} \tag{17}
\end{align*}
$$

Let $S$ denote any maximal subdomain of the cylinder

$$
Z^{\prime}=\left\{(x, t):|x|<r, t_{0}<t<t_{1}\right\}
$$

within which $u \neq k$; by continuity, $u-k$ has constant $\operatorname{sign}$ in S . Letting $S^{\prime}$ denote the part of the boundary of $S$ on which $t \neq t_{0}$, $t \neq t_{1}$, and $|x|<r$, we have by maximality

$$
u-k=0 \quad \text { on } \quad S^{\prime}
$$

We intend to show that

$$
\begin{align*}
\int_{\mathrm{S}}\left\{(u-k) f_{t}\right. & +\sum_{i}\left(\mathrm{~A}_{i}(x, t, u)-\mathrm{A}_{i}(x, t, k)\right) f_{x_{i}} \\
& \left.+\mathrm{C}^{*} f\right\} d x d t=\int_{\mathrm{E}^{n}} f(u-k) d x \left\lvert\, \begin{array}{l}
\mathrm{s}_{1}^{\prime \prime} \\
\mathrm{s}_{0}^{\prime \prime}
\end{array}\right. \tag{*}
\end{align*}
$$

$S_{0}^{\prime \prime}$ denoting the part of the boundary of $S$ on which $t=t_{0}$ and $S_{1}^{\prime \prime}$ the part on which $t=t_{1}$. (Either or both of these parts may be empty). Since sign $(u-k)$ is constant in $S$, when $\left(^{*}\right)$ is justified it will imply

$$
\begin{aligned}
\int_{\mathrm{S}}\left\{f_{t}|u-k|+\operatorname{sign}(u-k) \cdot\right. & {\left[f_{x_{i}}\left(\mathrm{~A}_{i}(x, t, u)-\mathrm{A}_{i}(x, t, k)\right)\right.} \\
& \left.\left.+\mathrm{C}^{*} f\right]\right\} d x d t=\left.\int_{\mathrm{E}^{n}} f|u-k| d x\right|_{\mathrm{s}_{0}^{\prime \prime}} ^{\mathrm{s}_{1}^{\prime \prime}}
\end{aligned}
$$

Then summing the relations of this form for all domains S as described will give

$$
\begin{aligned}
& \int_{\mathrm{Z}\left(t_{0}, t_{1}\right)}\left\{f_{t}|u-k|+\operatorname{sign}(u-k) \cdot\left[\left(\mathrm{A}_{i}(x, t, u)-\mathrm{A}_{i}(x, t, k)\right) f_{x_{i}}\right.\right. \\
&\left.\left.+\mathrm{C}^{*} f\right]\right\} d x d t=\left.\int_{\mathrm{E}^{n}} f|u-k| d x\right|_{t=t_{0}} ^{t_{1}}
\end{aligned}
$$

as claimed, since $f$ is zero for $|x| \geqslant r$.

We have still to justify (*). Let $S_{\nu}$ be an open subset of $S$ with piecewise smooth boundary $\dot{\mathrm{S}}_{\nu}, \nu=1,2, \ldots$, such that $\mathrm{S}_{\nu} \cup \dot{\mathrm{S}}_{\nu} \subset \mathrm{S}_{\nu+1}$ and $\bigcup_{\nu=1}^{\infty} S_{\nu}=S$. We require in addition for any point P of $\dot{\mathrm{S}}_{\nu}$

$$
\operatorname{dist}\left(\mathrm{P}, \mathrm{~S}^{\prime}\right) \leqslant \alpha d_{\nu}, d_{\nu}=\inf _{\mathrm{Q} \in \dot{\mathrm{~S}}_{\nu}} \operatorname{dist}\left(\mathrm{Q}, \mathrm{~S}^{\prime}\right)
$$

where $\alpha$ is a constant greater than 1 , and for an arbitrary set V

$$
\operatorname{dist}(P, V)=\inf _{R \in V}|P-R|
$$

( $\mathrm{S}_{\nu}$ might be, for instance, a union of open cubes of edge $2^{-\nu-1} n^{-1 / 2}$, covering the part of $S$ that is removed from $S^{\prime}$ by at least the distance $2^{-\nu}$, with $\alpha=2$ ). Set

$$
\mathrm{S}^{*}=\mathrm{Z}^{\prime} \backslash \mathrm{S}, \mathrm{~S}_{\nu}^{*}=\mathrm{Z}^{\prime} \backslash \mathrm{S}_{\nu}
$$

and define the functions

$$
\begin{array}{rlrlrl}
\tilde{u} & =u \quad \text { in } & \mathrm{S}, \quad \tilde{u}_{\nu} & =u \quad \text { in } \quad \mathrm{S}_{\nu} \\
& =k & \text { in } & \mathrm{S}^{*}, & & =k
\end{array} \quad \text { in } \quad \mathrm{S}_{v}^{*} .
$$

By our previous remarks, $\tilde{u}$ is Lipschitz continuous in $Z^{\prime}$, say with Lipschitz constant $\gamma$.

For $\nu=1,2, \ldots$, let $\zeta_{\nu}(x, t)$ be a function in $\mathrm{C}_{0}^{1}\left(\mathrm{E}^{n+1}\right)$ such that $\zeta_{\nu} \geqslant 0, \iint \zeta_{\nu} d x d t=1$,

$$
\zeta_{\nu}(x, t)=0 \quad \text { for } \quad|x|^{2}+t^{2}>\left(d_{\nu} / 2\right)^{2}
$$

and

$$
\left|\mathrm{D} \zeta_{\nu}(x, t)\right| \leqslant \beta / d_{\nu}
$$

where $\beta$ is a suitable constant, $D$ represents differentiation with respect to $t$ or a component of $x$, and the integration with respect to $x$ and $t$ is over $\mathrm{E}^{n+1}$. The continuously differentiable functions

$$
u_{\nu}=\zeta_{\nu} * \tilde{u}_{\nu}=\iint \zeta_{\nu}(x-y, t-s) \tilde{u}_{\nu}(y, s) d y d s
$$

then have the following properties :
a) $u_{\nu}=k$ at all points nearer $\mathrm{S}^{\prime}$ than the distance $d_{\nu} / 2$;
b) $\lim _{\nu \rightarrow \infty} u_{\nu}=u$ and $\lim _{\nu \rightarrow \infty} D u_{\nu}=\mathrm{Du}$ in S ;
c) the derivatives $\mathrm{Du}_{\nu}$ are bounded uniformly with respect to $\nu$. Properties (a) and (b) are clear. To justify (c), we write

$$
u_{\nu}=\zeta_{\nu} * \widetilde{u}+\zeta_{\nu} *\left(\tilde{u}_{\nu}-\widetilde{u}\right) \equiv \mathrm{Y}+\mathrm{Z}
$$

The first derivatives of Y are uniformly bounded, because $\tilde{u}$ is Lipschitz-continuous. Since $\widetilde{u}_{\nu}=\tilde{u}$ in $S^{*} \cup S_{\nu}$, we have

$$
\mathrm{DZ}=\int_{\mathrm{S} \backslash \mathrm{~s}_{\nu}} \mathrm{D} \zeta_{\nu}(x-y, t-s)[k-u(y, s)] d y d s
$$

and by the assumed estimates

$$
|\mathrm{DZ}| \leqslant\left(\beta / d_{\nu}\right) \gamma \alpha d_{\nu} \int_{\mathrm{S} \backslash \mathrm{~s}_{\nu}} d y d s=\alpha \beta \gamma \mu\left(\mathrm{S} \backslash \mathrm{~S}_{\nu}\right)
$$

$\mu$ referring to Lebesgue measure in $m+1$ dimensions. Consequently, $\mathrm{DZ} \rightarrow 0$ as $\nu \rightarrow \infty$, and we conclude, in particular, that the derivatives $D u_{\nu}$ are uniformly bounded, as asserted.

For this reason, if $L u_{\nu}$ is the expression obtained by replacing $u$ by $u_{\nu}$ in the first member of (17), then $\mathrm{L} u_{\nu}$ is bounded uniformly with respect to $\nu$. In addition, by property (b), $\lim _{\nu \rightarrow 0} L u_{\nu}=C^{*}$ in $S$. Therefore, Lebesgue's theorem gives :

$$
\lim _{\nu \rightarrow \infty} \int_{\mathrm{S}} f \mathrm{~L} u_{\nu} d x d t=\int_{\mathrm{S}} f \mathrm{C}^{*} d x d t
$$

On the other hand, partial integrations in the usual way result in the relation

$$
\begin{aligned}
& \int_{\mathrm{S}} f \mathrm{Lu}_{\nu} d x d t=-\int_{\mathrm{S}}\left\{\left(u_{\nu}-k\right) f_{t}+\sum_{i} f_{x_{i}}\left[\mathrm{~A}_{i}\left(x, t, u_{\nu}\right)\right.\right. \\
&\left.\left.\quad-\mathrm{A}_{i}(x, t, k)\right]\right\} d x d t+\left.\int f\left(u_{\nu}-k\right) d x\right|_{\mathrm{s}_{0}^{\prime \prime}} ^{\mathrm{s}_{1}^{\prime \prime}}
\end{aligned}
$$

other boundary terms disappearing because of (a). Letting $\nu \rightarrow \infty$ in this condition results in (*), completing our proof of Theorem 2.

A different proof is given by Kruzhkov [25].
Proof of Theorem 3. - Our aim is to prove that each strict solution of (A) is mean continuous in $x$, uniformly with respect to $t$. Let $f \in \mathrm{C}^{1}\left(\mathrm{Z}\left(t_{0}, t_{1}\right)\right)$, and suppose $f$ to vanish for $|x|>r$, where $r$
is some constant. We multiply both sides of equation (A) by $f$, integrate over the layer $\mathrm{Z}\left(t_{0}, t_{1}\right)$, and then integrate by parts to obtain :

$$
\int_{\mathrm{Z}\left(t_{0}, t_{1}\right)}\left(f_{t} u+\sum_{i} f_{x_{i}} \mathrm{~A}_{i}+f \mathrm{C}\right) d x d t=\left.\int_{\mathrm{E}^{n}} f u d x\right|_{t=t_{0}} ^{t_{1}} .
$$

An analogous relation will also hold in which $f$ is replaced by $g=\mathrm{T}_{-s} f, \mathrm{~T}_{s}$ denoting the translation defined by the condition $\mathrm{T}_{s} f(x, t)=f(x+s, t)$ for $s \in \mathrm{E}^{n}$. Since, for instance,

$$
\int_{\mathrm{E}^{n}}\left(\mathrm{~T}_{-s} f\right) u d x=\int_{\mathrm{E}^{n}} f\left(\mathrm{~T}_{s} u\right) d x
$$

we can write this second relation as

$$
\int_{Z\left(t_{0}, t_{1}\right)}\left(f_{t} \mathrm{~T}_{s} u+\sum_{i} f_{x_{i}} \mathrm{~T}_{s} \mathrm{~A}_{i}+f \mathrm{~T}_{s} \mathrm{C}\right) d x d t=\left.\int_{\mathrm{E}^{n}} f \mathrm{~T}_{s} u d x\right|_{t=t_{0}} ^{t_{1}} .
$$

Subtracting it from the first relation, we obtain :

$$
\begin{aligned}
\int_{\mathrm{Z}\left(t_{0}, t_{1}\right)}\left\{f_{t}\left(\mathrm{~T}_{s} u-u\right)+\sum_{i} f_{x_{i}}\left(\mathrm{~T}_{s} \mathrm{~A}_{i}-\mathrm{A}_{i}\right)\right. & \left.+f\left(\mathrm{~T}_{s} \mathrm{C}-\mathrm{C}\right)\right\} d x d t \\
= & \left.\int_{\mathrm{E}^{n}} f\left(\mathrm{~T}_{s} u-u\right) d x\right|_{t=t_{0}} ^{t_{1}}
\end{aligned}
$$

The rest of the argument depends upon a proper choice of $f$. Let

$$
w(x)=\left\{\begin{array}{l}
\operatorname{sign}\left(u\left(x+s, t_{1}\right)-u\left(x, t_{1}\right)\right) \text { for }|x-\xi| \leqslant \mathrm{N}\left(\mathrm{~T}-t_{1}\right) \\
0 \quad \text { for }|x-\xi|>\mathrm{N}\left(\mathrm{~T}-t_{1}\right)
\end{array}\right.
$$

and let $w_{m}, m=1,2, \ldots$, be continuously differentiable functions on $\mathrm{E}^{n}$ such that

$$
\left|w_{m}\right| \leqslant 1, w_{m}=0 \text { for }|x-\xi|>\mathrm{N}\left(\mathrm{~T}-t_{1}\right)+1 / m,
$$

and $\lim _{m \rightarrow \infty} w_{m}=w$ almost everywhere. (For instance,

$$
w_{m}(x)=m^{n} \int_{\mathrm{E}^{n}} k\left(m\left(x-x^{\prime}\right)\right) w\left(x^{\prime}\right) d x^{\prime}
$$

with continuously differentiable kernel $k$ such that $k \geqslant 0, k(x)=0$ for $|x| \geqslant 1, \int k(x) d x=1$ ). Soon we shall show that we can choose $f$ so that

$$
\begin{aligned}
& f\left(x, t_{1}\right)=w_{m}(x),\left|f\left(x, t_{0}\right)\right| \leqslant e^{b\left(t_{1}-t_{0}\right)} \\
& f\left(x, t_{0}\right)=0 \text { for }|x-\xi|>\mathrm{N}\left(\mathrm{~T}-t_{0}\right)+1 / m
\end{aligned}
$$

and the left member of (18) vanishes, $b$ denoting a constant. Then it will be clear from (18) that

$$
\int_{t=t_{1}} w_{m}\left(\mathrm{~T}_{s} u-u\right) d x=\int_{t=t_{0}} f\left(\mathrm{~T}_{s} u-u\right) d x
$$

Since the right side of this equality is not greater in absolute value than

$$
\left.e^{b\left(t_{1}-t_{0}\right)} \int_{\mathrm{D}\left(t_{0} ; \xi, \mathrm{U}_{m}\right)}\left|\mathrm{T}_{s} u-u\right| d x\right|^{t=t_{0}},
$$

where $\mathrm{U}_{m}=\mathrm{T}+1 / m$, letting $m \rightarrow \infty$ gives the required result.
To justify our choice of $f$, note first that

$$
\begin{aligned}
\mathrm{T}_{s} \mathrm{C}(t, u)-\mathrm{C}(t, u) & =\mathrm{C}(t, u(x+s, t))-\mathrm{C}(t, u(x, t)) \\
& =(u(x+s, t)-u(x, t)) \mathrm{C}^{\prime}(x, t, s),
\end{aligned}
$$

where

$$
\mathrm{C}^{\prime}(x, t, s)=\int_{0}^{1} \mathrm{C}_{u}(t, u(x, t)+\theta(u(x+s, t)-u(x, t)) d \theta .
$$

More briefly,

$$
\mathrm{T}_{s} \mathrm{C}-\mathrm{C}=\left(\mathrm{T}_{s} u-u\right) \mathrm{C}^{\prime},
$$

and, analogously,

$$
\mathrm{T}_{s} \mathrm{~A}_{i}-\mathrm{A}_{i}=\left(\mathrm{T}_{s} u-u\right) \mathrm{A}_{i}^{\prime}, \quad i=1, \ldots, n,
$$

$\mathrm{C}^{\prime}$ and $\mathrm{A}_{i}^{\prime}$ all being continuously differentiable in the domain considered. Schwarz's inequality applied to the integral expressions for the $\mathrm{A}_{i}^{\prime}$ shows that

$$
\begin{equation*}
\sum_{i=1}^{n}\left(\mathrm{~A}_{i}^{\prime}\right)^{2} \leqslant \mathrm{~N}^{2} \tag{19}
\end{equation*}
$$

Substituting from the previous results, we can write the expression in curly brackets in (18) as

$$
\left(\mathrm{T}_{s} u-u\right) \mathrm{L} f
$$

where

$$
\mathrm{L} f=f_{t}+\sum_{i} \mathrm{~A}_{i}^{\prime} f_{x_{i}}+\mathrm{C}^{\prime} f
$$

and we determine $f$ by the condition $\mathrm{L} f=0$,

$$
f\left(x, t_{1}\right)=w_{m}(x)
$$

Condition (19) and the nature of the support of $w_{m}$ show that $f\left(x, t_{0}\right)$ vanishes as required. The inequality stated for $f\left(x, t_{0}\right)$, with $b$ a bound for $\left|C_{u}\right|$ in the layer considered, follows from the assumption that $\left|w_{m}\right| \leqslant 1$. Thus, all the conditions demanded of $f$ are met, and Theorem 3 is proved.

This theorem can also be handled by such means as Kruzhkov [25] employed to show that generalized solutions depend uniquely and continuously upon the initial data. But those methods are as involved for strict as for generalized solutions, simplifications apparently not occuring in the case with which we are concerned.

## 3. Stratified solutions of initial value problems for equation (A).

In constructing stratified solutions of this problem, we shall smooth by means of averaging operations. Let $k(x)$ be a function of class $C^{1}$ in the cube

$$
\mathrm{C}:\left|x_{i}\right| \leqslant 1, \quad i=1, \ldots, n
$$

Assume that $k \geqslant 0$, that $\int_{\mathrm{C}} k d x=1$, and that $k$ is an even function of each individual coordinate $x_{i}$ of $x=\left(x_{1}, \ldots, x_{n}\right)$ when the other coordinates are held fixed. Then if, say, $f \in \mathrm{C}\left(\mathrm{E}^{n}\right)$-recall that $\mathrm{C}\left(\mathrm{E}^{n}\right)$ consists of the bounded, continuous functions on $\mathrm{E}^{n}$-define for $\varepsilon>0$

$$
\mathrm{K}_{\epsilon} f(x)=\left(\mathrm{K}_{\epsilon} f\right)(x)=\int_{\mathrm{C}} k(\xi) f(x+\varepsilon \xi) d \xi
$$

This operator has the following properties :

$$
\begin{align*}
\left|\mathrm{K}_{\epsilon} f\right|_{0} & \leqslant|f|_{0},  \tag{1}\\
\int_{|x| \leqslant r}\left|\mathrm{~K}_{\epsilon} f\right| d x & \leqslant \int_{|x| \leqslant r+\epsilon \sqrt{ } n}|f| d x  \tag{2}\\
\left|\operatorname{grad} \mathrm{~K}_{\epsilon} f\right| & \leqslant k_{1} \varepsilon^{-1}|f|_{0} \tag{3}
\end{align*}
$$

where $k_{1}$ is a constant depending on the function $k(\cdot)$. Properties (1) and (2) are immediate, and property (3) follows from the obvious generalization to $n$ dimensions of the following computation. Suppose $x$ to be a scalar quantity (i.e., $n=1$ ). Then

$$
\mathrm{K}_{\epsilon} f(x)=\int_{-1}^{1} k(\xi) f(x+\varepsilon \xi) d \xi=\varepsilon^{-1} \int_{x-\epsilon}^{x+\epsilon} k((y-x) / \varepsilon) f(y) d y
$$

and
$(d / d x) \mathrm{K}_{\epsilon} f(x)=\varepsilon^{-1}[k(1) f(x+\varepsilon)-k(-1) f(x-\varepsilon)]$

$$
-\varepsilon^{-2} \int_{x-\epsilon}^{x+\epsilon} k^{\prime}((y-x) / \varepsilon) f(y) d y
$$

while the integral on the right is equal to

$$
\varepsilon^{-1} \int_{-1}^{1} k^{\prime}(\xi) f(x+\varepsilon \xi) d \xi
$$

Therefore,

$$
\left|(d / d x) K_{\epsilon} f\right|_{0} \leqslant\left[|k(1)|+|k(-1)|+\int_{-1}^{1}\left|k^{\prime}(\xi)\right| d \xi\right] \varepsilon^{-1}|f|_{0}
$$

justifying (3) in this case.
Four more properties of $K_{\epsilon}$ will be used. In the notation

$$
\bar{f}(x)=|f(x)|
$$

the fourth property is that

$$
\begin{equation*}
\left|K_{\epsilon} f(x)\right| \leqslant K_{\epsilon} \bar{f}(x) . \tag{4}
\end{equation*}
$$

The fifth, referring to the translation operator $\mathrm{T}_{z}$ defined for $z \in E^{n}$ by the rule $\mathrm{T}_{z} f(x)=f(x+z)$, is :

$$
\begin{equation*}
\mathrm{K}_{\epsilon} \mathrm{T}_{z}=\mathrm{T}_{z} \mathrm{~K}_{\epsilon} \tag{5}
\end{equation*}
$$

The sixth property is that for any constant $k$,

$$
\begin{equation*}
\mathrm{K}_{\epsilon} f+k=\mathrm{K}_{\epsilon}(f+k) \tag{6}
\end{equation*}
$$

These are all trivial. The seventh property pertains to the order of weak approximation of $f$ by $\mathrm{K}_{\epsilon} f$ : for all $g \in \mathrm{C}_{0}^{2}\left(\mathrm{E}^{n}\right)$,

$$
\begin{equation*}
\left|\int_{\mathbf{E}^{n}} g\left(\mathrm{~K}_{\epsilon} f-f\right) d x\right| \leqslant \mathrm{C}_{2}|g|_{\mathbb{C}^{2}\left(\mathbf{E}^{n}\right)}|f|_{0} a(\varepsilon) \tag{7}
\end{equation*}
$$

where $C_{2}$ and the function $a(\varepsilon)$ depend only upon $n$, and

$$
\lim _{\epsilon \rightarrow 0} \varepsilon^{-1} a(\varepsilon)=0
$$

This property will be justified in the next chapter.
Kuznetsov used averaging operators such as $K_{\epsilon}$, but required the kernel $k$ to belong to $\mathrm{C}_{0}^{1}\left(\mathrm{E}^{n}\right)$.

Let $\mathrm{T}>0$ and $0<h \leqslant 1 / 2 a$ with $\mathrm{T} / \mathrm{h}$ an integer. (The constant $a$ is that of Theorem 1 , Section 2 ). We shall construct a stratified solution of (A), (B) in $\mathrm{Z}(\mathrm{T})$ with layer thickness $h$. Smoothing will be performed with $K_{\epsilon}$, where

$$
\varepsilon=2 \mathrm{~A} k_{1} \phi\left(\mathrm{~T} ; 0,\left|u_{0}\right|_{0}\right) h
$$

The resulting stratified solution $v$ will be subject to the estimate

$$
\begin{equation*}
|\nu|_{\mathrm{C}^{0}(\mathrm{Z}(0, \mathrm{~T}))} \leqslant \phi\left(\mathrm{T} ; 0,\left|u_{0}\right|_{0}\right) \tag{8}
\end{equation*}
$$

For convenience, set $\phi(t)=\phi\left(t ; 0,\left|u_{0}\right|_{0}\right)$. Since

$$
\left|\operatorname{grad} K_{\epsilon} u_{0}\right|_{0} \leqslant k_{1} \varepsilon^{-1}\left|u_{0}\right|_{0}<k_{1} \varepsilon^{-1} \phi(\mathrm{~T})
$$

Theorem 1 of Section 2 shows that a strict solution $u_{1}$ of (A) satisfying the initial condition

$$
u_{1}(x, 0)=K_{\epsilon} u_{0}(x)
$$

exists in a layer of thickness

$$
\frac{1}{\mathrm{~A} k_{1} \varepsilon^{-1} \cdot \phi(\mathrm{~T})+a}=\frac{2 h}{1+2 a h}
$$

since $h \leqslant 1 / 2 a$, this thickness is $\geqslant h$, and $u$ exists, in particular, in $\mathrm{Z}_{1}$.

Furthermore,

$$
\left|u_{1}\right|_{\mathrm{c}^{0}\left(z_{1}\right)} \leqslant \phi(h)
$$

To proceed inductively, let $m$ be an arbitrary positive integer $<\mathrm{T} / h$, and suppose $u_{m}$ to belong to $\mathrm{C}^{1}\left(\mathrm{Z}_{m}\right)$ and to satisfy the condition

$$
\begin{equation*}
\left|u_{m}\right|_{\mathbf{C}^{0}\left(\mathrm{z}_{m}\right)} \leqslant \phi(m h) \tag{8}
\end{equation*}
$$

Since $\phi(m h)<\phi(\mathrm{T})$, considerations like those in the case $m=1$ show that a strict solution $u_{m+1}$ of (A) exists in $\mathrm{Z}_{m+1}$ and conforms to the restriction

$$
\left|u_{m+1}\right|_{\mathrm{C}^{0}\left(\mathrm{z}_{m+1}\right)} \leqslant \phi\left((m+1) h ; m h, w_{m}\right)
$$

where $w_{m}=\phi(m h)$. To complete the induction, it suffices to prove that

$$
\phi\left((m+1) h ; m h, w_{m}\right)=\phi((m+1) h)
$$

For this purpose, define the sequence

$$
\mathrm{F}_{0}=\left|u_{0}\right|_{0}, \mathrm{~F}_{m}=\phi\left(m h ;(m-1) h, \mathrm{~F}_{m-1}\right), m=1,2, \ldots
$$

By definition of $\phi$,

$$
\int_{\mathrm{F}_{i-1}}^{\mathrm{F}_{i}} d v / \mathrm{E}(v)=h \quad \text { for } \quad i=1, \ldots, m
$$

Summing gives us

$$
\int_{\mathrm{F}_{0}}^{\mathrm{F}_{m}} d v / \mathrm{E}=m h
$$

while, again by definition,

$$
\int_{\mathrm{F}_{0}}^{\phi(m h)} d v / \mathrm{E}=m h
$$

It follows that $\mathrm{F}_{\boldsymbol{m}}=\phi(m h)$, as contended. Thus, the induction is complete and inequality $(8)_{m}$, in particular, proved.

We now define a stratified function $v$ in $Z(T)$ as

$$
v=u_{m} \quad \text { in } \quad \mathrm{Z}_{m}^{\prime} \quad \text { for } \quad m=1,2, \ldots, \mathrm{~T} / h
$$

inequality (8) $)_{m}$ implying inequality (8). We shall also denote $v$ by $v_{h}$.

## 4. Precompactness of stratified solutions.

Let $\mathrm{S}(r), r>0$, denote the sphere,

$$
\mathrm{S}(r)=\left\{x: x \in \mathrm{E}^{n},|x|<r\right\}
$$

in $\mathrm{E}^{n}$ of radius $r$. In this section, we shall prove the following result :
Precompactness Theorem. A bounded, measurable function $u$ in $\mathrm{Z}(\mathrm{T})$ and a null sequence $\left\{h_{k}\right\}$ exist such that, for $r>0$, for $t=0$, and for almost all $t$ in the interval $(0, \mathrm{~T})$,

$$
\lim _{k \rightarrow \infty} \int_{\mathrm{S}(r)}\left|v_{h_{k}}(x, t)-u(x, t)\right| d x=0
$$

Our demonstration is in several steps. Let $v=v_{h}$ denote a stratified solution, as constructed in Section 3, with arbitrary h. Let $f \in \mathrm{C}_{0}^{2}\left(\mathrm{E}^{n}\right)$, and for $0 \leqslant t<\mathrm{T}$ set

$$
\mathrm{V}(t)=\mathrm{V}_{h}(t)=\int_{\mathrm{E}^{n}} v_{h}(x, t) f(x) d x
$$

We intend to show that the variation

## var V

of V in the interval $0 \leqslant t<\mathrm{T}$ is finite with a bound that is independent of $h$. Since $v_{h}$ is stratified, $V$ is continuous except for possible jumps at $h, 2 h, \ldots, \mathrm{~T}-h$, and its continuous variation in ( $0, \mathrm{~T}$ ) is given by

$$
\sum_{m=1}^{\mathrm{T} / h} \int_{(m-1) h}^{m h}\left|\mathrm{~V}^{\prime}(t)\right| d t
$$

We can obtain an upper bound on $\left|\mathrm{V}^{\prime}\right|$. The stratified solution $v=v_{h}$ being continuously differentiable for $(m-1) h<t<m h$ and satisfying equation (A), we have

$$
\begin{aligned}
\mathrm{V}^{\prime}(t) & =\int(\partial v / \partial t) f d x \\
& =\int\left(\mathrm{C}-\sum_{i} \partial \mathrm{~A}_{i} / \partial x_{i}\right) f d x \\
& =\int\left(\mathrm{C} f+\sum_{i} f_{x_{i}} \mathrm{~A}_{i}\right) d x
\end{aligned}
$$

Therefore,

$$
\left|\mathrm{V}^{\prime}\right| \leqslant c^{\prime}|f|_{\mathrm{C}^{1}\left(\mathbb{E}^{n}\right)}
$$

where $c^{\prime}$ is a constant depending on $\phi\left(\mathrm{T} ; 0,\left|u_{0}\right|_{0}\right)$ and the support of $f$, and the continuous variation of V in $(0, \mathrm{~T})$ does not exceed

$$
c^{\prime} \mathrm{T}|f|_{\mathrm{c}^{1}\left(\mathbf{E}^{n}\right)}
$$

We estimate the discontinuous variation of $V$ from the formula $\mathrm{V}(m h+0)-\mathrm{V}(m h-0)=\int_{\mathrm{E}^{n}}\left(\mathrm{~K}_{\epsilon} u_{m}(x, m h)-u_{m}(x, m h)\right) f(x) d x$, in which $u_{m}$ is as defined in the previous section. By property (7) of Section 3,

$$
|\mathrm{V}(m h+0)-\mathrm{V}(m h-0)| \leqslant \mathrm{C}_{2}|f|_{\mathrm{C}^{2}\left(\mathbf{E}^{n}\right)} \phi\left(\mathrm{T} ; 0,\left|u_{0}\right|_{0}\right) a(\varepsilon)
$$

Since $\varepsilon$ is proportional to $h$, and $a(\varepsilon)$ is of higher than first order in $\varepsilon$, summing these inequalities for $m=1,2, \ldots, \mathrm{~T} / h$ gives us as upper bound for the discontinuous variation of V in the interval considered

$$
\text { const. }|f|_{\mathrm{C}^{2}\left(\mathrm{E}^{n}\right)} \phi\left(\mathrm{T} ; 0,\left|u_{0}\right|_{0}\right)
$$

Combining this with our previous estimate of the continuous variation of V now gives a bound for var V that is uniform with respect to $h$, as desired.

It follows that a null sequence $\left\{h_{k}\right\}$ exists such that the sequence $\left\{\mathrm{V}_{h_{k}}(t)\right\}$ has a limit as $k \rightarrow \infty$ for every $t$ in the interval $(0, \mathrm{~T})$ (see Natanson [29]).

Furthermore, if $f_{j} \in \mathrm{C}_{0}^{2}\left(\mathrm{E}^{n}\right), j=1,2, \ldots$, by a diagonal process we can find a null sequence, which we again denote by $\left\{h_{k}\right\}$, such that the limits as $k \rightarrow \infty$ of the functions

$$
\mathrm{V}_{h_{k}, j}(t)=\int_{\mathrm{E}}{ }^{n} v_{h_{k}}(x, t) f_{j}(x) d x
$$

exist for all $j=1,2, \ldots$, and all $t$ in the interval $(0, T)$. As $f_{j}$ we can choose a sequence dense in $L^{1}\left(E^{n}\right)$. Since the $v_{h}$ are uniformly bounded (see (8), Section 3), we thereby arrive at the following fact :

Proposition 1. - If $f$ is any bounded, measurable function on $\mathrm{E}^{n}$ with compact support, then

$$
\int_{\mathrm{E}^{n}} v_{h_{k}}(x, t) f(x) d x
$$

converges as $k \rightarrow \infty$ for $0 \leqslant t<\mathrm{T}$.
The next step is to prove :

Proposition 2. - Let $t$ be a fixed value in ( $0, \mathrm{~T}$ ), and consider any subsequence $\left\{h_{\boldsymbol{k}^{\prime}}\right\}$ of $\left\{h_{\boldsymbol{k}}\right\}$. A subsequence $\left\{h_{\boldsymbol{k}^{\prime \prime}}\right\}$ of $\left\{h_{\boldsymbol{k}^{\prime}}\right.$, and a bounded, measurable function $w$ on $\mathrm{E}^{n}$ exist such that

$$
\lim _{k^{\prime \prime} \rightarrow \infty}\left|v_{h_{k^{\prime}}}(\cdot, t)-w\right|_{L^{1}(\mathrm{~S}(r))}=0
$$

for every $r>0$.

Proof. - The main step in justifying this statement is to show that for arbitrary $r$ and $t$ as indicated, the integrals

$$
\int_{\mathrm{S}(r)}\left|v_{h}(x+z, t)-v_{h}(x, t)\right| d x
$$

converge to zero uniformly with respect to $h$ as $|z| \rightarrow 0, z \in \mathrm{E}^{n}$. A result of M. Riesz [30] then will imply that a function $w \in \mathrm{~L}^{1}(\mathrm{~S}(r))$ and a subsequence $\left\{\widetilde{h}_{k^{\prime \prime}}\right\}$ of the sequence $\left\{h_{k^{\prime}}\right\}$ exist such that

$$
\lim _{k^{\prime \prime} \rightarrow \infty}\left|\mathcal{V}_{\tilde{h}_{k^{\prime \prime}}}(\cdot, t)-w\right|_{L^{1}(\mathrm{~S}(r))}=0
$$

A diagonal procedure will produce a further subsequence $\left\{h_{k^{\prime \prime}}\right\}$ of $\left\{h_{k^{\prime}}\right\}$ and a function $w$ bounded and measurable in $\mathrm{E}^{n}$ that satisfy the conclusion of the proposition.

Let

$$
\mathrm{S}_{t, \delta}=\left\{x: x \in \mathrm{E}^{n},|x| \leqslant \mathrm{N}(\mathrm{~T}-t)+\delta+\mathrm{R}\right\}
$$

and set

$$
\begin{aligned}
\mathrm{I}_{\delta}(t)=\mathrm{I}_{\delta}(t, z ; v) & =\int_{\mathrm{S}_{t, \delta}}|v(x+z, t)-v(x, t)| d x \\
& =\int_{\mathrm{S}_{t, \delta}}\left|\mathrm{~T}_{z} v-v\right| d x
\end{aligned}
$$

where $\mathrm{R}>0, \delta \geqslant 0, z \in \mathrm{E}^{n}, 0 \leqslant t<\mathrm{T}$, and $v=v_{h}$ as above. It suffices for the previous argument to prove that for arbitrary R and $t$

$$
\begin{equation*}
\lim _{|z| \rightarrow 0} \mathrm{I}_{0}\left(t, z ; v_{h}\right)=0 \text { uniformly with respect to } h . \tag{1}
\end{equation*}
$$

Theorem 3, Section 2, shows that

$$
\begin{aligned}
& \mathrm{I}_{0}(t) \leqslant e^{\mathrm{C}_{1}[t-(m-1) h]} \mathrm{I}_{0}((m-1) h) \\
& \text { for }(m-1) h \leqslant t<m h, \quad m=1,2, \ldots, \mathrm{~T} / h
\end{aligned}
$$

( T is arbitrary in Theorem 3 ; we replace it at present by $\mathrm{T}+\mathrm{R} / \mathrm{N}$ ).
In addition, we can estimate $\mathrm{I}_{\delta}(m h), m=1, \ldots, \mathrm{~T} / h$, as follows. On the plane $t=m h$, by definition of $v$,

$$
v=\mathrm{K}_{\epsilon} u_{m-1}
$$

therefore,

$$
\mathrm{T}_{z} v=\mathrm{T}_{z} \mathrm{~K}_{\epsilon} u_{m-1}=\mathrm{K}_{\epsilon} \mathrm{T}_{z} u_{m-1}
$$

on this plane, and we have for $\delta \geqslant 0$

$$
\begin{aligned}
\mathrm{I}_{\delta}(m h) & =\left.\int_{\mathrm{S}_{m h, \delta}}\left|\mathrm{~T}_{z} \mathrm{~K}_{\epsilon} u_{m-1}-\mathrm{K}_{\epsilon} u_{m-1}\right| d x\right|^{t=m h} \\
& =\left.\int_{\mathrm{S}_{m h, \delta}}\left|\mathrm{~K}_{\epsilon}\left(\mathrm{T}_{z} u_{m-1}-u_{m-1}\right)\right| d x\right|^{t=m h} \\
& \leqslant \int_{\left.\mathrm{S}_{m h, \delta+\epsilon^{\prime}}\left|\mathrm{T}_{z} u_{m-1}-u_{m-1}\right| d x\right|^{t=m h}} \\
& =\lim _{t \rightarrow m h-0} \mathrm{I}_{\delta+\epsilon^{\prime}}(t)
\end{aligned}
$$

$\varepsilon$ where $\varepsilon^{\prime}=\varepsilon \sqrt{ } n$. This and inequality (2) imply that

$$
\mathrm{I}_{\delta}(m h) \leqslant e^{\mathrm{C}_{1} h} \mathrm{I}_{\delta+\epsilon^{\prime}}((m-1) h) \quad \text { for } \quad m=1, \ldots, \mathrm{~T} / h
$$

It follows that if $q h \leqslant t<(q+1) h$, then

$$
\mathrm{I}_{0}(t) \leqslant e^{\mathrm{C}_{1}(t-q h)} \mathrm{I}_{0}(q h) \leqslant e^{\mathrm{C}_{1} t} \mathrm{I}_{q \epsilon^{\prime}}(0) \leqslant e^{\mathrm{C}_{1} t} \mathrm{I}_{\left(\epsilon^{\prime} / h\right) t}(0)
$$

or, in other words,

$$
\mathrm{I}_{0}\left(t, z ; v_{h}\right) \leqslant \mathrm{I}_{\left(\epsilon^{\prime} / h\right) t}\left(0, z ; \mathrm{K}_{\epsilon} u_{0}\right) e^{\mathrm{C}_{1} t} \quad \text { for } \quad 0 \leqslant t<\mathrm{T}
$$

This estimate justifies (1) and thus completes the proof of Proposition 2.

Propositions 1 and 2 enable us to prove :
Proposition 3. - Suppose that for all $r>0$

$$
\int_{\mathrm{S}(r)}\left|v_{h_{k^{\prime}}}(x, t)-w(x)\right| d x \rightarrow 0 \quad \text { as } \quad k^{\prime} \rightarrow \infty
$$

where $w$ is a bounded, measurable function on $\mathrm{E}^{n}, t$ is a fixed value in $(0, \mathrm{~T})$, and $\left\{h_{k^{\prime}}\right\}$ a subsequence of $\left\{h_{k}\right\}$. Then if $f$ is any bounded, measurable function on $\mathrm{E}^{n}$ with compact support

$$
\begin{equation*}
\lim _{k \rightarrow \infty} \int_{\mathrm{E}^{n}} f(x)\left(v_{h_{k}}(x, t)-w(x)\right) d x=0 \tag{3}
\end{equation*}
$$

Proof. - Taking $r$ so large that $\mathrm{S}(r)$ contains the support of $f$, we have immediately

$$
\lim _{k^{\prime} \rightarrow \infty} \int f(x) v_{h_{k^{\prime}}}(x, t) d x=\int f(x) w(x) d x
$$

the domain of integration in each case being $\mathrm{E}^{\boldsymbol{n}}$. By Proposition 1, however,

$$
\lim _{k \rightarrow \infty} \int f(x) v_{h_{k}}(x, t) d x
$$

exists, and the conclusion stated is obvious.
We can now finish proving the theorem. At most one function $w$ can satisfy (3) for all admitted $f$. Hence, the function $w$ in Proposition 2 -a different function for each $t$ - is unique and independent of the subsequence $\left\{v_{h_{k}}\right\}$. Thus, every subsequence of $v_{h_{k}}(\cdot, t)$ contains a subsequence of itself that converges to $w$ in $\mathrm{L}^{1}(\mathrm{~S}(r))$ for $r \rightarrow 0$. This means that every subsequence of the sequence of real numbers

$$
\left|v_{h_{k}}(\cdot, t)-w\right|_{\mathrm{L}^{1}(\mathrm{~S}(r))} k=1,2, \ldots,
$$

contains a null subsequence of itself. It follows that the sequence itself has the limit 0 , and this implies that for $r>0$ Cauchy's condition

$$
\lim _{\substack{k \rightarrow \infty \\ \ell \rightarrow \infty}}\left|v_{h_{k}}(\cdot, t)-v_{h_{\ell}}(\cdot, t)\right|_{\mathrm{L}^{1}(\mathrm{~S}(r))}=0
$$

is satisfied with any $t$ in ( $0, \mathrm{~T}$ ). Integrating gives us for the cylinder

$$
\mathrm{C}(r, \mathrm{~T})=\{(x, t):|x|<r, 0 \leqslant t<\mathrm{T}\}
$$

the relation

$$
\begin{aligned}
\lim _{k \rightarrow \infty} \mid v_{h_{k}} & -\left.v_{h_{\ell}}\right|_{\mathrm{L}^{1}(\mathrm{C}(r, \mathrm{~T}))} \\
& =\lim _{k \rightarrow \infty} \int_{0}^{\mathrm{T}}\left|v_{h_{k}}(\cdot, t)-v_{h_{\ell}}(\cdot, t)\right|_{\mathrm{L}^{1}(\mathrm{~S}(r))} d t \\
& =0 .
\end{aligned}
$$

By completeness, a function $\mathrm{W} \in \mathrm{L}^{1}(\mathrm{C}(r, \mathrm{~T}))$ exists such that

$$
\lim _{k \rightarrow \infty}\left|v_{h_{k}}-W\right|_{L^{1}(\mathrm{C}(r, \mathrm{~T}))}=0
$$

Hence, for a suitable subsequence $\left\{h_{k^{*}}\right\}$ of $\left\{h_{k}\right\}$

$$
\lim _{k^{*} \rightarrow \infty}\left|v_{h_{k^{*}}}(\cdot, t)-\mathrm{W}(\cdot, t)\right|_{\mathrm{L}^{1}(\mathrm{~S}(r))}=0
$$

for almost all $t$ in $(0, \mathrm{~T})$. On the other hand, for each such $t$, a function $w(\cdot, t)$-that previously denoted by $w$ - exists such that

$$
\lim _{k \rightarrow \infty} \mid v_{h_{k}}(\cdot, t)-w(\cdot, t)_{L^{1}(\mathrm{~S}(r))}=0 .
$$

It follows that for almost all $t$ in $(0, \mathrm{~T})$, the functions $w(\cdot, t)$ and $\mathrm{W}(\cdot, t)$ are equivalent members of $\mathrm{L}^{1}(\mathrm{~S}(r))$ and thus may be identified. Therefore, for almost all $t$ in $(0, \mathrm{~T})$

$$
\lim _{k \rightarrow \infty}\left|v_{h_{k}}(\cdot, t)-\mathrm{W}(\cdot, t)\right|_{\mathrm{L}^{1}(\mathrm{~S}(r))}=0
$$

This being true for arbitrary $r, \mathrm{~W}$ can be regarded as the restriction to $\mathrm{C}(r, \mathrm{~T})$ of a bounded, measurable function $u$ on $\mathrm{Z}(\mathrm{T})$, and this $u$ obviously satisfies the demands of the theorem.

## 5. The existence of a generalized solution of (A), (B).

We shall now conclude the proof of Theorem A, Section 1. In the last section, we obtained in an arbitrary layer $\mathrm{Z}(\mathrm{T})$ a bounded measurable function $u$ such that for all $r>0$

$$
\lim _{k \rightarrow \infty}\left|v_{h_{k}}-u\right|_{L^{1}(\mathcal{C}(r, T))}=0
$$

$\left\{h_{\boldsymbol{k}}\right\}$ being a suitable null sequence and the $\nu_{h_{\boldsymbol{k}}}$ stratified solutions as described. Our aim is to prove $u$ to be a generalized solution of (A), (B) under the definition in Section 1. Choose arbitrarily a constant $k$ and a function $f \in \mathrm{C}_{0}^{1}(\mathrm{Z}(\mathrm{T}))$, and let $r$ be so large that the cylinder $\mathrm{C}(r, \mathrm{~T})$ contains the support of $f$. Theorem 2, Section 2, says that for each $h=h_{j}$

$$
\begin{aligned}
\int_{z_{i}}\left\{f_{t}\left|v_{h}-k\right|\right. & +\operatorname{sign}\left(v_{h}-k\right) \cdot\left[\sum _ { i } f _ { x _ { i } } \left(\mathrm{A}_{i}\left(t, v_{h}\right)\right.\right. \\
& \left.\left.-\mathrm{A}_{i}(t, k)\right)+f \mathrm{C}\right] d x d t \\
& =\left.\int_{\mathrm{E}^{n}} f\left|v_{h}-k\right| d x\right|_{t=(i-1) h} ^{t=i n}
\end{aligned}
$$

Summing over $i$ gives us

$$
\begin{align*}
\int_{\mathrm{Z}(\mathrm{~T})} & \left\{f_{t}\left|v_{h}-k\right|+\operatorname{sign}\left(v_{h}-k\right) \cdot\left[\sum _ { i } f _ { x _ { i } } \left(\mathrm{A}_{i}\left(t, v_{h}\right)\right.\right.\right.  \tag{1}\\
& \left.\left.-\mathrm{A}_{i}(t, k)\right)+f \mathrm{C}\right] d x d t+\int_{\mathbb{E}^{n}} f(x, 0)\left|v_{h}(x, 0)-k\right| d x=R_{h},
\end{align*}
$$

where

$$
\mathrm{R}_{h}=\sum_{i=1}^{\mathrm{T} / h} \int_{\mathrm{E}^{n}} f(x, i h)\left[\left|v_{h}(x, i h-0)-k\right|-\left|v_{h}(x, i h+0)-k\right|\right] d x .
$$

If $k \rightarrow \infty$, the left member of (1) tends toward the left member of inequality (3), Section 1. Hence, $u$ will certainly be a generalized solution of (A), (B) if $\lim _{h \rightarrow 0} \inf _{h} \geqslant 0$. To discuss this, set

$$
f(x)=f(x, i h), v(x)=v_{h}(x, i h-0) .
$$

Then

$$
v_{h}(x, \text { ih }+0)=K_{\epsilon} v(x)
$$

and we have

$$
\left|v_{h}(x, i h+0)-k\right|=\left|\mathrm{K}_{\epsilon} v-k\right|=\left|\mathrm{K}_{\epsilon}(v-k)\right| \leqslant \mathrm{K}_{\epsilon}(|v-k|) .
$$

For the $i-t h$ term in $\mathrm{R}_{h}$,

$$
\begin{aligned}
\int f\left(|v-k|-\mid \mathrm{K}_{\epsilon} v\right. & -k \mid) d x \geqslant \int f\left(|v-k|-\mathrm{K}_{\epsilon}(|v-k|)\right) d x \\
& \geqslant-|f|_{\mathrm{c}^{2}(\mathrm{Z}(\mathrm{~T}))} \mathrm{C}_{2}\left(\phi\left(\mathrm{~T} ; 0,\left|u_{0}\right|_{0}\right)+|k|\right) a(\varepsilon)
\end{aligned}
$$

by property (7), Section 3 ; the bound (8) in that section also has been used. Multiplying the last expression by $T / h$ gives a lower bound for $\mathrm{R}_{h}$. Since $h$ is proportional to $\epsilon$ and $\epsilon^{-1} a(\epsilon) \rightarrow 0$, it follows that the lower limit of $\mathrm{R}_{h}$ as $h \rightarrow 0$ is nonnegative, and condition (3a), Section 2, follows.

To justify the second condition (relation (3b), Section 1) for generalized solutions, consider again stratified solutions $v_{h_{k}}$ with layer thickness $h_{k}$ such that $\lim _{k \rightarrow 0} v_{h_{k}}=u$ almost everywhere. Denoting by $\mathscr{E}$ the set of $t$ in $[0, \mathrm{~T}]$ for which this limit relation fails to hold for almost all $x$, we have $\mu \mathscr{\delta}=0$. We first prove $u_{0}$ to be the weak limit of the $v_{h}(\cdot, t)$ as $t$ approaches 0 in the complement of $\mathscr{E}$. This is meant in the following sense :

Proposition 4. - If $f(x)$ is bounded, measurable, and has compact support in $\mathrm{E}^{n}$, then

$$
\lim _{t \rightarrow 0}^{\prime} \int_{\mathbb{E}^{n}} f(x) u(x, t) d x=\int_{\mathbf{E}^{n}} f(x) u_{0}(x) d x
$$

the prime indicating that $t$ while approaching zero is positive and avoids the null set $\mathcal{E}$.

Proof. -- Since $u_{0}$ and $u$ are bounded, it suffices to prove the lemma for $f \in \mathrm{C}_{0}^{2}\left(\mathrm{E}^{n}\right)$, which henceforth we shall assume. If $v(x, t)$ is a strict solution of $(\mathrm{A})$ in a layer $\mathrm{Z}\left(t_{0}, t_{1}\right)$, integration within this layer of the condition $f\left(v_{t}+\sum_{i}\left(\frac{\partial \mathrm{~A}_{i}}{\partial x_{i}}\right)-\mathrm{C}\right)=0$ and integration by parts give us

$$
\left.\left.\int_{\mathrm{E}^{n}} f(x) v(x, t) d x\right|_{t=t_{0}} ^{t_{1}}=\int_{\mathrm{Z}\left(t_{0}, t_{1}\right)} \sum_{i} f_{x_{i}} \mathrm{~A}_{i}+f \mathrm{C}\right) d x d t
$$

We apply this result to the layers of a stratified solution $v_{h_{\boldsymbol{k}}}$. Summing over the layers that are contained within $\mathrm{Z}(0, t)$, where $t$ is an arbitrary ordinate such that $0<t<\mathrm{T}$, we obtain :

$$
\begin{aligned}
& \int_{\mathrm{E}^{n}} f(x)\left(v_{h_{k}}(x, t)-\mathrm{S}_{\epsilon_{k}} u_{0}(x)\right) d x= \\
& \qquad \int_{\mathrm{Z}(0, t)}\left(\sum_{i} f_{x_{i}} \mathrm{~A}_{i}+f \mathrm{C}\right) d x d t+\mathrm{R}_{k}(t)
\end{aligned}
$$

where $\varepsilon_{k}$ is proportional to $h_{k}$ (Section 2), and as noted previously

$$
\left|\mathrm{R}_{k}(t)\right| \leqslant|f|_{\mathrm{C}^{2}\left(\mathrm{E}^{n}\right)} b_{k} \quad \text { with } \quad \lim _{k \rightarrow 0} b_{k}=0
$$

The integral on the right side can be estimated by $c|f|_{\mathbf{C}^{1}\left(\mathrm{E}^{n}\right)} t$, $c$ depending on upper bounds for the $\left|\mathrm{A}_{i}\right|$ and $|\mathrm{C}|$ and on the support of $f$. Denoting the support of $f$ by $\mathrm{S}_{f}$, we have also

$$
\left|\int_{\mathbf{E}^{n}} f\left(u_{0}-\mathrm{S}_{\epsilon_{k}} u_{0}\right) d x\right| \leqslant|f|_{\mathrm{C}^{0}\left(\mathbf{E}^{n}\right)} c_{k}
$$

where $c_{k}=\int_{\mathrm{S}_{f}}\left|u_{0}-\mathrm{S}_{\epsilon_{k}} u_{0}\right| d x \rightarrow 0$ as $k \rightarrow \infty$. Hence in sum,

$$
\left|\int_{\mathrm{E}^{n}} f(x)\left(v_{h_{k}}(x, t)-u_{0}(x)\right) d x\right| \leqslant|f|_{\mathrm{C}^{2}\left(\mathrm{E}^{n}\right)}\left(c t+b_{k}+c_{k}\right)
$$

For $t \notin \mathscr{E}$, letting $k \rightarrow \infty$ gives

$$
\left|\int_{\mathrm{E}^{n}} f(x)\left(u(x, t)-u_{0}(x)\right) d x\right| \leqslant|f|_{\mathrm{C}^{2}\left(\mathrm{E}^{n}\right)} c t
$$

from which Proposition 4 follows in the case in which $f \in \mathrm{C}_{0}^{2}$. As observed, the entire proposition follows from this case.

Now consider any null sequence $\left\{t_{m}\right\}$ of positive values not in $\delta$. Theorem 3, Section 2, insures that a subsequence $\left\{t_{\boldsymbol{m}^{\prime}}\right\}$ and a bounded measurable function $u^{*}(x)$ exist such that

$$
\lim _{m^{\prime} \rightarrow \infty}\left|u\left(\cdot, t_{m^{\prime}}\right)-u^{*}\right|_{L^{1}(\mathrm{~S}(r))}=0 \quad \text { for } \quad r>0
$$

This and Proposition 4 imply that $\int_{\mathrm{E}^{n}} f(x)\left(u^{*}(x)-u_{0}(x)\right) d x=0$
for arbitrary bounded, measurable $f$ with compact support, and we conclude that $u^{*}=u_{0}$ almost everywhere. Therefore,

$$
\lim _{m^{\prime} \rightarrow 0}\left|u\left(\cdot, t_{m^{\prime}}\right)-u_{0}\right|_{L^{1}(\mathrm{~S}(r))}=0 \quad \text { for } \quad r>0
$$

Since the original sequence $\left\{t_{m}\right\}$ was arbitrary, we conclude that

$$
\lim _{t \rightarrow 0}^{\prime}\left|u(\cdot, t)-u_{0}\right|_{\mathrm{L}^{1}(\mathrm{~S}(r))}=0 \quad \text { for } \quad r>0
$$

which is (3b), the final condition to be justified.

## CHAPTER 3

## METHODS OF SMOOTHING

In this chapter, we first derive formulas for the modifications in an arbitrary function produced by averaging (Sections 1 and 2). These lead us immediately to the estimate applied in Chapter 1 (inequality (7), Section 3) ; Kuznetsov [27] had used a similar inequality, which he proved differently. In Section 3, we go to other smoothing methods of interest in connection with solutions of finite variation, which will be discussed in Chapters 4 and 5.

## 1. Symmetric averaging operators.

Let

$$
\mathrm{C}_{\epsilon}=\left\{x: x=\left(x_{1}, \ldots, x_{n}\right),\left|x_{i}\right| \leqslant \varepsilon, i=1, \ldots, n\right\}
$$

denote the cube in $\mathrm{E}^{n}$ with center at the origin and edges of length $2 \varepsilon(\varepsilon>0)$ parallel to the coordinate axes. Suppose $k$ to be in $\mathrm{L}^{1}\left(\mathrm{C}_{\epsilon}\right)$ and to be an even function with respect to each argument $x_{i}$, $i=1, \ldots, n$. For convenience, we also assume :

$$
k \geqslant 0, \int_{\mathrm{C}_{\epsilon}} k(\xi) d \xi=1
$$

If $u$ denotes an arbitrary bounded, measurable function on $E^{n}$, we define the average of $u$ with kernel $k$ to be the integral

$$
\mathrm{K} u(x)=(\mathrm{K} u)(x)=\int_{\mathrm{C}_{\epsilon}} k(\xi) u(x+\xi) d \xi
$$

For the deviation of $u$ from its average we shall prove that

$$
\begin{equation*}
\mathrm{K} u-u=\sum_{i=1}^{n} \mathrm{D}_{i}^{2} \mathrm{~J}_{i} \tag{1}
\end{equation*}
$$

where $\mathrm{J}_{1}, \ldots, \mathrm{~J}_{n}$ denote bounded, measurable functions defined almost everywhere on $\mathrm{E}^{n}$ such that

$$
\begin{equation*}
\left|J_{i}\right| \leqslant \varepsilon^{2}|u|_{0}, i=1, \ldots, n \tag{2}
\end{equation*}
$$

and $\mathrm{D}_{i}=\partial / \partial x_{i}$ in the distribution sense. If $u \in \mathrm{C}^{1}\left(\mathrm{E}^{n}\right)$, each $\mathrm{J}_{i}$ admits the differentiations indicated, and formula (1) holds literally.

Through integration by parts, this formula implies inequality (7), Section 3, Chapter 1, with $a(\varepsilon)=\varepsilon^{2}$.

Our proof of (1) is based on the one dimensional case, in which case $x=x_{1}$ and we define

$$
\mathrm{J}_{1}(x)=\mathrm{J}(x)=\int_{-\epsilon}^{\epsilon} \mathrm{M}(|s|) u(x+s) d s
$$

where for $0 \leqslant r \leqslant \varepsilon$,

$$
\mathrm{M}(r)=\int_{r}^{\epsilon} d r^{\prime} \int_{r^{\prime}}^{\epsilon} k\left(r^{\prime \prime}\right) d r^{\prime \prime}
$$

The formula in this case asserts that $\mathrm{K} u-u=d^{2} \mathrm{~J} / d x^{2}$. It suffices to establish it for $u \in \mathrm{C}^{1}\left(\mathrm{E}^{n}\right)$. Writing

$$
\mathrm{J}(x)=\int_{-\epsilon}^{0} \mathrm{M}(-s) u(x+s) d s+\int_{0}^{\epsilon} \mathrm{M}(s) u(x+s) d s
$$

we obtain

$$
\begin{aligned}
\mathrm{DJ}=d \mathrm{~J} / d x & =\int_{-\epsilon}^{0} \mathrm{M}(-s) u^{\prime}(x+s) d s+\int_{0}^{\epsilon} \mathrm{M}(s) u^{\prime}(x+s) d s \\
& =\int_{-\epsilon}^{0} \mathrm{M}^{\prime}(-s) u(x+s) d s-\int_{0}^{c} \mathrm{M}^{\prime}(s) u(x+s) d s
\end{aligned}
$$

and by similar manipulations

$$
\begin{aligned}
\mathrm{D}^{2} \mathrm{~J}=d^{2} \mathrm{~J} / d x^{2} & =2 \mathrm{M}^{\prime}(0) u(x)+\int_{-\mathrm{c}}^{0} \mathrm{M}^{\prime \prime}(-s) u(x+s) d s \\
& +\int_{0}^{\epsilon} \mathrm{M}^{\prime \prime}(s) u(x+s) d s=-u(x)+\mathrm{K} u(x)
\end{aligned}
$$

since $\mathrm{M}^{\prime \prime}=k$, which is even. Thus, (1) is justified in the one dimensional case ; inequality (2) is immediate.

A change in notation will help in extending this result to $n$ dimensions. Still with $x=x_{1}$ and $\mathrm{D}=d / d x_{1}$, define

$$
\mathrm{I} k(x)=-\mathrm{M}^{\prime}(|x|)=\int_{|x|}^{\epsilon} k(y) d y
$$

and let $\delta$ denote Dirac's point distribution. The formula in one dimension we have just proved can be written :

$$
\begin{aligned}
\mathrm{D}^{2} \int_{-\epsilon}^{\epsilon} \mathrm{I}^{2} k(y) \cdot u(x+y) d y & =\int_{-\epsilon}^{\epsilon} k(y) u(x+y) d y \\
& -\int_{-\epsilon}^{\epsilon} \delta(y) u(x+y) d y \\
& =\int_{-\epsilon}^{\epsilon}(k(y)-\delta(y)) u(x+y) d y .
\end{aligned}
$$

If the integral

$$
\mathrm{A} k=\int_{-\epsilon}^{\epsilon} k(y) d y
$$

is not necessarily normalized to be 1 , as we previously required, by considering $k / \mathrm{A} k$ in place of $k$ we have, more generally,

$$
\mathrm{D}^{2} \int_{-\epsilon}^{\epsilon} \mathrm{I}^{2} k(y) \cdot u(x+y) d y=\int_{-\epsilon}^{\epsilon}(k(y)-\delta(y)(\mathrm{A} k)) u(x+y) d y .(3)
$$

To prove (1) in the general case, let

$$
\begin{aligned}
& \mathrm{I}_{i} k(x)=\int_{\left|x_{i}\right|}^{\epsilon} k\left(x_{1}, \ldots, x_{i-1}, s, x_{i+1}, \ldots, x_{n}\right) d s \\
& \mathrm{~A}_{i} k(x)=\int_{-\epsilon}^{\epsilon} k\left(x_{1}, \ldots, x_{i-1}, s, x_{i+1}, \ldots, x_{n}\right) d s, \quad i=1, \ldots, n
\end{aligned}
$$

and

$$
k_{0}=k, k_{j}=\left(\prod_{i=1}^{j} \mathrm{~A}_{i}\right) k, j=1, \ldots, n
$$

Then $\mathrm{A}_{i} k$ is independent of the $i$-th coordinate, $k_{j}$ is independant of the first $j$ coordinates, and $k_{n}=1$. Formula (3) shows that if $\psi_{i}$ is an arbitrary function - or distribution- on $\mathrm{C}_{\epsilon}$, but independent of the $i$-th coordinate, we have

$$
\begin{aligned}
\mathrm{D}_{i}^{2} \int_{\mathrm{C}_{\epsilon}} \mathrm{I}_{i}^{2} k(y) \cdot \psi_{i}(y) & u(x+y) d y \\
& =\int_{\mathrm{C}_{\epsilon}}\left[k(y)-\delta\left(y_{i}\right) \mathrm{A}_{i} k(y)\right] \psi_{i}(y) u(x+y) d y
\end{aligned}
$$

Replacing $k$ by $k_{i-1}$ and $\psi_{i}(y)$ by $\prod_{k=1}^{i-1} \delta\left(y_{k}\right)$ gives us the following formulas we shall need :
$\mathrm{D}_{i}^{2} \int_{\mathrm{C}_{\epsilon}} \mathrm{I}_{i}^{2} k_{i-1}(y){\underset{q=1}{i-1} \delta\left(y_{q}\right) u(x+y) d y}^{\eta_{i}}$
$=\int_{\mathrm{C}_{\epsilon}}\left[k_{i-1}(y)-\delta\left(y_{i}\right) k_{i}(y)\right] \prod_{q=1}^{i-1} \delta\left(y_{q}\right) u(x+y) d y, i=1, \ldots, n$
Now to prove (1), we represent the left side as a telescoping sum :

$$
\begin{aligned}
\mathrm{K} u(x)-u(x) & =\int k(y) u(x+y) d y \\
& -\int k_{n}(y) \prod_{q=1}^{n} \delta\left(y_{q}\right) u(x+y) d y \\
& =\int\left[k(y)-k_{1}(y) \delta\left(y_{1}\right)\right] u(x+y) \\
& +\int\left[k_{1}(y)-k_{2}(y) \delta\left(y_{2}\right)\right] \delta\left(y_{1}\right) u(x+y) d y \\
& +\int\left[k_{2}(y)-k_{3}(y) \delta\left(y_{3}\right)\right] \delta\left(y_{1}\right) \delta\left(y_{2}\right) u(x+y) d y \\
& +\cdots \\
& +\int\left[k_{n-1}(y)-k_{n}(y) \delta\left(y_{n}\right)\right] \prod_{q=1}^{n-1} \delta\left(y_{q}\right) u(x+y) d y
\end{aligned}
$$

all integrations being over $C_{\epsilon}$. Formula (1), i.e.,

$$
\mathrm{K} u-u=\sum_{i} \mathrm{D}_{i}^{2} \mathrm{~J}_{i}
$$

with

$$
\begin{equation*}
\mathrm{J}_{i}(x)=\int_{\mathrm{C}_{\epsilon}} \mathrm{I}_{i}^{2} k_{i-1}(y) \prod_{q=1}^{i-1} \delta\left(y_{q}\right) u(x+y) d y \tag{5}
\end{equation*}
$$

results, in view of (4). Since $\mathrm{I}_{i} k_{i-1} \leqslant \frac{1}{2} \mathrm{~A}_{i} k_{i-1}=k_{i} / 2$ and $k_{i}$ is independent of the first $i$ coordinates, we have $\mathrm{I}_{i}^{2} k_{i-1} \leqslant \frac{1}{2} \mathrm{I}_{i} k_{i} \leqslant \frac{\varepsilon}{2} k_{i}$ and

$$
\left|\mathrm{J}_{i}\right| \leqslant \varepsilon^{2}|u|_{0} \int k_{i} d y_{i+1} \ldots d y_{n}=\varepsilon^{2}|u|_{0}
$$

t he last integration being over the domain $\left|y_{p}\right| \leqslant \varepsilon, p=i+1, \ldots, n$. Thus inequality (2) is proved.

Later in this chapter, we shall need an estimate we now derive concerning $\mathrm{D}_{i} \mathrm{~J}_{i}$. For fixed $i=1, \ldots, n$, consider the rectangular portion of hyperplane

$$
\mathrm{H}_{i}=\left\{\left(x_{1}, \ldots, x_{n}\right): x_{i}=c_{i}, a_{k} \leqslant x_{k} \leqslant b_{k}\right.
$$

for $k=1, \ldots, n$, but $k \neq i\}$,
where $c_{i}$ and the $a_{k}$ and $b_{k}$ are constants with $a_{k}<b_{k}$. Obviously, $\mathrm{H}_{i}$ is contained in the rectangular parallelepiped

$$
\mathrm{Q}_{i}^{\epsilon}=\left\{\left(x_{1}, \ldots, x_{n}\right):\left|x_{i}-c_{i}\right| \leqslant \varepsilon, a_{k}-\varepsilon \leqslant x_{k} \leqslant b_{k}+\varepsilon\right.
$$

for $k=1, \ldots, n$, but $k \neq i\}$.
The estimate we shall require is that for $u \in C^{1}\left(\mathrm{Q}_{i}^{\epsilon}\right)$,

$$
\begin{equation*}
\int_{\mathrm{H}_{i}}\left|\mathrm{D}_{i} \mathrm{~J}_{i}\right| d x_{i}^{\prime} \leqslant(\varepsilon / 2) \int_{\mathrm{Q}_{i}^{\epsilon}}\left|\mathrm{D}_{i} u\right| d x \tag{6}
\end{equation*}
$$

where $d x_{i}^{\prime}=d x_{1} \ldots d x_{i-1} d x_{i+1} \ldots d x_{n}$ is the element of volume in $\mathrm{H}_{i}$. To obtain this, we first differentiate $\mathrm{J}_{\boldsymbol{i}}$ with respect to $x_{i}$ under the sign of integration and use again the inequality

$$
\mathrm{I}_{i}^{2} k_{i-1} \leqslant \varepsilon k_{i} / 2
$$

to obtain

$$
\begin{aligned}
\left|\mathrm{D}_{i} \mathrm{~J}_{i}(x)\right| \leqslant(\varepsilon / 2) & \int_{\mathrm{C}_{\epsilon}} k_{i}(y) \prod_{q=1}^{i-1} \delta\left(y_{q}\right)\left|\mathrm{D}_{i} u(x+y)\right| d y \\
& =(\varepsilon / 2) \int k_{i}(y)\left(\int_{-\epsilon}^{\epsilon}\left|\mathrm{D}_{i} u(x+y)\right| d y_{i}\right) \prod_{q=1}^{i-1} \delta\left(y_{q}\right) d y_{i}^{\prime}
\end{aligned}
$$

all $y_{k}$ with $k \neq i$ being integrated over the domain $\left|y_{k}\right| \leqslant \varepsilon$. Integrating with respect to $x$ and $y$ gives

$$
\begin{aligned}
\int_{\mathrm{H}_{i}} & \left|\mathrm{D}_{i} \mathrm{~J}_{i}\right| d x_{i}^{\prime} \\
& \leqslant \frac{1}{2} \varepsilon \int k_{i}(y) \prod_{q=1}^{i-1} \delta\left(y_{q}\right)\left[\int_{\mathrm{H}_{i}}\left(\int_{-\epsilon}^{\epsilon}\left|\mathrm{D}_{i} u(x+y)\right| d y_{i}\right) d x_{i}^{\prime}\right] d y_{i}^{\prime}
\end{aligned}
$$

while the integral in square brackets is $\leqslant \int_{\mathrm{Q}_{i}^{\epsilon}}\left|\mathrm{D}_{\boldsymbol{i}} u\right| d x$. Inequality (6) follows at once.

## 2. Radially symmetric averaging operators.

An averaging kernel $k$ is radially symmetric if $k(x)$ is a function of $|x|$, say

$$
k(x)=\mathrm{R}(|x|)
$$

In that case, we also call the operator $K$ radially symmetric. For radially symmetric averages, a more elegant deviation formula than the previous (equation (1), Section 1) can be given, as follows.

Let

$$
\mathrm{S}_{\epsilon}=\left\{x: x \in \mathrm{E}^{n},|x| \leqslant \varepsilon\right\}
$$

denote the ball in $\mathrm{E}^{n}$ with center at the origin and radius $\varepsilon>0$. Suppose $k$ to be in $\mathrm{L}^{1}\left(\mathrm{~S}_{\epsilon}\right)$ and to be radially symmetric. We also assume :

$$
k \geqslant 0, \int_{\mathbf{S}_{\epsilon}} k(x) d x=1
$$

If $u$ denotes an arbitrary bounded, measurable function on $\mathrm{E}^{n}$, we define the average of $u$ with this kernel $k$ to be the integral

$$
\mathrm{K} u(x)=(\mathrm{K} u)(x)=\int_{\mathbf{S}_{\epsilon}} k(y) u(x+y) d y
$$

For the deviation $\mathrm{K} u-u$ we shall prove in this case that

$$
\begin{equation*}
\mathrm{K} u(x)-u(x)=\Delta \int_{\mathrm{S}_{\epsilon}} \mathrm{M}(|y|) u(x+y) d y \tag{1}
\end{equation*}
$$

where $\Delta=\sum_{i=1}^{n} \partial^{2} / \partial x_{i}^{2}$ in the distribution sense, and for $r>0$

$$
\mathrm{M}(r)=\int_{r}^{\epsilon}\left[\frac{1}{w_{n} \rho^{n-1}}-\int_{0}^{\rho}\left(\frac{\sigma}{\rho}\right)^{n-1} \mathrm{R}(\sigma) d \sigma\right] d \rho
$$

$w_{n}$ denoting the surface area of the unit sphere $|x|=1$ in $\mathrm{E}^{n}$. If $u \in C^{1}\left(\mathrm{E}^{n}\right)$, the differentiations indicated can actually be carried out, and formula (1) holds in the strict sense.

The function $M$ is the solution of the second order singular differential equation

$$
\begin{equation*}
\left(\mathrm{M}^{\prime}(r) r^{n-1}\right)^{\prime}=r^{n-1} \mathrm{R}(r) \tag{2}
\end{equation*}
$$

satisfying the initial conditions

$$
\begin{equation*}
M(\varepsilon)=0 \tag{3}
\end{equation*}
$$

and

$$
\begin{equation*}
\mathrm{M}^{\prime}(\varepsilon)=0 \tag{4}
\end{equation*}
$$

Verifying these statements is straightforward, the condition

$$
\int_{\mathrm{S}_{\epsilon}} k(x) d x=w_{n} \int_{0}^{\epsilon} \mathrm{R}(r) r^{n-1} d r=1
$$

being used in connection with (4). Two other important properties of $M$ are :

$$
\begin{equation*}
\lim _{r \rightarrow 0} r^{n-1} \mathrm{M}(r)=0, w_{n} \lim _{r \rightarrow 0} r^{n-1} \mathrm{M}^{\prime}(r)=-1 \tag{5}
\end{equation*}
$$

It suffices to prove (1) for $u \in \mathrm{C}^{2}\left(\mathrm{E}^{n}\right)$. For such $u$, define

$$
\mathrm{J}(x)=\int_{\mathrm{S}_{\epsilon}} \mathrm{M}(|y|) u(x+y) d y
$$

In spherical coordinates, we have

$$
\mathrm{J}(x)=w_{n} \int_{0}^{\epsilon} \mathrm{M}(r) r^{n-1} u(x, r) d r
$$

$u(x, r)$ denoting the mean value of $u$ on the sphere with center $x$ and radius $r$ :

$$
u(x, r)=\frac{1}{w_{n}} \int_{|\eta|=1} u(x+r \eta) d \mathrm{~S}_{\eta}, r>0
$$

By $d \mathrm{~S}_{\eta}$ is meant the element of area of the sphere, $|\eta|=1$ over which the integration takes place. Since $u$ is continuous, we also define

$$
u(x, 0)=u(x)
$$

It is known (see Courant-Hilbert [5] Vol. II, p. 699) that

$$
r^{n-1} \Delta u(x, r)=\left(r^{n-1} u_{r}(x, r)\right)_{r}
$$

$\Delta$ being the Laplacean operator with respect to $x, x$ and $r$ being regarded as independent variables (in all, $n+1$ variables), and the subscript $r$ indicating partial differentiation with respect to $r$. Using this formula, we have

$$
\Delta \mathrm{J}(x)=w_{n} \int_{0}^{\epsilon} \mathrm{M}(r) r^{n-1} \Delta u(x, r) d r=w_{n} \int_{0}^{\epsilon} \mathrm{M}(r)\left(r^{n-1} u_{r}(x, r)\right)_{r} d r
$$

Integrating by parts gives us, in view of (3) and (5),

$$
\Delta \mathrm{J}(x)=-w_{n} \int_{0}^{\epsilon} \mathrm{M}^{\prime}(r) r^{n-1} u_{r}(x, r) d r
$$

Again integrating by parts and using (4) and (5) we obtain

$$
\begin{aligned}
\Delta \mathrm{J}(x) & =-u(x, 0)+w_{n} \int_{0}^{\epsilon}\left(r^{n-1} \mathrm{M}^{\prime}(r)\right)^{\prime} u(x, r) d r \\
& =-u(x)+w_{n} \int_{0}^{\epsilon} \mathrm{R}(r) r^{n-1} u(x, r) d r \\
& =-u(x)+\int_{\mathrm{S}_{\epsilon}} \mathrm{R}(|y|) u(x+y) d y
\end{aligned}
$$

Thus, formula (1) is proved.

## 3. Other methods of smoothing.

Averaging is only one kind of process meeting the rather delicate requirements -such as inequality (7), Section 3, Chapter 1 - of stratified solutions. Additional ways of smoothing and otherwise modifying functions while satisfying the requirements will be described below ; it is believed that some will be of interest in other problems.

These more general processes pertain to functions of locally finite variation in the sense of Tonelli and Cesari, although in this section we shall keep to the class $C^{1}\left(E^{n}\right)$. (From now on, we revert to the usual meaning of $\mathrm{C}^{k}\left(\mathrm{E}^{n}\right)$, dropping our previous requirement that the derivatives of orders up to $k$ of a member of this class be bounded). If $f \in \mathrm{C}^{1}\left(\mathrm{E}^{n}\right)$ and D is a domain of $\mathrm{E}^{n}$, we define the $T C$-variation of $f$ on D to be

$$
\mathrm{V}(f ; \mathrm{D})=\sum_{i=1}^{n} \int_{\mathrm{D}}\left|\partial f / \partial x_{i}\right| d x
$$

Our criterion pertains to a system of grid points in $\mathrm{E}^{\boldsymbol{n}}$ that must be selected in advance. Changing our previous notation, here let

$$
x=\left(x^{(1)}, x^{(2)}, \ldots, x^{(n)}\right)
$$

denote an arbitrary point in $\mathrm{E}^{n}$. For $i=1, \ldots, n$, divide the $x^{(i)}$ axis by partition points

$$
x_{i_{i}}^{(i)}, j_{i}=0, \pm 1, \pm 2, \ldots,
$$

where

$$
\ldots<x_{-2}^{(i)}<x_{-1}^{(i)}<x_{0}^{(i)}<x_{1}^{(i)}<x_{2}^{(i)}<\ldots
$$

and

$$
x_{-k}^{(i)} \rightarrow-\infty, x_{k}^{(i)} \rightarrow \infty \quad \text { as } \quad k \rightarrow \infty .
$$

The grid points

$$
x_{j}=\left(x_{j_{1}}^{(1)}, x_{j_{2}}^{(2)}, \ldots, x_{i_{n}}^{(n)}\right)
$$

in which $j$ stands for the multi-index $\left(j_{1}, j_{2}, \ldots, j_{n}\right)$, are the vertices of the rectangular parallelepipeds

$$
C_{j}=\left\{x: x_{j_{i}}^{(i)} \leqslant x^{(i)} \leqslant x_{j_{i}+1}^{(i)} \quad \text { for } \quad i=1, \ldots, n\right\}
$$

For any multi-index $j=\left(j_{1}, \ldots, j_{n}\right)$ and for $k=1, \ldots, n$, define the translation

$$
\mathrm{T}_{k} j=\left(j_{1}, \ldots, j_{k-1}, j_{k}+1, j_{k+1}, \ldots, j_{n}\right)
$$

For any quantity $a_{j}$ indexed by $j$, define

$$
\mathrm{T}_{k} a_{j}=a_{\mathrm{T}_{k^{j}}}
$$

for a function $f$ on the grid points $x_{j}$ define

$$
\mathrm{T}_{k} f\left(x_{j}\right)=f\left(\mathrm{~T}_{k} x_{j}\right)
$$

and so forth.
Let Q be a rectangular parallelepiped in $\mathrm{E}^{n}$ that is the union of certain $C_{j}$, say

$$
\mathrm{Q}=\left\{x: x_{i^{\prime}}^{(i)} \leqslant x^{(i)} \leqslant x_{i_{i}^{\prime}}^{(i)}, i=1, \ldots, n\right\}
$$

where $j_{i}^{\prime \prime} \geqslant j_{i}^{\prime}+2$. Let

$$
\begin{aligned}
& \mathrm{Q}_{*}=\left\{x: x_{i_{i}^{\prime}+1}^{(i)} \leqslant x^{(i)} \leqslant x_{i_{i}^{\prime \prime}-1}^{(i)} \quad \text { for } \quad i=1, \ldots, n\right\}, \\
& \mathrm{Q}_{\epsilon}=\left\{x: x_{i_{i}^{\prime}}^{(i)}-\varepsilon \leqslant x^{(i)} \leqslant x_{j_{i}^{\prime}}^{(i)}+\varepsilon \quad \text { for } \quad i=1, \ldots ; n\right\},
\end{aligned}
$$

and

$$
|\Delta|_{\mathrm{Q}}=\max \left(x_{j_{i}+1}^{(i)}-x_{j_{i}}^{(i)}\right),(\Delta)_{\mathrm{Q}}=\min \left(x_{j_{i}+1}^{(i)}-x_{j_{i}}^{(i)}\right),
$$

the maximum and the minimum being with respect to $i$ and $j_{i}$ such that $j_{i}^{\prime} \leqslant j_{i}<j_{i}^{\prime \prime}, i=1, \ldots, n$.

Using some of this notation, we now state the basic result of this section. It takes the form of a condition that a function $v$ approximate zero in a certain sense.

Thiorem ol approximation. - Let $v \in \mathrm{C}^{1}\left(\mathrm{E}^{n}\right)$, and let Q denote a parallelepiped as previously defined. Suppose that for each $\mathrm{C}_{j}$ in Q a set of $n$ real numbers $\mathrm{B}_{1 j}, \mathrm{~B}_{2 j}, \ldots, \mathrm{~B}_{n j}$ exists such that

$$
\begin{equation*}
\int_{\mathrm{C}_{j}} v d x=\sum_{k=1}^{n}\left(\mathrm{~T}_{k}-\mathrm{I}\right) \mathrm{B}_{k j} \tag{1}
\end{equation*}
$$

where I denotes the identity. Suppose also that

$$
\begin{equation*}
\sum_{k=1}^{n} \sum_{j}^{\mathrm{Q}}\left|\mathrm{~B}_{k j}\right| \leqslant c \mathrm{~V}\left(v ; \mathrm{Q}_{\epsilon}\right)|\Delta|_{\mathrm{Q}}, \tag{2}
\end{equation*}
$$

where $\varepsilon$ and $c$ are constants, and $\sum_{j}^{Q}$ indicates summation over such $j$ that $\mathrm{C}_{j} \subset \mathrm{Q}$. Then

$$
\begin{equation*}
\int_{\mathrm{Q}^{*}}|v| d x \leqslant(1+2 c) \mathrm{V}\left(v ; \mathrm{Q}_{\epsilon}\right)|\Delta|_{\mathrm{Q}} \tag{3}
\end{equation*}
$$

and for $f \in \mathrm{C}^{1}\left(\mathrm{E}^{n}\right)$ with $f=0$ in the complement of $\mathrm{Q}^{*}$

$$
\begin{equation*}
\left|\int_{\mathrm{Q}} f v d x\right| \leqslant(1+3 c) \mathrm{V}\left(v ; \mathrm{Q}_{\epsilon}\right)|f|_{\mathrm{C}^{1}\left(\mathrm{E}^{n}\right)}|\Delta|_{\mathrm{Q}}^{2} \tag{4}
\end{equation*}
$$

(It would suffice to assume $v \in \mathrm{C}^{1}\left(\mathrm{Q}_{\epsilon}\right)$ ).

To prove this, we introduce the step function

$$
v_{\Delta}(x)=\frac{1}{\left|\mathrm{C}_{j}\right|} \int_{\mathrm{C}_{j}} v(y) d y=v_{j} \quad \text { for } \quad x \in \mathrm{C}_{j}^{0}
$$

where $C_{j}^{0}$ denotes the interior of $C_{j}$ and $\left|C_{j}\right|$ its measure. We shall first show that

$$
\begin{equation*}
\int_{\mathrm{Q}}\left|v_{\Delta}-v\right| d x \leqslant \mathrm{~V}(v ; \mathrm{Q})|\Delta|_{\mathrm{Q}} \tag{5}
\end{equation*}
$$

Indeed, we have

$$
\int_{\mathrm{Q}}\left|v_{\Delta}-v\right| d x=\sum_{j}^{\mathrm{Q}} \int_{\mathrm{C}_{j}}\left|v_{j}-v(x)\right| d x
$$

while

$$
\begin{aligned}
\int_{\mathrm{C}_{j}}\left|v(x)-v_{j}\right| d x & =\int_{\mathrm{C}_{j}}\left|\frac{1}{\left|\mathrm{C}_{j}\right|} \int_{\mathrm{C}_{j}}(v(x)-v(y)) d y\right| d x \\
& \leqslant \frac{1}{\left|\mathrm{C}_{j}\right|} \int_{\mathrm{C}_{j}} \int_{\mathrm{C}_{j}}|v(x)-v(y)| d x d y
\end{aligned}
$$

Since

$$
\begin{aligned}
|f(x)-f(y)| & \leqslant\left|f\left(x^{(1)}, x^{(2)}, \ldots, x^{(n)}\right)-f\left(y^{(1)}, x^{(2)}, \ldots, x^{(n)}\right)\right| \\
& +\left|f\left(y^{(1)}, x^{(2)}, \ldots, x^{(n)}\right)-f\left(y^{(1)}, y^{(2)}, \ldots, x^{(n)}\right)\right| \\
& +\ldots \\
& +\left|f\left(y^{(1)}, \ldots, y^{(n-1)}, x^{(n)}\right)-f\left(y^{(1)}, \ldots, y^{(n)}\right)\right|
\end{aligned}
$$

we easily see that

$$
\int_{\mathrm{C}_{j}}\left|v(x)-v_{j}\right| d x \leqslant|\Delta|_{\mathrm{Q}} \mathrm{~V}\left(v ; \mathrm{C}_{j}\right),
$$

and (5) follows by summing over $j$.
Inequality (3) now results from the decomposition

$$
\int_{\mathrm{Q}^{*}}|v| d x \leqslant \int_{\mathrm{Q}^{*}}\left|v-v_{\Delta}\right| d x+\int_{\mathrm{Q}^{*}}\left|v_{\Delta}\right| d x=\mathrm{G}_{1}+\mathrm{G}_{2}
$$

since $G_{1}$ is given by (5) and $G_{2}$ by the calculation :

$$
\begin{aligned}
\mathrm{G}_{2} & =\sum_{j}^{\mathrm{Q}_{*}} \int_{\mathrm{C}_{j}}\left|v_{\Delta}\right| d x=\sum_{i}^{\mathrm{Q}_{*}}\left|\int_{\mathrm{C}_{j}} v_{\Delta} d x\right|=\sum_{j}^{\mathrm{Q}_{*}}\left|\int_{\mathrm{C}_{j}} v d x\right| \\
& =\sum_{j}^{\mathrm{Q}_{*}}\left|\sum_{k=1}^{n}\left(\mathrm{~T}_{k}-\mathrm{I}\right) \mathrm{B}_{k j}\right| \leqslant 2 \sum_{j}^{\mathrm{Q}} \sum_{k=1}^{n}\left|\mathrm{~B}_{k j}\right| \\
& \leqslant 2 c \mathrm{~V}\left(v ; \mathrm{Q}_{\epsilon}\right)|\Delta|_{\mathrm{Q}}
\end{aligned}
$$

To prove (4), we write

$$
\int_{\mathrm{Q}} f v d x=\int_{\mathrm{Q}^{*}} f_{\Delta} v d x+\int_{\mathrm{Q}^{*}}\left(f-f_{\Delta}\right) v d x=\mathrm{J}_{1}+\mathrm{J}_{2}
$$

Since $\left|f-f_{\Delta}\right| \leqslant|f|_{c^{1}}|\Delta|_{\mathrm{Q}}$ in Q , we have from (3)

$$
\left|\mathrm{J}_{2}\right| \leqslant|f|_{\mathrm{C}^{1}}|\Delta|_{\mathrm{Q}} \int_{\mathrm{Q}^{*}}|v| d x \leqslant(1+2 c)|f|_{\mathrm{c}^{1}} \mathrm{~V}\left(v ; \mathrm{Q}_{\epsilon}\right)|\Delta|_{\mathrm{Q}}^{2}
$$

$\mathrm{C}^{1}$ abbreviating $\mathrm{C}^{1}\left(\mathrm{E}^{n}\right)$. Since $\int_{\mathrm{Q}^{*}} f_{\Delta}\left(v-v_{\Delta}\right) d x=0$, we also obtain

$$
\mathrm{J}_{1}=\int_{\mathrm{Q}^{*}} f_{\Delta} v_{\Delta} d x=\sum_{j}^{\mathrm{Q}_{*}} f_{j} \sum_{k=1}^{n}\left(\mathrm{~T}_{k}-\mathrm{I}\right) \mathrm{B}_{k j}
$$

where $f_{j}$ denotes the value of $f_{\Delta}$ in $\mathrm{C}_{j}$. Thus,

$$
\begin{aligned}
\mathrm{J}_{1} & =\sum_{j}^{\mathrm{Q}_{*}} \sum_{k=1}^{n} f_{j}\left(\mathrm{~T}_{k}-\mathrm{I}\right) \mathrm{B}_{k j} \\
& =-\sum_{j} \sum_{k=1}^{n}\left(\mathrm{I}-\mathrm{T}_{k}^{-1}\right) f_{j} \cdot \mathrm{~B}_{k j}
\end{aligned}
$$

$f_{j}$ being zero outside $\mathrm{Q}_{*}$. Therefore,

$$
\begin{aligned}
\left|\mathrm{J}_{1}\right| & \leqslant|\Delta|_{\mathrm{Q}}|f|_{\mathrm{c}^{1}} \sum_{k=1}^{n} \sum_{j}^{\mathrm{Q}}\left|\mathrm{~B}_{k j}\right| \\
& \leqslant c|f|_{\mathrm{c}^{1}} \mathrm{~V}\left(v ; \mathrm{Q}_{\epsilon}\right)|\Delta|_{\mathrm{Q}}^{2}
\end{aligned}
$$

inequality (4) following from this and the previous estimate of $\mathrm{J}_{2}$.

Notice that if $f \in \mathrm{C}^{1}(\mathrm{Q})$ with $f=0$ in the complement of $\mathrm{Q}_{*}$ and $v \in \mathrm{C}^{1}(\mathrm{Q})$, then

$$
\begin{equation*}
\left|\int_{\mathrm{Q}} f\left(v-v_{\Delta}\right) d x\right| \leqslant|f|_{\mathrm{C}^{1}(\mathrm{Q})} \mathrm{V}(v ; \mathrm{Q})|\Delta|_{\mathrm{Q}}^{2} \tag{6}
\end{equation*}
$$

In fact, the left member of (6) is equal to

$$
\left|\int_{\mathrm{Q}}\left(f-f_{\Delta}\right)\left(v-v_{\Delta}\right) d x\right| \leqslant|\Delta|_{\mathrm{Q}}|f|_{\mathrm{C}^{1}(\mathrm{Q})} \int_{\mathrm{Q}}\left|v-v_{\Delta}\right| d x
$$

and (5) applies.
In using the Approximation Theorem just proved, we replace $v$ by the deviation $u-S u$ of an arbitrary function $u$ from a presumed approximation $S u$. So applied, this theorem covers, in particular, the averaging operators $K$ of Sections 1 and 2. To verify this, we derive from the representation (1) of Section 1 for an arbitrary function $u \in \mathrm{C}^{1}\left(\mathrm{E}^{n}\right)$ the relation

$$
\begin{aligned}
\int_{\mathrm{C}_{j}}(\mathrm{~K} u-u) d x & =\sum_{k=1}^{n} \int_{\mathrm{C}_{j}} \mathrm{D}_{k}^{2} \mathrm{~J}_{k} d x \\
& =\sum_{k}\left(\int_{\mathrm{C}_{\mathrm{T}_{k}, k}, k} \mathrm{D}_{k} \mathrm{~J}_{k} d x_{k}^{\prime}-\int_{\mathrm{C}_{j, k}} \mathrm{D}_{k} \mathrm{~J}_{k} d x_{k}^{\prime}\right)
\end{aligned}
$$

where

$$
\mathrm{C}_{j, k}=\left\{x: x^{(k)}=x_{j_{k}}^{(i)}, x_{j_{i}}^{(i)} \leqslant x^{(i)} \leqslant x_{j_{i}+1}^{(i)} \quad \text { for } \quad i \neq k\right\} .
$$

This is of the form (1), with $\mathrm{K} u-u=v$ and

$$
\mathrm{B}_{k j}=\int_{\mathrm{C}_{j, k}} \mathrm{D}_{k} \mathrm{~J}_{k} d x_{k}^{\prime}
$$

Furthermore,

$$
\begin{aligned}
\sum_{k} \sum_{j}^{\mathrm{Q}}\left|\mathrm{~B}_{k j}\right| & \leqslant \sum_{k} \sum_{j}^{\mathrm{Q}} \int_{\mathrm{c}_{j, k}}\left|\mathrm{D}_{k} \mathrm{~J}_{k}\right| d x_{k}^{\prime} \\
& =\sum_{k} \sum_{j_{k}=j_{k}^{\prime}}^{i_{k}^{\prime \prime}} \int_{\mathrm{H}_{k, j_{k}}}\left|\mathrm{D}_{k} \mathrm{~J}_{k}\right| d x_{k}^{\prime}
\end{aligned}
$$

where

$$
\mathrm{H}_{k, j_{k}}=\left\{x: x^{(k)}=x_{j_{k}}^{(k)}, x_{\dot{j}_{i}^{\prime}}^{(i)} \leqslant x^{(i)} \leqslant x_{j_{i}^{\prime}}^{(i)} \quad \text { for } \quad i \neq k\right\} .
$$

Hence, by inequality (6), Section 1,

$$
\sum_{k} \sum_{j}^{\mathrm{Q}}\left|\mathrm{~B}_{k j}\right| \leqslant \frac{n}{2} \varepsilon \sum_{k} \int_{\mathrm{Q}_{\epsilon}}\left|\mathrm{D}_{k} u\right| d x=\frac{n}{2} \varepsilon \mathrm{~V}\left(u ; \mathrm{Q}_{\epsilon}\right)
$$

provided

$$
\varepsilon \leqslant 1 / 2 \min _{k}\left(x_{j_{k}^{\prime \prime}}^{(k)}-x_{j_{k}}^{(k)}\right)
$$

Thus, condition (2) holds with $c=\frac{n}{4}$.
We now give an example of a smoothing process that is not averaging in the sense we have considered, but falls under the previous Approximation Theorem. The effect of this process is to approximate an arbitrary $C^{\prime}$ function - more generally, a function of locally bounded variation in the sense of Tonelli and Cesari- by what we shall call a "weathered" step function, i.e., a smooth function coinciding with a step function except near the discontinuities of the latter. Let $\left\{x_{j}\right\}$ be a system of mesh points such as we have described, and relative to this system define the operator $B_{\Delta}$ on the space of bounded, measurable functions $v$ on $\mathrm{E}^{n}$ by

$$
\mathrm{B}_{\Delta} v=v_{\Delta} .
$$

Let $Q$ denote a parallelepiped of the previously indicated type, and choose $\varepsilon$ such that

$$
0<\varepsilon<1 / 2(\Delta)_{\mathrm{Q}}
$$

Let $K_{\epsilon}$ denote an averaging operator, of the kind discussed in Sections 1 or 2 , over a cube or ball of respective edge or diameter $2 \varepsilon$. The smoothing operator we wish to mention here is defined by :

$$
\mathrm{S}_{\epsilon}=\mathbf{K}_{\epsilon} \mathbf{B}_{\Delta}
$$

We apply this operator to functions $u \in C^{1}\left(\mathrm{E}^{n}\right)$. Since

$$
\int_{\mathrm{C}_{\boldsymbol{j}}} \mathrm{B}_{\Delta} u d x=\int_{\mathrm{C}_{\boldsymbol{j}}} u d x
$$

we have

$$
\int_{\mathrm{C}_{j}}\left(\mathrm{~S}_{\epsilon} u-u\right) d x=\int_{\mathrm{C}_{j}}\left(\mathrm{~K}_{\epsilon} \mathrm{B}_{\Delta} u-\mathrm{B}_{\Delta} u\right) d x
$$

while the last integral can be expressed in the form

$$
\sum_{k=1}^{n}\left(\mathrm{~T}_{k}-\mathrm{I}\right) \mathrm{B}_{k j}
$$

where

$$
\begin{equation*}
\sum_{k} \sum_{j}^{\mathrm{Q}}\left|\mathrm{~B}_{k j}\right| \leqslant\left(\frac{n}{2}\right) \varepsilon \mathrm{V}\left(u ; \mathrm{Q}^{*}\right), \tag{7}
\end{equation*}
$$

Q* being the parallelepiped

$$
\mathrm{Q}^{*}=\left\{x: x_{i_{i}^{\prime}-1}^{(i)} \leqslant x^{(i)} \leqslant x_{j_{i}^{\prime}+1}^{(i)} \quad \text { for } \quad i=1, \ldots, n\right\}
$$

This is a consequence of the fact that averaging operators are subject to the Approximation Theorem, which however we have proved only if $u \in C^{1}$. We extend that theorem as follows. Requiring $0<\eta<(\Delta)_{\mathrm{Q}} / 2$, define the approximations $w \equiv w_{\eta}=\mathrm{K}_{\eta} u_{\Delta}$. Rather tedious calculations like the previous enable us to show that

$$
\mathrm{V}\left(w ; \mathrm{Q}_{\epsilon}\right) \leqslant \mathrm{V}\left(u ; \mathrm{Q}^{*}\right)
$$

On the other hand, $w$ being in $C^{1}$ satisfies a condition of the form

$$
\int_{\mathrm{C}_{j}}\left(\mathrm{~K}_{\epsilon} w-w\right) d x=\sum_{k}\left(\mathrm{~T}_{k}-\mathrm{I}\right) \mathrm{B}_{k j}
$$

where

$$
\sum_{k} \sum_{j}^{\mathrm{Q}}\left|\mathrm{~B}_{k j}\right| \leqslant\left(\frac{n}{2}\right) \varepsilon \mathrm{V}\left(w ; \mathrm{Q}_{\epsilon}\right) \leqslant\left(\frac{n}{2}\right) \varepsilon \mathrm{V}\left(u ; \mathrm{Q}^{*}\right)
$$

As $\eta \rightarrow 0$, by means of expression (5), Section 1 , we can verify that the $\mathrm{B}_{k j}$ pertaining to $w$ have limits. These limits, again called $\mathrm{B}_{k j}$, must satisfy (7) and the condition

$$
\int_{\mathrm{C}_{j}}\left(\mathrm{~K}_{\epsilon} u_{\Delta}-u_{\Delta}\right) d x=\sum_{k}\left(\mathrm{~T}_{k}-\mathrm{I}\right) \mathrm{B}_{k j}
$$

This outcome justifies our use of the Approximation Theorem, and we may conclude that

$$
\int_{\mathrm{C}_{j}}\left(\mathrm{~S}_{\epsilon} u-u\right) d x=\sum_{k}\left(\mathrm{~T}_{k}-\mathrm{I}\right) \mathrm{B}_{k j}
$$

as asserted.

CHAPTER 4

## CONSERVATION LAWS WITH INITIAL DATA OF BOUNDED VARIATION

The condition of bounded variation early became prominent in the one-dimensional theory of nonlinear conservation laws and was first used in several dimensions, in the sense of Tonelli and Cesari, by Conway and Smoller [4]. K. Kojima [20] extended their methods and results to conservation laws that might depend upon $x$; we excluded such $x$-dependence from equation (A), Chapter 2. N.N. Kuznetsov [27] used the bounded variation property in his stratifying approach. In this chapter, we wish to justify a variety of stratifying methods based on the smoothing procedures of Chapter 3 when the initial data are of locally bounded variation in the sense of Tonelli and Cesari.

## 1. Functions of bounded variation in the sense of Tonelli and Cesari.

Let $x=\left(x_{1}, \ldots, x_{n}\right)$ be an arbitrary point of $n$-dimensional real Euclidean space $\mathrm{E}^{n}$, let $x_{k}^{\prime}$ denote the point in $\mathrm{E}^{n-1}$ that is obtained by suppressing the $k$-th coordinate of $x$ : for instance,

$$
x_{1}^{\prime}=\left(x_{2}, \ldots, x_{n}\right), x_{2}^{\prime}=\left(x_{1}, x_{3}, \ldots, x_{n}\right)
$$

For $\mathrm{Z} \subset \mathrm{E}^{n}$ and for any function $f$ on $\mathrm{E}^{n}$, define

$$
\mathrm{V}_{1}\left(f ; \mathrm{Z} ; x_{1}^{\prime}\right)=\sup \sum_{i=1}^{p}\left|f\left(s_{i}, x_{1}^{\prime}\right)-f\left(s_{i-1}, x_{1}^{\prime}\right)\right|
$$

taking the supremum with respect to all finite increasing sequences of real numbers $s_{i}, i=0,1, \ldots, p$, such that

$$
\left(s_{i}, x_{1}^{\prime}\right) \notin \mathrm{Z} \quad \text { for } \quad i=0,1, \ldots, p
$$

and also with respect to all positive integers $p$. We call $\mathrm{V}_{1}\left(f ; \mathrm{Z} ; x_{1}^{\prime}\right)$ "the variation of $f$ with respect to $x_{1}$, for fixed $x_{1}^{\prime}$, outside $Z$ ". For
$k=1, \ldots, n$, we define "the variation of $f$ with respect to $x_{k}$, for fixed $x_{k}^{\prime}$, outside Z " analogously, and denote it by $\mathrm{V}_{k}\left(f ; \mathrm{Z} ; x_{k}^{\prime}\right)$. A function $f$ defined in $\mathrm{E}^{n}$ is said to have bounded TC-variation -bounded variation in the sense of Tonelli and Cesari- if a set $\mathrm{Z} \subset \mathrm{E}^{n}$ of measure zero exists such that

$$
\mathrm{V}_{k}\left(f ; \mathrm{Z} ; x_{k}^{\prime}\right) \in \mathrm{L}_{1}\left(\mathrm{E}^{n-1}\right) \quad \text { for } \quad k=1,2, \ldots, n
$$

For such a function, we set

$$
\mathrm{V}(f)=\inf \sum_{k=1}^{n} \int \mathrm{~V}_{k}\left(f ; \mathrm{Z} ; x_{k}^{\prime}\right) d x_{k}^{\prime}
$$

the infimum being over all subsets Z of $\mathrm{E}^{n}$ of measure zero, and the integration being over $\mathrm{E}^{n-1}$. We call $\mathrm{V}(f)$ the $T C$-variation of $f$ on $\mathrm{E}^{n}$.

A function $f$ defined in $\mathrm{E}^{n}$ has locally bounded TC-variation if

$$
\mathrm{V}\left(f \chi_{\mathrm{Q}}\right)<\infty \quad \text { for any cube } \quad \mathrm{Q} \subset \mathrm{E}^{n}
$$

where for any set $E \subset E^{n} \chi_{E}$ denotes the characteristic function of $E$ :

$$
\begin{array}{rlll}
\chi_{\mathrm{E}}(x) & =1 & \text { if } & x \in \mathrm{E} \\
& =0 & \text { if } & x \notin \mathrm{E} .
\end{array}
$$

If $f \in \mathrm{C}^{1}\left(\mathrm{E}^{n}\right)$ with $|f| \leqslant M$, and Q is a cube in $\mathrm{E}^{n}$ with edges of length $s$, then

$$
\mathrm{V}\left(f \chi_{\mathrm{Q}}\right) \leqslant 2 n s^{n-1} \mathrm{M}+\mathrm{V}(f ; \mathrm{Q})
$$

where $\mathrm{V}(f ; \mathrm{Q})$ denotes the TC -variation of $f$ on Q as defined in Section 3, Chapter 3.

To H. Federer [8], E. de Giorgi [12, 13], W.H. Fleming [11], and K. Krickeberg [22] (see also H. Federer [9], Section 4.5)) is due the following criterion of precompactness.

Precompactness Theorem. - Let G denote an infinite family of locally integrable functions defined on $\mathrm{E}^{n}$. Assume that any cube $\mathrm{Q} \subset \mathrm{E}^{n}$ determines constants M and N such that, for any $f \in \mathrm{G}$,

$$
|f(x)| \leqslant M \quad \text { for } \quad x \in \mathrm{Q}
$$

and

$$
\mathrm{V}\left(f \chi_{\mathrm{Q}}\right) \leqslant \mathrm{N} .
$$

Then a locally integrable function $f_{0}$ and $a$ sequence $f_{n} \in G$, $n=1,2, \ldots$, exist such that

$$
\lim _{n \rightarrow \infty} \int_{Q}\left|f_{n}-f_{0}\right| d x=0
$$

for every cube $Q \subset \mathrm{E}^{n}$. The limit $f_{0}$ has locally bounded TC-variation.

## 2. Suitable smoothing operators.

Let $B$ denote the class of functions defined on $\mathrm{E}^{n}$ that are bounded and have bounded, continuous first partial derivatives. For $f \in \mathrm{~B}$, let

$$
|f|_{0}=\sup _{x \in \mathbf{E}^{n}}|f(x)|
$$

if $f$ is a member of B with continuous derivatives of orders up to $k$, .t $|f|_{k}$ analogously denote the supremum of the moduli of the derivatives of $f$ of order $k$ at points of $\mathrm{E}^{n}$.

We permit here families of smoothing operators $S_{\epsilon}, \varepsilon>0$, that map $B$ into $B$ and satisfy seven conditions, pertaining to an arbitrary $f \in \mathrm{~B}$, as follows :
i) Pointwise bounds are preserved : $\left|\mathrm{S}_{\epsilon} f\right|_{0} \leqslant|f|_{0}$.
ii) Sharp gradients are blunted : $\left|\operatorname{grad} \mathrm{S}_{\epsilon} f\right|_{0} \leqslant \mathrm{C}_{1}|f|_{0} / \varepsilon$, where $C_{1}$ is a constant independent of $f, x$, and $\varepsilon$.
iii) If $\bar{f}$ denotes the function defined by $\bar{f}(x)=|f(x)|$, then

$$
\left|\mathrm{S}_{\epsilon} f(x)\right| \leqslant \mathrm{S}_{\epsilon} \bar{f}(x) \quad \text { for } \quad x \in \mathrm{E}^{n}
$$

iv) The smoothed function approximates the original weakly to better than first order in $\varepsilon$. More precisely, for all $\phi \in \mathrm{C}_{0}^{2}\left({ }^{1}\right)$,

$$
\int_{\mathrm{E}^{n}} \phi\left(\mathrm{~S}_{\epsilon} f-f\right) d x \leqslant \mathrm{C}_{2}|\phi|_{2} a(\varepsilon)
$$

( ${ }^{1}$ ) $\mathrm{C}_{0}^{2}$ here denotes the class of functions that are of class $\mathrm{C}^{2}$ in $\mathrm{E}^{n}$ and have compact support.
where $a(\cdot)$ is a function independant of $\phi$ and $f$ such that

$$
\varepsilon^{-1} a(\varepsilon) \rightarrow 0 \quad \text { as } \quad \varepsilon \rightarrow 0
$$

and $\mathrm{C}_{2}$ is a constant independent of $\phi$ and $\varepsilon$, but possibly depending on $|f|_{0}$ and $\mathrm{V}(f ; \operatorname{supp} \phi)$, where supp $\phi$ denotes the support of $\phi$. The last three conditions pertain to an arbitrary cube

$$
\mathrm{Q}=\left\{x: a_{i} \leqslant x_{i} \leqslant b_{i}, \quad i=1, \ldots, n\right\}
$$

with edges paraliel to the axes. Let

$$
\mathrm{Q}_{\epsilon}=\left\{x: a_{i}-\varepsilon \leqslant x_{i} \leqslant b_{i}+\varepsilon, \quad i=1, \ldots, n\right\}
$$

denote the $\varepsilon$-neighborhood of Q . We require for any such cube Q and $f \in \mathrm{~B}$ :
v) $\int_{\mathrm{Q}}\left|\mathrm{S}_{\epsilon} f\right| d x \leqslant \int_{\mathrm{Q}}|f| d x+c_{1} \varepsilon|f|_{0}, c_{1}$ being a constant independent of $f$ and $\varepsilon$, but depending upon Q ;
vi) $\int_{Q}\left|S_{\epsilon} f-f\right| d x \leqslant c_{2} \varepsilon$, where the constant $c_{2}$ may depend upon $\mathrm{Q},|f|_{0}$, and $\mathrm{V}\left(f ; \mathrm{Q}_{\epsilon}\right)$, but otherwise is independent of $f$ and $\varepsilon$;
vii) $\mathrm{V}\left(\mathrm{S}_{\epsilon} f ; \mathrm{Q}\right) \leqslant \mathrm{V}\left(f ; \mathrm{Q}_{\epsilon}\right)$.

These hypotheses are satisfied by the averaging operators of Sections 1-2, Chapter 3, and also by additional smoothing operators such as are described in Section 3 of that chapter. Smoothing operators that lead to "weathered step functions", in particular, satisfy these hypotheses.

## 3. The TC-variation of strict solutions of first order quasi-linear equations.

Conway and Smoller first showed that the TC-variation of an arbitrary weak solution of a first order conservation law behaves regularly, and their result was apparently new even in the case of a solution that is strict. Kuznetsov first proved the analogous thing for strict solutions of first order quasi-linear equations, applying his result to weak solutions of conservation laws by a method of
stratifying. In this section, we give our approach to these problems treated by Kuznetsov. Our calculations are essentially like his, but show to advantage, for instance, in the following special case. For convenience limiting ourselves to three independent variables, consider the partial differential equation

$$
\begin{equation*}
u_{t}+f(u) u_{x}+g(u) u_{y}=0 \tag{1}
\end{equation*}
$$

where $f$ and $g$ are of class $C^{1}$ for all real values of their argument $u$. With $t_{0} \geqslant 0$, we take an initial condition of the form

$$
\begin{equation*}
u\left(x, y, t_{0}\right)=w(x, y) \tag{2}
\end{equation*}
$$

assuming $w$ to be of class $C^{1}$ and to have compact support. We shall prove that $\mathrm{V}(u)=\mathrm{V}(w)$ on every plane $t=$ constant in the layer within which $u$ is a strict solution of (1). More precisely, our result is as follows :

Theorfm 1. - Under the previous assumptions, a strict solution $u(x, y, t)$ of (1), (2) exists in the layer

$$
t_{0} \leqslant t \leqslant t_{0}+h
$$

where

$$
\begin{aligned}
& h=1 /\left[2|\nabla w|_{0} \max _{|s| \leqslant|w|_{0}} \sqrt{\left.f^{\prime}(s)^{2}+g^{\prime}(s)^{2}\right]}\right. \\
& \qquad|\nabla w|_{0}=\sup _{x, y} \sqrt{w_{x}(x, y)^{2}+w_{y}(x, y)^{2}}
\end{aligned}
$$

The solution has compact support in this layer and satisfies the conditions

$$
\begin{equation*}
|u(\cdot, \cdot, t)|_{0} \leqslant|w|_{0} \tag{3}
\end{equation*}
$$

and

$$
\begin{align*}
& \int\left|u_{x}(x, y, t)\right| d x d y=\int\left|w_{x}(x, y)\right| d x d y \\
& \int\left|u_{y}(x, y, t)\right| d x d y=\int\left|w_{y}(x, y)\right| d x d y \tag{4}
\end{align*}
$$

the integrations being over the supports of the functions involved.

Proof. - The method of characteristics gives us a formula for the solution as follows. For an arbitrary point $(\xi, \eta)$ of the initial plane $t=t_{0}$, let

$$
x(t) \equiv x(t ; \xi, \eta), y(t) \equiv y(t ; \xi, \eta), \mathrm{U}(t) \equiv \mathrm{U}(t ; \xi, \eta)
$$

denote solutions of the characteristic differential equations

$$
\frac{d x}{d t}=f(\mathrm{U}), \frac{d y}{d t}=g(\mathrm{U}), \frac{d \mathrm{U}}{d t}=0
$$

satisfying the initial conditions

$$
x\left(t_{0}\right)=\xi, y\left(t_{0}\right)=\eta, \mathrm{U}\left(t_{0}\right)=w(\xi, \eta)
$$

We have immediately

$$
\mathrm{U}(t) \equiv \mathrm{U}(t ; \xi, \eta)=w(\xi, \eta)
$$

Therefore,

$$
\begin{align*}
& x(t) \equiv x(t ; \xi, \eta)=\xi+f(w(\xi, \eta))\left(t-t_{0}\right)  \tag{5}\\
& y(t) \equiv y(t ; \xi, \eta)=\eta+g(w(\xi, \eta))\left(t-t_{0}\right)
\end{align*}
$$

For each fixed $t$, consider the transformation

$$
\begin{equation*}
x=x(t ; \xi, \eta), y=y(t ; \xi, \eta) \tag{6}
\end{equation*}
$$

For the Jacobian determinant

$$
\mathrm{J}=\frac{\partial(x, y)}{\partial(\xi, \eta)}=x_{\xi} y_{\eta}-x_{\eta} y_{\xi}
$$

of this transformation, an explicit calculation shows :

$$
\mathrm{J}=1+\left(f^{\prime}(w) w_{\xi}+g^{\prime}(w) w_{\eta}\right)\left(t-t_{0}\right)
$$

Hence, with $h$ defined as above,

$$
\mathrm{J} \geqslant \frac{1}{2} \quad \text { if } \quad 0 \leqslant t-t_{0} \leqslant h
$$

This implies that the transformation (6), in which $t$ is to be regarded as a parameter, is invertible in the layer $0 \leqslant t-t_{0} \leqslant h$. We denote the inverse transformation by

$$
\xi=\xi(x, y, t), \eta=\eta(x, y, t)
$$

Cauchy's theory says that the function

$$
u(x, y, t)=w(\xi(x, y, t), \eta(x, y, t))
$$

is a solution of (1), (2) in the layer $0 \leqslant t-t_{0} \leqslant h$ and also that this solution is unique.

It is clear that $u$ has compact support and satisfies condition (3). Conditions (4) will result from the relations

$$
\begin{equation*}
u_{x}=\mathrm{J}^{-1} w_{\xi}, u_{y}=\mathrm{J}^{-1} w_{\eta} \tag{8}
\end{equation*}
$$

which we now justify. To prove the first of these relations, we write

$$
u_{x}=w_{\xi} \xi_{x}+w_{\eta} \eta_{x}
$$

and make the substitutions

$$
\begin{aligned}
& \xi_{x}=\mathrm{J}^{-1} y_{\eta}=\mathrm{J}^{-1}\left(1+g^{\prime}(w) w_{\eta}\left(t-t_{0}\right)\right) \\
& \eta_{x}=-\mathrm{J}^{-1} y_{\xi}=-\mathrm{J}^{-1} g^{\prime}(w) w_{\xi}\left(t-t_{0}\right)
\end{aligned}
$$

The result is the first formula of (8), and the second is obtained in a similar way.

By the first formula of (8), we have

$$
\int\left|u_{x}(x, y, t)\right| d x d y=\int\left|w_{\xi}(\xi, \eta)\right| \mathrm{J}^{-1} d x d y=\int\left|w_{\xi}(\xi, \eta)\right| d \xi d \eta
$$

giving the first formula of (4). The second formula of (4) is obtained from the second formula of (8). Thus the theorem is proved.

The idea just employed in the special case of an equation of the form (1) can be extended to the more general equations

$$
\begin{equation*}
u_{t}+\sum_{i=1}^{n^{\prime}} a_{i}(x, t, u) u_{x_{i}}=c(x, t, u) \tag{9}
\end{equation*}
$$

under assumptions (i) and (ii) formulated in the proof of Theorem 1, Section 2, Chapter 2. Concerning the initial condition

$$
\begin{equation*}
u\left(x, t_{0}\right)=w(x) \tag{10}
\end{equation*}
$$

we assume, as in the earlier section, that $w$ and grad $w$ are bounded and continuous in $\mathrm{E}^{n}$. Thus, in particular, (iii) $w$ has locally bounded TC-variation.

Repeating the proof of the theorem referred to, we find that a strict solution of (9), (10) exists in a layer $\mathrm{Z}\left(t_{0}, t_{0}+h\right)$, where $h$ is an arbitrary positive number not greater than $1 /\left(\mathrm{A}|\operatorname{grad} w|_{0}+a\right)$; $a$ and A are constants depending on the coefficients in (9) and on $|w|_{0}$, but not otherwise on $w$. Furthermore, the solution is subject to the bound

$$
|u(x, t)| \leqslant \phi\left(t, t_{0},|w|_{0}\right)<\phi\left(t_{0}+h, t_{0},|w|_{0}\right)
$$

for

$$
\begin{equation*}
(x, t) \in \mathrm{Z}\left(t_{0}, t_{0}+h\right) \tag{11}
\end{equation*}
$$

Let

$$
\mathrm{N}^{\prime}=\sup _{Z\left(t_{0}, t_{0}+h ; \phi\left(t_{0}+h ; t_{0},|w|_{0}\right)\right)}\left[\sum_{i} a_{i}^{2}\right]^{1 / 2}
$$

With $t_{1}>t_{0}+h, x_{1} \in \mathrm{E}^{n}$, and $\mathrm{N} \geqslant \mathrm{N}^{\prime}$, define the disk

$$
\mathrm{D}_{t}=\left\{x:\left|x-x_{1}\right| \leqslant \mathrm{N}\left(t_{1}-t\right)\right\} .
$$

Since $N \geqslant N^{\prime}$, as is well known -and follows from the proof of the theorem in Chapter 2 to which we have been referring- the values of $u(x, t)$ for $x \in \mathrm{D}_{t}$ and $t_{0} \leqslant t \leqslant t_{0}+h$ are determined by the values of $w$ on $\mathrm{D}_{t_{0}}$. Our aim is the following estimate, essentially first given by Kuznetsov :

Theorem 2. - Under the indicated hypotheses, the solution of (9) and (10) satisfies a condition of the form

$$
\mathrm{V}\left(u ; \mathrm{D}_{t}\right) \leqslant e^{\mathrm{B}\left(t-t_{0}\right)} \mathrm{V}\left(w ; \mathrm{D}_{t_{0}}\right)+b\left(e^{\mathrm{B}\left(t-t_{0}\right)}-1\right)
$$

where $b=\left(\omega_{n} / n\right) \mathrm{N}^{n}\left(t_{1}-t_{0}\right)^{n}$, and B depends upon bounds for the first derivatives of the coefficients in (9). Here, $\omega_{n}$ represents the area of the unit sphere in $\mathrm{E}^{n}$.

Proof. - We continue to use the terminology and results of Section 2, Chapter 2. Let

$$
\mathrm{P}(t) \equiv \mathrm{V}\left(u ; \mathrm{D}_{t}\right)=\int_{\mathrm{D}_{t}} \sum_{i}\left|u_{x_{i}}(x, t)\right| d x
$$

Changing from $x$ to $\xi$, we have in terms of the quantities $\mathrm{P}_{i}=p_{i} \mathrm{~J}$,

$$
\mathrm{P}(t)=\int_{\mathrm{D}_{t_{0}}} \sum_{i}\left|p_{i}(t ; \xi)\right| \mathrm{J}(t ; \xi) d \xi=\int_{\mathrm{D}_{t_{0}}} \sum_{i}\left|\mathrm{P}_{i}(t ; \xi)\right| d \xi
$$

for $t_{0} \leqslant t \leqslant t_{0}+h, \mathrm{~J}$ being positive in that layer.
Our aim is to estimate the last expression with the help of the expansion

$$
\begin{equation*}
\mathrm{P}_{i}(t ; \xi)=\mathrm{P}_{i}\left(t_{0} ; \xi\right)+\int_{t_{0}}^{t} \dot{\mathrm{P}}_{i}(s ; \xi) d s \tag{12}
\end{equation*}
$$

where $\dot{\mathrm{P}}_{i}=d \mathrm{P}_{i} / d t$. We have $\dot{\mathrm{P}}_{i}=\dot{p}_{i} \mathrm{~J}+p_{i} \dot{\mathrm{~J}}$, and by substitution from equations (11) and (14) of Section 2, Chapter 2,

$$
\dot{\mathrm{P}}_{i}=\sum_{k}\left[a_{k, x_{k}} \mathrm{P}_{i}-a_{k, x_{i}} \mathrm{P}_{k}\right]+c_{u} \mathrm{P}_{i}+c_{x_{i}} \mathrm{~J}
$$

the terms quadratic in the $p_{k}$ all dropping out. Therefore,

$$
\sum_{i}\left|\dot{\mathrm{P}}_{i}\right| \leqslant \mathrm{B} \sum_{i} \mathrm{P}_{i}+\mathrm{BJ},
$$

where $B$ is a constant depending upon upper bounds for the absolute values of the first derivatives of the $a_{i}$ and $c$. Hence, and because $\mathrm{P}_{i}\left(t_{0} ; \xi\right)=w_{x_{i}}(\xi)$, we have from (12)

$$
\sum_{i}\left|\mathrm{P}_{i}(t ; \xi)\right| \leqslant \sum_{i}\left|w_{x_{i}}(\xi)\right|+\mathrm{B} \int_{t_{0}}^{t}\left[\sum_{i}\left|\mathrm{P}_{i}(s ; \xi)\right|+\mathrm{J}(s ; \xi)\right] d s .
$$

Integrating with respect to $\xi$ over $D_{t_{0}}$ and changing the order of the integrations with respect to $s$ and $\xi$ gives us :

$$
\mathrm{P}(t) \leqslant \mathrm{V}\left(w ; \mathrm{D}_{t_{0}}\right)+\mathrm{B} \int_{t_{0}}^{t} \mathrm{P}(s) d s+\mathrm{B} \int_{t_{0}}^{t}\left(\int_{\mathrm{D}_{t_{0}}} \mathrm{~J}(s ; \xi) d \xi\right) d s
$$

The last term is easily calculated, and we have :

$$
\begin{aligned}
\int_{t_{0}}^{t}\left(\int_{\mathrm{D}_{t_{0}}} \mathrm{~J}(s ; \xi) d \xi\right) d s= & \int_{t_{0}}^{t}\left(\int_{\mathrm{D}_{s}} d x\right) d s \\
& =\left(\omega_{n} / n\right) \int_{t_{0}}^{t}\left[\mathrm{~N}\left(t_{1}-s\right)\right]^{n} d s \leqslant b\left(t-t_{0}\right)
\end{aligned}
$$

$b$ being the constant previously defined. If we replace the last term on the right side of (13) by this estimate, we arrive by using Gronwall's reasoning at the inequality asserted in the theorem.

## 4. Stratified solutions and weak solutions.

Given $\mathrm{T}>0$, determine constants A and $a$ as in the previous section. Fix $h$ as a positive number not greater than $1 /(2 a)$ and such that, for convenience, $T / h$ is an integer. The procedures of Section 3, Chapter 2, enable us to construct a stratified solution $v$ with layer thickness $h$ if we perform our smoothing with $\mathrm{S}_{\epsilon}$, where

$$
\varepsilon=2 \mathrm{C}_{1} \mathrm{~A} \phi\left(\mathrm{~T} ; 0,\left|u_{0}\right|_{0}\right) h .
$$

We will have

$$
\begin{equation*}
|v|_{\mathrm{C}^{0}(\mathrm{Z}(0, \mathrm{~T}))} \leqslant \phi\left(\mathrm{T} ; 0,\left|u_{0}\right|_{0}\right) \tag{1}
\end{equation*}
$$

and, in every stratum, will be able to use the determination

$$
\mathrm{N}=\sup _{\mathrm{Z}\left(0, \mathrm{~T} ; \phi\left(\mathrm{T} ; 0,\left|u_{0}\right|_{0}\right)\right)}\left[\sum_{i} a_{i}^{2}\right]^{1 / 2}
$$

when applying Theorem 2 of the previous section.
To describe the TC-variation of $v$, we introduce suitable "stepped" cones, as follows. With arbitrary $t_{1}>\mathrm{T}$ and $x_{1} \in \mathrm{E}^{n}$, let

$$
\mathrm{S}_{k}=\left\{(x, t):\left|x-x_{1}\right| \leqslant \mathrm{N}\left(t_{1}-t\right)\right.
$$

$$
\left.+\left(\mathrm{Th}^{-1}-k\right) \varepsilon, k h \leqslant t<(k+1) h\right\}
$$

be a frustum of a "cone of determinacy" in $Z_{k+1}$. For

$$
k h \leqslant \tau<(k+1) h, \quad \text { let } \quad \mathrm{E}(\tau)=\left\{(x, \tau):(x, \tau) \in \mathrm{S}_{k}\right\}
$$

be a "horizontal" section of $\mathrm{S}_{\boldsymbol{k}}$ at the ordinate $\tau$. We denote the base of $S_{k}$ by

$$
\mathrm{E}_{k}=\mathrm{E}(k h)=\left\{(x, k h):\left|x-x_{1}\right| \leqslant \mathrm{N}\left(t_{1}-k h\right)+\left(\mathrm{Th}^{-1}-k\right) \varepsilon\right\}
$$

and we denote the top of $S_{k-1}$ by

$$
\mathrm{F}_{k}=\left\{(x, k h):\left|x-x_{1}\right| \leqslant \mathrm{N}\left(t_{1}-k h\right)+\left(\mathrm{Th}^{-1}-k+1\right) \varepsilon\right\} .
$$

The union of the $\mathrm{S}_{k}$ for $k=0,1, \ldots, \mathrm{~T} / h$ is a figure we might describe as a portion of a stepped cone. We shall prove that for any of its horizontal sections $\mathrm{E}(t)$, we have :

$$
\begin{equation*}
\mathrm{V}(v ; \mathrm{E}(t)) \leqslant e^{\mathrm{B} t} \mathrm{~V}\left(u_{0} ; \mathrm{E}(0)\right)+b\left(e^{\mathrm{B} t}-1\right) \tag{2}
\end{equation*}
$$

We have first

$$
\mathrm{V}\left(v ; \mathrm{E}_{k}\right)=\mathrm{V}(v ; \mathrm{E}(k h))=\mathrm{V}\left(\mathrm{~S}_{\epsilon} v ; \mathrm{E}(k h-0)\right)
$$

the right hand member representing $\lim _{\sigma \downarrow 0} \mathrm{~V}\left(\mathrm{~S}_{\epsilon} v ; \mathrm{E}(k h-\sigma)\right.$; this is by definition of $v(x ; k h)$ in a stratified solution. Property (vii) of smoothing operators (Section 2) shows that

$$
\mathrm{V}\left(\mathrm{~S}_{\epsilon} v ; \mathrm{E}(k h-0)\right) \leqslant \mathrm{V}\left(v ; \mathrm{F}_{k}\right)
$$

while by Theorem 2, Section 3,

$$
\mathrm{V}\left(v ; \mathrm{F}_{k}\right) \leqslant e^{\mathrm{B} h} \mathrm{~V}\left(v ; \mathrm{E}_{k-1}\right)+b\left(e^{\mathrm{B} h}-1\right)
$$

In sum,

$$
\mathrm{V}\left(v ; \mathrm{E}_{k}\right) \leqslant e^{\mathrm{B} h} \mathrm{~V}\left(v ; \mathrm{E}_{k-1}\right)+b\left(e^{\mathrm{B} h}-1\right) \quad \text { for } \quad k=1,2, \ldots, \mathrm{~T} / h
$$

From this, by induction, we have
$\mathrm{V}\left(v ; \mathrm{E}_{k}\right) \leqslant e^{\mathrm{B} k h} \mathrm{~V}\left(v ; \mathrm{E}_{0}\right)+b\left(e^{\mathrm{B} k h}-1\right) \quad$ for $\quad k=1,2, \ldots \mathrm{~T} / h$, which is of the desired form (2) for $t=k h$. For $k h<t<(k+1) h$, Theorem 2, Section 3, tells us that

$$
\mathrm{V}(v ; \mathrm{E}(t)) \leqslant e^{\mathrm{B}(t-k h)} \mathrm{V}\left(v ; \mathrm{E}_{k}\right)+b\left(e^{\mathrm{B}(t-k h)}-1\right)
$$

and substitution for $\mathrm{V}\left(v ; \mathrm{E}_{\boldsymbol{k}}\right)$ from the previous result leads to (2) in the present case as well. In this way, inequality (2) is completely verified.

With inequalities (1) and (2) established, the precompactness theorem of Section 1 can be applied and then the limit function shown to be a weak solution, for instance by the methods of Chapter 2. We omit further details.

## CHAPTER 5

## EQUATIONS OF HAMILTON-JACOBI TYPE

The global theory of Hamilton-Jacobi equations is very similar to that of quasilinear conservation laws. That this should be so is obvious in the one-dimensional case, the Hamilton-Jacobi equation $u_{t}+f\left(u_{x}\right)=0$, for instance, leading to the conservation law

$$
\left.v_{t}+f(v)\right)_{x}=0
$$

with $v=u_{x}$.
Multi-dimensional equations of Hamilton-Jacobi type have been treated by several authors : S.N. Kruzhkov [23] employed artificial viscosity, E.D. Conway and E. Hopf [3] used a formulation within the calculus of variations, E. Hopf [17] adapted the method of envelopes ( ${ }^{1}$ ), A. Douglis [7] set up approximating problems in which $x$ is discrete, and W.H. Fleming [11,1-4] applied the stochastic calculus of variations and also control theory and the theory of differential games. Uniqueness theorems have been given by A. Douglis [7] and S.N. Kruzhkov [24].

## 1. Aims and assumptions.

Equations of the Hamilton-Jacobi type

$$
\begin{equation*}
u_{t}+f(x, t, u, \operatorname{grad} u)=0 \tag{F}
\end{equation*}
$$

permit the method of stratified solutions to be applied more simply than do conservation laws. For given $\mathrm{T}>0$, suppose $f$ to be of class $C^{2}$ in the $(2 n+2)$-dimensional region

$$
\mathrm{W}(\mathrm{~T})=\left\{(x, t, u, p): x \in \mathrm{E}^{n}, 0 \leqslant t \leqslant \mathrm{~T}, u \in \mathrm{E}^{1}, p \in \mathrm{E}^{n}\right\}
$$

and suppose that for any constant V the derivatives

$$
f_{p_{i}}(x, t, u, p), i=1, \ldots, n
$$

$\left({ }^{1}\right)$ Applied to boundary problems by S. Aigawa and N. Kikuchi [1-1] and by S.H. Benton, Jr. [1-3].
are bounded on the subregion

$$
\left\{(x, t, u, p): x \in \mathrm{E}^{n}, 0 \leqslant t \leqslant \mathrm{~T},|u| \leqslant \mathrm{V},|p| \leqslant \mathrm{V}\right\}
$$

We have used the notation $f_{p_{i}}=\partial f / \partial p_{i}$ and shall also abbreviate $\partial f / \partial t, \partial f / \partial x_{i}, \partial f / \partial u$ by $f_{t}, f_{x_{i}}, f_{u}$, respectively, with analogous notation for second derivatives.

Prescribing an initial condition of the form

$$
\begin{equation*}
u(x, 0)=u_{0}(x) \tag{G}
\end{equation*}
$$

where

$$
\begin{equation*}
\left|u_{0}(x)\right| \leqslant \mathrm{U},\left|u_{0}(x)-u_{0}\left(x^{\prime}\right)\right| \leqslant \mathrm{U}\left|x-x^{\prime}\right| \quad \text { for } \quad x, x^{\prime} \in \mathrm{E}^{n} \tag{1}
\end{equation*}
$$

with constant U , we shall ask for a Lipschitz continuous function $u$ satisfying (G) in $\mathrm{E}^{\boldsymbol{n}}$ and (F) almost everywhere in

$$
\mathrm{Z}(\mathrm{~T})=\left\{(x, t): x \in \mathrm{E}^{n}, 0 \leqslant t \leqslant \mathrm{~T}\right\} .
$$

To obtain such solutions, we make further assumptions concerning $f$.
The first additional assumption is the following convexity condition : a positive constant $\mu$ shall exist such that for any vector $\xi=\left(\xi_{1}, \ldots, \xi_{n}\right) \in \mathrm{E}^{n}$ we have

$$
\begin{equation*}
\sum_{i, j} f_{p_{i} p_{j}}(x, t, u, p) \xi_{i} \xi_{j} \geqslant \mu|\xi|^{2} \quad \text { for } \quad(x, t, u, p) \in \mathrm{W}(\mathrm{~T}) \tag{2}
\end{equation*}
$$

where $|\xi|^{2}=\xi_{1}^{2}+\ldots+\xi_{n}^{2}$. This and the previous suppositions already imply that at most one Lipschitz continuous function $u$ can satisfy (G) in $E^{n}$, (F) almost everywhere in $Z(T)$, and a "semiconcavity" requirement of the form
$u(x+\xi, t)+u(x-\xi, t)-2 u(x, t) \leqslant \mathrm{K}(t)|\xi|^{2}$ for $x, \xi \in \mathrm{E}^{n},(3)$
where $\mathrm{K}(t)<\infty$ for $t>0$. A Lipschitz continuous function $u$ that fulfills all these conditions is called a generalized solution of (F) and (G) in $\mathrm{Z}(\mathrm{T})$.

Now we make assymptions concerning the growth of $f$ and its first and second derivatives. We presume a positive nondecreasing function $E(\rho)$ to exist for $\rho \geqslant 0$ such that

$$
\begin{equation*}
\int_{a}^{\infty} d \rho / \mathrm{E}(\rho)=\infty \quad \text { for any } \quad a>0 \tag{4}
\end{equation*}
$$

and that if $\Psi$ represents any of the $n+1$ expressions

$$
-f+\sum_{j} p_{i} f_{p_{j}}, f_{x_{i}}+p_{i} f_{u}, i=1, \ldots, n
$$

then

$$
\begin{equation*}
|\Psi(x, t, u, p)| \leqslant \mathrm{E}(\rho) \tag{5}
\end{equation*}
$$

for

$$
\begin{equation*}
x \in \mathrm{E}^{n}, 0 \leqslant t \leqslant \mathrm{~T},|u| \leqslant \rho,|p| \leqslant \rho \tag{6}
\end{equation*}
$$

This assumption will be responsible for an a priori estimate of the Lipschitz constant for a solution of (F), (G). The next, our final supposition, guarantees that strict solutions will have domains of suitable widths. It is that if $\Phi$ represents $f$ or any one of its first or second derivatives, then

$$
\begin{equation*}
\sup |\Phi(x, t, u, p)| \leqslant \mathrm{B}(\rho) \tag{7}
\end{equation*}
$$

the supremum being taken for the values (6), and B( $\rho$ ) being a nondecreasing function finite for $\rho \geqslant 0$.

Our principal aim is to prove :
Throrfm C. Under the assumptions stated, a generalized solution of (F), (G) exists in $\mathrm{Z}(\mathrm{T})$.

## 2. Strict solutions of equation ( $F$ ).

By a strict solution of ( F ) we mean a continuously differentiable function that satisfies (F) at all points of its domain. We are concerned here with strict solutions $u$ of (F) defined in layers of the type

$$
\mathrm{Z}\left(t_{0}, t_{0}+h\right)=\left\{(x, t): x \in \mathrm{E}^{n}, t_{0} \leqslant t \leqslant t_{0}+h\right\}
$$

and satisfying initial conditions of the form

$$
\begin{equation*}
u\left(x, t_{0}\right)=w(x) \tag{1}
\end{equation*}
$$

where $w$ is an arbitrary bounded function of class $C^{2}$ with bounded first and second derivatives in $E^{n}$. Let $W$ and $M$ be constants such that

$$
\begin{equation*}
|w| \leqslant \mathrm{W},\left|\partial w / \partial x_{i}\right| \leqslant \mathrm{W},\left|\partial^{2} w / \partial x_{i} \partial x_{j}\right| \leqslant \mathrm{M} \text { in } \mathrm{E}^{n} \tag{2}
\end{equation*}
$$

for $i, j=1, \ldots, n$.
With fixed $t_{0}$ in the interval $0 \leqslant t_{0}<\mathrm{T}$, let $y\left(t ; t_{0}\right)$ be the function that satisfies the differential equation $d y / d t=\mathrm{E}(y)$ and the initial condition $y\left(t_{0}\right)=\mathrm{W}$. Since this function is characterized by the relation $\int_{\mathrm{W}}^{y\left(t ; t_{0}\right)} d z / \mathrm{E}(z)=t-t_{0}$, in view of condition (4), Section 1 , it exists for $t \geqslant t_{0}$; it is positive and increasing. Let A be a constant such that

$$
\begin{aligned}
\left|f_{p_{i}}(x, t, u, p)\right| \leqslant & \frac{\mathrm{A}}{\sqrt{n}} \text { for } x \in \mathrm{E}^{n}, t_{0} \leqslant t \leqslant \mathrm{~T} \\
& |u| \leqslant y\left(\mathrm{~T} ; t_{0}\right),|p| \leqslant y\left(\mathrm{~T} ; t_{0}\right), i=1, \ldots, n .
\end{aligned}
$$

Any solid conical layer of the form

$$
\begin{aligned}
\mathrm{L} \equiv \mathrm{~L}(h) & \equiv \mathrm{L}\left(x_{1}, t_{1} ; t_{0}, h\right) \\
& =\left\{(x, t):\left|x-x_{1}\right| \leqslant \mathrm{A}\left(t_{1}-t\right), t_{0} \leqslant t \leqslant t_{0}+h\right\}
\end{aligned}
$$

where $x_{1} \in \mathrm{E}^{n}$ and $t_{0}<t_{0}+h<\mathrm{T}<t_{1}$, will be called a (conical) layer of determinacy of $u$. As will be seen, the values of $u$ in L will be determined by those of $w$ prescribed on the base

$$
\mathrm{B} \equiv \mathrm{~B}\left(x_{1}, t_{1} ; t_{0}\right)=\left\{\left(x, t_{0}\right):\left|x-x_{1}\right| \leqslant \mathrm{A}\left(t_{1}-t_{0}\right)\right\}
$$

We can now formulate our main results concerning strict solutions of (F). In the following theorem, C and $\mathrm{C}^{\prime}$ denote appropriate constants depending on W and on bounds for $f$ and its partial derivatives of first and second order ; C and $\mathrm{C}^{\prime}$ do not depend upon $t_{0}, \mathrm{M}$, or $\mu$.

Theorlem. - Under the previous hypotheses, the initial value problem (F), (G) has a strict solution $u$ in $Z\left(t_{0}, t_{0}+h\right)$, where
i) $h=1 /[\mathrm{C}(1+\mathrm{M})]$.

In $\mathrm{Z}\left(t_{0}, t_{0}+h\right)$,
ii) $|u(x, t)| \leqslant y\left(t ; t_{0}\right),\left|\partial u(x, t) / \partial x_{i}\right| \leqslant y\left(t ; t_{0}\right)$

$$
\text { for } i=1, \ldots, n,
$$

and
iii) $\left|\partial^{2} u / \partial x_{i} \partial x_{j}\right| \leqslant 2 n \mathrm{M}+1 \quad$ for $\quad i, j=1, \ldots, n$.

With an arbitrary unit vector $z=\left(z_{1}, \ldots, z_{n}\right),|z|=1$, define

$$
\begin{gathered}
\mathrm{P}^{*}(t ; z) \equiv \mathrm{P}^{*}(t ; z \mid u)=\sup _{\mathrm{E}^{n}} \sum_{i, j} \frac{\partial^{2} u}{\partial x_{i} \partial x_{j}} z_{i} z_{j} \\
\mathrm{P}(t) \equiv \mathrm{P}(t ; z)=\mathrm{P}(t ; z \mid u)=\max \left(\mathrm{C}^{\prime} / \mu, \mathrm{P}^{*}(t ; z \mid u) .\right.
\end{gathered}
$$

Then we have

$$
\text { iv) } \frac{\mu \mathrm{P}(t)-\mathrm{C}^{\prime}}{\mu \mathrm{P}(t)+\mathrm{C}^{\prime}} \leqslant e^{-\mathrm{C}^{\prime}\left(t-t_{0}\right)} \frac{\mu \mathrm{P}\left(t_{0}\right)-\mathrm{C}^{\prime}}{\mu \mathrm{P}\left(t_{0}\right)+\mathrm{C}^{\prime}} \quad \text { for } \quad t_{0} \leqslant t \leqslant t_{0}+h
$$

Proof. - The characteristic equations corresponding to (F) are the differential equations

$$
\begin{align*}
\dot{x}_{i}=f_{p_{i}}, \dot{p}_{i} & =-f_{x_{i}}-f_{u} p_{i} \equiv f_{i}  \tag{3}\\
\dot{v} & =\sum_{i} p_{i} f_{p_{i}}-f \equiv f_{0}, \quad i=1, \ldots, n
\end{align*}
$$

in which $=d / d t$, and $f$ and its derivatives have $(x, t, v, p)$ as argument. Selecting an arbitrary point $\xi=\left(\xi_{1}, \ldots, \xi_{n}\right)$, we impose the initial conditions :
$x_{i}=\xi_{i}, v=w(\xi), p_{i}=\partial w(\xi) / \partial \xi_{i} \quad$ for $\quad t=t_{0}, \quad i=1, \ldots, n$.
We shall see that the problem (3), (4) has a unique solution for $t_{0} \leqslant t \leqslant \mathrm{~T}$.

For us to be assured of a solution

$$
\begin{equation*}
x_{i}(t ; \xi), v(t ; \xi), p_{i}(t ; \xi), \quad i=1, \ldots, n \tag{5}
\end{equation*}
$$

of the initial value problem (3), (4) in the interval $t_{0} \leqslant t \leqslant \mathrm{~T}$, it will suffice to be able to find a priori bounds for $|v|$ and $|p|$ in any interval $t_{0} \leqslant t \leqslant t^{\prime}$ in which the solution (5) exists. For convenience, set $p_{0}=v$ and $\mathrm{P}=\left(p_{0}, p_{1}, \ldots, p_{n}\right)$, then writing the last $n+1$ characteristic equations in the form

$$
\dot{p}_{a}=f_{a}(x, t, \mathrm{P}), \quad a=0,1, \ldots, n
$$

Integrating and using (4) gives us the integral relations

$$
p_{a}(t ; \xi)=w_{a}(\xi)+\int_{t_{0}}^{t} f_{a}(x(s ; \xi), s, \mathrm{P}(s ; \xi)) d s, \quad a=0,1, \ldots, n
$$

supposedly valid for $t \leqslant t \leqslant t^{\prime}$; here,

$$
w_{0}=w \quad \text { and } \quad w_{i}=w_{x_{i}}, i=1, \ldots, n
$$

Defining

$$
\rho(t) \equiv \rho(t ; \xi)=\max _{a=0,1, \ldots, n} \sup _{t_{0}<s<t}\left|p_{a}(s ; \xi)\right|
$$

in view of (2) and of assumption (5), Section 1, we have

$$
\left|p_{a}(t ; \xi)\right| \leqslant \mathrm{W}+\int_{t_{0}}^{t} \mathrm{E}(\rho(s)) d s
$$

This implies that

$$
\rho(t) \leqslant \mathrm{W}+\int_{t_{0}}^{t} \mathrm{E}(\rho(s)) d s
$$

We conclude by Gronwall's reasoning that $\rho(t) \leqslant y\left(t ; t_{0}\right)$ in the interval $t_{0} \leqslant t \leqslant t^{\prime}$. Thus, all the dependent variables (5) have a priori bounds in this interval, and the bounds do not depend upon $t^{\prime}$. From well known results on the maximum extent of solutions of ordinary differential equations we can conclude that the problem (3), (4) has a unique solution in the interval $t_{0} \leqslant t \leqslant T$ and the estimates

$$
\begin{equation*}
|v(t ; \xi)| \leqslant y\left(t ; t_{0}\right),\left|p_{i}(t ; \xi)\right| \leqslant y\left(t ; t_{0}\right), \quad i=1, \ldots, n \tag{6}
\end{equation*}
$$

hold there.
Eventually, $v$ is identified with the solution $u$ of (F), (G) and $p_{i}$ with $\partial u / \partial x_{i}$ in a suitable layer $\mathrm{Z}\left(t_{0}, t_{0}+h\right)$. To estimate the height $h$ of this layer, we are led to consider new quantities $p_{i j}$, $i, j=1, \ldots, n$, which will ultimately be identified with the second derivatives $\partial^{2} u / \partial x_{i} \partial x_{j}$. These $p_{i j}$ satisfy ordinary differential equations obtained by applying $\partial^{2} / \partial x_{i} \partial x_{j}$ formally to equation ( $F$ ) and the substituting $p_{i j}$ for $\partial^{2} u / \partial x_{i} \partial x_{j}, p_{i}$ for $\partial u / \partial x_{i}$, and $\dot{p}_{i j} \equiv d p_{i j} / d t$ for $\sum_{k}\left(\partial p_{i j} / \partial x_{k}\right) f_{p_{k}}+\partial p_{i j} / \partial t$. These equations are

$$
\begin{aligned}
& \dot{p}_{i j}=\mathrm{C}_{i j}+\sum_{k} \mathrm{D}_{i k} p_{k j}+\sum_{k} \mathrm{D}_{i k} p_{k i}-f_{u} p_{i j} \\
&-\sum_{k, \ell} f_{p_{k} p_{\ell}} p_{i k} p_{j \ell}, \quad i, j=1, \ldots, n,(7)
\end{aligned}
$$

where

$$
\begin{gathered}
-\mathrm{C}_{i j}=f_{x_{i} x_{j}}+f_{u x_{i}} p_{j}+f_{u x_{j}} p_{i}+f_{u u} p_{i} p_{j} \\
-\mathrm{D}_{i k}=f_{x_{i} p_{k}}+p_{i} f_{\boldsymbol{p}_{k} u}
\end{gathered}
$$

In these equations, $p_{i}$ stands for $p_{i}(t ; \xi)$, and the derivatives of $f$ have the argument $(x(t ; \xi), t, v(t ; \xi), p(t ; \xi))$. We understand the $p_{i j}$ to be functions of $(t ; \xi)$ determined by (7) and the initial conditions

$$
\begin{equation*}
p_{i j}\left(t_{0} ; \xi\right)=w_{x_{i} x_{j}}(\xi) \tag{8}
\end{equation*}
$$

In view of (6) and assumption (7), Section 1, a constant $C$ independent of $M$ and $\mu$ exists such that the right side of (7) is not greater in absolute value than $\mathrm{C}\left(1+\sum_{k, \ell}\left|p_{k \ell}\right|\right)^{2}$. Hence, any solution of (7) defined, say, for $t_{0} \leqslant t \leqslant t^{\prime}$ will satisfy the inequalities

$$
\left|\dot{p}_{i j}\right| \leqslant \mathrm{C}\left(1+\sum_{k, \ell}\left|p_{k \ell}\right|\right)^{2}
$$

and by (8), on the other hand,

$$
\left|p_{i j}\left(t_{0} ; \xi\right)\right| \leqslant \mathrm{M}
$$

Summing over $i$ and $j$ and reasoning as in the proof of Gronwall's inequality shows that $\sum_{i, j}\left|p_{i j}\right| \leqslant q$, where $q \equiv q(t)$ satisfies the integral relation

$$
q(t)=n \mathrm{M}+n \mathrm{C} \int_{t_{0}}^{t}(1+q(s))^{2} d s
$$

and, equivalently, the differential equation $d q / d t=n \mathrm{C}(1+q)^{2}$ and the initial condition $q\left(t_{0}\right)=n \mathrm{M}$. The latter imply that

$$
q(t)=\frac{n \mathrm{M}+n \mathrm{C}(1+n \mathrm{M})\left(t-t_{0}\right)}{1-n \mathrm{C}(1+n \mathrm{M})\left(t-t_{0}\right)}
$$

if $t$ is such that the denominator is positive. Setting

$$
\begin{equation*}
h=1 /[2 n \mathrm{C}(1+n \mathrm{M})] \tag{9}
\end{equation*}
$$

we conclude, in particular, that
$\left|p_{i j}(t, \xi)\right| \leqslant 2 n \mathrm{M}+1 \quad$ if $\quad t_{0} \leqslant t \leqslant t_{0}+h, \quad i, j=1, \ldots, n$,
the right side of this inequality being the maximum of $q$ for the indicated domain of $t$.

In view of the a priori estimate (10) just achieved, standard methods now show that the initial value problem (7), (8) has a unique solution in the interval $t_{0} \leqslant t \leqslant t_{0}+h$; it is subject to (10).

The next step is to show that the mapping $x=x(t ; \xi)$ can be inverted to give $\xi$ as a function of $x, t$. To do so, choose $x_{1} \in \mathrm{E}^{n}$ and $t_{1} \geqslant \mathrm{~T}$ arbitrarily, and require that $\left(\xi, t_{0}\right) \in \mathrm{B} \equiv \mathrm{B}\left(x_{1}, t_{1} ; t_{0}\right)$. Let

$$
\begin{aligned}
\mathrm{Y}=(t ; \xi)=\left(t ; \xi_{1}, \ldots, \xi_{n}\right), \mathrm{X}=(x, t)=\left(x_{1}, \ldots, x_{n}, t\right) \\
\mathrm{X}(t ; \xi)=\left(x_{1}(t ; \xi), \ldots, x_{n}(t ; \xi), t\right)
\end{aligned}
$$

and define the set

$$
\mathrm{K}(h)=\left\{\mathrm{Y}=(t ; \xi):\left(\xi, t_{0}\right) \in \mathrm{B}, t_{0} \leqslant t \leqslant t_{0}+h\right\}
$$

In view of (2)", the transformation

$$
\begin{equation*}
X=X(Y) \tag{11}
\end{equation*}
$$

maps $\mathrm{K}(h)$ into $\mathrm{L}(h)$. Below we shall prove that the functional determinant of the transformation has a positive lower bound $\alpha$ :

$$
\begin{equation*}
\frac{\partial(\mathrm{X})}{\partial(\mathrm{Y})} \geqslant \alpha \quad \text { for } \quad \mathrm{Y} \in \mathrm{~K}(h) \tag{12}
\end{equation*}
$$

as will be seen, $\alpha$ can be selected to be independent of $M$. It will follow from (12) that the mapping (11) is locally one-to-one. Furthermore, it is easily seen that the range of the mapping is $L(h)$. Indeed, a characteristic curve issuing from an arbitrary point of $\mathrm{L}(h)$ has spatial coordinates satisfying the conditions $\dot{x}_{i}=f_{p_{i}}$, $i=1, \ldots, n$, and, owing to the definition of A, therefore will not escape from $\mathrm{L}(h)$ through the conical sides as the parameter $t$ decreases to the value $t_{0}$. Hence, the characteristics that emanate from
the points of B cover $\mathrm{L}(h)$, which means that (11) maps $\mathrm{K}(h)$ onto $\mathrm{L}(h)$, a simply connected region. Hence, this locally one-to-one mapping is globally one-to-one from $\mathrm{K}(h)$ onto $\mathrm{L}(h)$. Both the mapping and its inverse are continuous and have continuous first and second partial derivatives. This mapping (11) embraces the transformation

$$
\begin{equation*}
x=x(t ; \xi) \tag{13}
\end{equation*}
$$

which, for fixed $t$ in the interval $t_{0} \leqslant t \leqslant t_{0}+h$, maps a subset $\mathrm{S}(t)$ of

$$
\mathrm{S}=\left\{\xi:\left|\xi-x_{1}\right| \leqslant \mathrm{A}\left(t_{1}-t_{0}\right)\right\}
$$

onto

$$
\mathrm{S}_{t}=\left\{x:\left|x-x_{1}\right| \leqslant \mathrm{A}\left(t_{1}-t\right)\right\} .
$$

The inverse transformation to (11) embraces the inverse transformation to (13), which we write as

$$
\begin{equation*}
\xi=\xi(x, t), \tag{14}
\end{equation*}
$$

mapping $S_{t}$ onto $S(t)$. The previous remarks concerning (11) and its inverse imply that both (13) and (14) are continuous and have continuous first and second partial derivatives. Consequently, by traditional theory, the function

$$
u(x, t)=v(t, \xi(x, t))
$$

is a solution of $(\mathrm{F})$ in $\mathrm{L}(h)$, where also

$$
u_{x_{i}}(x, t)=p_{i}(t, \xi(x, t)), u_{x_{i} x_{j}}(x, t)=p_{i j}(t, \xi(x, t))
$$

for $i, j=1, \ldots, n$. Since $x_{1}$ and $t_{1}$ are arbitrary, the same is true in $\mathrm{Z}\left(t_{0}, t_{0}+h\right)$. Finally, $u$ satisfies the initial condition (G).

We have still to prove (12). The functional determinant on the left equals

$$
\mathrm{J} \equiv \mathrm{~J}(t) \equiv \mathrm{J}(\xi, t)=\frac{\partial(x(t ; \xi))}{\partial(\xi)}
$$

the functional determinant for (13). By (4), $\mathrm{J}\left(t_{0}\right)=1$, and by continuity a positive number $h^{\prime}$ exists such that $\mathrm{J}(\xi, t)>0$ in $\mathrm{K}\left(h^{\prime}\right)$. Denoting by $h_{1}$ the supremum of all $h^{\prime}$ such that $\mathbf{J}(\xi, t)>0$ in $\mathrm{K}\left(h^{\prime}\right)$ and $h^{\prime} \leqslant h$, we shall prove that $h_{1}=h$. For any $h^{\prime}$ as described, the argument of the previous paragraph shows that the mapping (13) has
an inverse (14) for which $u(x, t)=v(t ; \xi(x, t)$ solves (F), (G), in $\mathrm{L}\left(h^{\prime}\right)$ and $u_{x_{i}}=p_{i}, u_{x_{i} x_{j}}=p_{i j}$ there. Hence, in $L\left(h^{\prime}\right)$ we have the differentiation formulas

$$
\begin{gathered}
v_{\xi_{j}} \equiv \partial v / \partial \xi_{j}=\sum_{k} \frac{\partial u}{\partial x_{k}} \frac{\partial x_{k}}{\partial \xi_{j}}=\sum_{k} p_{k} x_{k, \xi_{j}} \\
p_{i, \xi_{j}} \equiv \partial p_{i} / \partial \xi_{j}=\sum_{k} \frac{\partial^{2} u}{\partial x_{i} \partial x_{k}} \frac{\partial x_{k}}{\partial \xi_{j}}=\sum_{k} p_{i k} x_{k, \xi_{j}}
\end{gathered}
$$

where $x_{k, \xi_{j}}=\partial x_{k} / \partial \xi_{j}$. Using these when differentiating with respect to $\xi_{j}$ in the first $n$ characteristic equations (3) gives us

$$
\dot{x}_{i, \xi_{j}}=\sum_{k} f_{p_{i} x_{k}} x_{k, \xi_{j}}+f_{p_{i} u} \sum_{k} p_{k} x_{k, \xi_{j}}+\sum_{k} f_{p_{i} p_{k}} \sum_{\ell} p_{k \ell} x_{\ell, \xi_{j}}
$$

Hence, differentiating $\mathbf{J}$ by rows and cancelling as is appropriate, we obtain

$$
\mathrm{j}=\sum_{i}\left(f_{p_{i} x_{i}}+f_{p_{i} u} p_{i}+\sum_{k} p_{i k} f_{p_{i} p_{k}}\right) \mathrm{J} .
$$

Previous estimations show that the quantities in parentheses are bounded in absolute value by an expression of the form $(a+b \mathrm{M}) / n$, where $a$ and $b$ are constants independant of $M$ and $\mu$. Therefore,

$$
\mathrm{j} \geqslant-(a+b \mathrm{M}) \mathrm{J}
$$

and since $\mathrm{J}\left(\xi, t_{0}\right)=1$, we can deduce that

$$
\mathrm{J}(\xi, t) \geqslant e^{-(a+b \mathrm{M})\left(t-t_{0}\right)} \quad \text { in } \quad \mathrm{K}(h)
$$

In $\mathrm{K}(h)$, we have $0 \leqslant t-t_{0} \leqslant h$, while by (9) $h=1 /[2 n \mathrm{C}(1+n \mathrm{M})]$. We conclude that a positive constant $\alpha$ independent of M and $\mu$ exists such that $\mathrm{J} \geqslant \alpha$ in $\mathrm{K}(h)$ and therefore such that (12) holds, as required. We have thus established statements (i), (ii), (iii) of the theorem being considered.

To prove (iv), let $z=\left(z_{1}, \ldots, z_{n}\right)$ be an arbitrary unit vector : $|z|=1$. Multiplying both sides of equation (7) by $z_{i} z_{j}$ and summing gives

$$
\begin{align*}
& \dot{\mathrm{P}}^{*}=\sum_{i, j} \dot{p}_{i j} z_{i} z_{j}=\sum_{i, j} \mathrm{C}_{i j} z_{i} z_{j}+2 \sum_{k} \sum_{i} z_{i} \mathrm{D}_{i k} \sum_{j} p_{k j} z_{j} \\
&-f_{u} \mathrm{P}^{*}-\sum_{k, \ell} f_{p_{k} p_{\ell}}\left(\sum_{i} p_{i k} z_{i}\right)\left(\sum_{j} p_{j \ell} z_{\jmath}\right) \tag{15}
\end{align*}
$$

where $\mathrm{P}^{*}=\sum_{i, j} p_{i j} z_{i} z_{j}$. Constants $c^{\prime}$ and $c^{\prime \prime}$ independent of M and $\mu$ exist such that $\sum_{i, j} \mathrm{C}_{i j}^{2} \leqslant\left(c^{\prime}\right)^{2}, f_{u}^{2} \leqslant c^{\prime \prime}$, and $\sum_{i, k} \mathrm{D}_{i k}^{2} \leqslant c^{\prime \prime}$, the last condition implying

$$
2\left|\sum_{k} \sum_{i} z_{i} \mathrm{D}_{i k} \sum_{j} p_{k j} z_{j}\right| \leqslant 4 c^{\prime \prime} / \mu+(\mu / 4) \sum_{k}\left(\sum_{j} p_{k j} z_{j}\right)^{2} .
$$

In view of assumption (2), Section 1, we also have

$$
\sum_{k, \ell} f_{p_{k} p \ell}\left(\sum_{i} p_{i k} z_{i}\right)\left(\sum_{j} p_{j \ell} z_{j}\right) \geqslant \mu \sum_{k}\left(\sum_{j} p_{k j} z_{j}\right)^{2} .
$$

Hence, (15) leads to the inequality

$$
\dot{\mathrm{P}}^{*} \leqslant c / \mu+(\mu / 4) \mathrm{P}^{* 2}-(3 \mu / 4) \sum_{k}\left(\sum_{i} p_{k j} z_{j}\right)^{2}
$$

where $c$ is a bound, for instance, for $c^{\prime} \mu+5 c^{\prime \prime}$ and is taken to be independent of $\mu$ for $\mu$ less than an agreed upper bound. Moreover, $\sum_{k}\left(\sum_{i} p_{i k} z_{i}\right)^{2}=\left[\sum_{k} z_{k}^{2}\right]\left[\sum_{k}\left(\sum_{i} p_{i k} z_{i}\right)^{2}\right]$

$$
\geqslant\left[\sum_{k} z_{k} \sum_{i} p_{i k} z_{i}\right]^{2}=\mathrm{P}^{* 2}
$$

It follows that

$$
\begin{equation*}
\mathrm{P}^{*} \leqslant c / \mu-(\mu / 2) \mathrm{P}^{* 2} \tag{16}
\end{equation*}
$$

In order to discuss this, set

$$
s=\sqrt{c / 2 t}, s_{0}=\sqrt{c / 2} t_{0}, q=(\mu / \sqrt{2 c}) \mathrm{P}^{*}, q_{0}=q\left(s_{0}\right),
$$

inequality (16) becoming

$$
\begin{equation*}
d q / d s \leqslant 1-q^{2} . \tag{17}
\end{equation*}
$$

We shall prove from it that if

$$
\mathrm{Q}=\max (1, q), \mathrm{Q}_{0}=\max \left(1, q_{0}\right),
$$

then

$$
\begin{equation*}
\frac{\mathrm{Q}-1}{\mathrm{Q}+1} \leqslant \frac{\mathrm{Q}_{0}-1}{\mathrm{Q}_{0}+1} e^{-2\left(s-s_{0}\right)} . \tag{18}
\end{equation*}
$$

Four cases arise.
Case 1. $\left|q_{0}\right|<1$. In this case, (17) leads to $d q /\left(1-q^{2}\right) \leqslant d s$ and, by integration, to the inequality

$$
\frac{1+q}{1-q} \leqslant \frac{1+q_{0}}{1-q_{0}} \zeta,
$$

where $\zeta=e^{2\left(s-s_{0}\right)}$. This implies that

$$
q \leqslant \frac{\zeta-1+(\zeta+1) q_{0}}{\zeta+1+(\zeta-1) q_{0}}
$$

and thus that $q \leqslant 1$ in this case.
Case 2. $q_{0}<-1$. As $s$ increases, $q$ achieves no value exceeding $q_{0}$, in view of the sign of $d q / d s$. Hence, $q<-1$ in this case.

Case 3. $q_{0}>1$. Here, we have $d q /\left(q^{2}-1\right) \leqslant-d s$, and therefore

$$
\frac{q-1}{q+1} \leqslant \frac{q_{0}-1}{q_{0}+1} \zeta^{-1}
$$

or, equivalently,

$$
\begin{equation*}
q \leqslant \frac{(\zeta+1) q_{0}+\zeta-1}{(\zeta-1) q_{0}+\zeta+1} \tag{19}
\end{equation*}
$$

Case 4. $\left|q_{0}\right|=1$. Either $q \leqslant 1$ or else $q$ assumes a value $q_{1}>1$. Letting $s_{1}$ be the first value above $s_{0}$ for which this occurs, we have by Case 3 that for $s \geqslant s_{1}$

$$
q \leqslant \frac{\left(\zeta_{1}+1\right) q_{1}+\zeta_{1}-1}{\left(\zeta_{1}-1\right) q_{1}+\zeta_{1}+1},
$$

where $\zeta_{1}=\exp \left[2\left(s-s_{1}\right)\right]$. Since we may select $q_{1}$ arbitrarily close to 1 , we conclude that in fact $q \leqslant 1$ for $s \geqslant s_{0}$.

Except in Case 3, we have just seen that $q \leqslant 1$. Hence, in all four cases,

$$
q \leqslant \frac{(\zeta+1) \mathrm{Q}_{0}+\zeta-1}{(\zeta-1) \mathrm{Q}_{0}+\zeta+1}
$$

the right side of this inequality exceeding the right side of (19) and also exceeding 1 . For the last reason,

$$
\mathrm{Q} \leqslant \frac{(\zeta+1) \mathrm{Q}_{0}+\zeta-1}{(\zeta-1) \mathrm{Q}_{0}+\zeta+1} .
$$

Equivalently,

$$
\frac{Q-1}{Q+1} \leqslant \frac{Q_{0}-1}{Q_{0}+1} \zeta^{-1}
$$

which implies (iv) with $\mathrm{C}^{\prime}=\sqrt{2 c}$. The theorem is now completely proved.

## 3. Stratified solutions of (F), (G).

To construct stratified solutions of this problem we employ smoothing operators $S_{\epsilon}$ with the properties that for an arbitrary function $f \in \mathrm{C}^{2}\left(\mathrm{E}^{n}\right)$ :

$$
\begin{gather*}
\left|\mathrm{S}_{\epsilon} f\right|_{0} \leqslant|f|_{0},  \tag{1}\\
\left|\left(\mathrm{~S}_{\epsilon} f\right)_{x_{i} x_{j}}\right|_{0} \leqslant k_{1} \varepsilon^{-1}|f|_{1} \quad \text { for } \quad i, j=1, \ldots, n,  \tag{2}\\
\mathrm{P}^{*}\left(t ; z \mid \mathrm{S}_{\epsilon} f\right) \leqslant \mathrm{P}^{*}(t ; z \mid f) \quad \text { for all unit vectors } z \in \mathrm{E}^{n},  \tag{3}\\
\left|\int\left(\mathrm{~S}_{\epsilon} f-f\right) \phi d x\right| \leqslant \mathrm{C}_{\phi}|f|_{0} a(\varepsilon) \quad \text { for any } \quad \phi \in \mathrm{C}_{0}^{2}\left(\mathrm{E}^{n}\right) . \tag{4}
\end{gather*}
$$

Here, as on previous occasions, $\mathrm{C}^{2}\left(\mathrm{E}^{n}\right)$ denotes the class of bounded functions with bounded, continuous derivatives of first and second orders in $\mathrm{E}^{n} ;|f|_{0}=\sup _{x \in \mathrm{E}^{n}}|f(x)| ;|f|_{1}=\max _{i=1, \ldots, n} \sup _{\mathrm{E}^{n}}\left|f_{x_{i}}\right| ; k_{1}$ is an absolute constant $; \mathrm{C}_{\phi}$ is a constant depending on $\phi$ only $; \mathrm{P}^{*}$ is the functional defined in the theorem of Section $2 ; \lim _{\epsilon \rightarrow 0} a(\varepsilon) / \varepsilon=0$. Averaging operators, for instance, will satisfy all these conditions, see Chapter 3.

Let us fix once and for all the thickness $T$ of the layer $Z(T)$ within which we are to solve the given problem. Thereby the constants $C$ and $C^{\prime}$ in the theorem of Section 2 are determined. We shall construct a stratified solution of (F), (G) with stratum thickness $h$ such that $\mathrm{C} h \leqslant 1 / 2$ and that, for convenience, $\mathrm{T} / h$ is an integer. We determine $\varepsilon$ by the condition

$$
h=1 /\left\{\mathrm{C}\left[1+k_{1} y(\mathrm{~T}) / \varepsilon\right]\right\}
$$

and thus as

$$
\varepsilon=\frac{k_{1} y(\mathrm{~T}) \mathrm{C} h}{1-\mathrm{C} h}
$$

where $y(t)$ is the increasing nonnegative function defined by the relation $\int_{U}^{y(t)} d \rho / E(\rho)=t$. In view of (2) and part (i) in the theorem of the previous section, an inductive procedure analogous to that of Section 3, Chapter 2, shows that, for $m=1,2, \ldots, T / h$, a strict solution $u_{m}$ of (F) in $Z_{m}\left({ }^{1}\right)$ exists that satisfies the initial condition

$$
u_{m}(x,(m-1) h)=\mathrm{S}_{\epsilon} u_{m-1}(x,(m-1) h)
$$

and obeys the conditions

$$
\left|u_{m}\right|_{\mathrm{c}^{0}\left(\mathrm{z}_{m}\right)} \leqslant y(m h),\left|u_{m, x_{i}}\right|_{\mathrm{C}^{0}\left(\mathrm{z}_{m}\right)} \leqslant y(m h) \text { for } i=1, \ldots, n,
$$

$u_{0}(x, 0)$ standing for $u_{0}(x)$. Thus, the stratified solution

$$
v=u_{m} \quad \text { in } \quad \mathrm{Z}_{m}^{\prime} \quad\left({ }^{1}\right) \text { for } m=1, \ldots, \mathrm{~T} / h
$$

satisfies the inequalitics

$$
\begin{equation*}
|\nu|_{\mathrm{C}^{0}(\mathrm{Z}(\mathrm{~T}))} \leqslant y(\mathrm{~T}),\left|v_{x_{i}}\right|_{\mathrm{C}^{0}(\mathrm{Z}(\mathrm{~T}))} \leqslant y(\mathrm{~T}) . \tag{5}
\end{equation*}
$$

Condition (iv) of the theorem referred to tells us that for

$$
\begin{gather*}
(m-1) h \leqslant t \leqslant m h \\
\frac{\mu \mathrm{P}(t)-\mathrm{C}^{\prime}}{\mu \mathrm{P}(t)+\mathrm{C}^{\prime}} \leqslant e^{-\mathrm{C}^{\prime}(t-(m-1) h)} \frac{\mu \mathrm{P}_{m-1}-\mathrm{C}^{\prime}}{\mu \mathrm{P}_{m-1}+\mathrm{C}^{\prime}} \tag{6}
\end{gather*}
$$

where $\mathrm{P}(t)=\mathrm{P}\left(t ; z \mid u_{m}\right), \mathrm{P}_{k}=\mathrm{P}\left(k h ; z \mid u_{k}\right)$, and also gives us the relations
$\left({ }^{1}\right)$ The notation is that of Chapter 2, Section 1 (p. 44).

$$
\begin{aligned}
& \frac{\mu \mathrm{P}_{k+1}-\mathrm{C}^{\prime}}{\mu \mathrm{P}_{k+1}+\mathrm{C}^{\prime}}=\frac{\mu \mathrm{P}\left((k+1) h ; z \mid u_{k+1}\right)-\mathrm{C}^{\prime}}{\mu \mathrm{P}\left((k+1) h ; z \mid u_{k+1}\right)+\mathrm{C}^{\prime}} \\
& \quad \leqslant e^{-\mathrm{C}^{\prime} h \frac{\mu \mathrm{P}\left(k h ; z \mid u_{k+1}\right)-\mathrm{C}^{\prime}}{\mu \mathrm{P}\left(k h ; z \mid u_{k+1}\right)+\mathrm{C}^{\prime}}}
\end{aligned}
$$

for $k=0,1, \ldots, \mathrm{~T} / h-1$.
But inequality (3) implies that $\mathrm{P}\left(k h ; z \mid u_{k+1}\right) \leqslant \mathrm{P}\left(k h ; z \mid u_{k}\right)=\mathrm{P}_{k}$. Hence, the last member of the previous inequality is

$$
\leqslant e^{-\mathrm{C}^{t h}}\left[\mu \mathrm{P}_{k}-\mathrm{C}^{\prime}\right] /\left[\mu \mathrm{P}_{k}+\mathrm{C}^{\prime}\right]
$$

and in sum we have the recursion

$$
\frac{\mu \mathrm{P}_{k+1}-\mathrm{C}^{\prime}}{\mu \mathrm{P}_{k+1}+\mathrm{C}^{\prime}} \leqslant e^{-\mathrm{C}^{\prime} h} \frac{\mu \mathrm{P}_{k}-\mathrm{C}^{\prime}}{\mu \mathrm{P}_{k}+\mathrm{C}^{\prime}} \text { for } k=1,1, \ldots, \mathrm{~T} / h-1 .
$$

Multiplying these for $k=0,1, \ldots, m-2$ and then also multiplying by (6) shows that

$$
\frac{\mu \mathrm{P}(t)-\mathrm{C}^{\prime}}{\mu \mathrm{P}(t)+\mathrm{C}^{\prime}} \leqslant e^{-\mathrm{C}^{\prime} t} \frac{\mu \mathrm{P}_{0}-\mathrm{C}^{\prime}}{\mu \mathrm{P}_{0}+\mathrm{C}^{\prime}} \leqslant e^{-\dot{c}^{\prime} t} .
$$

from which we immediately deduce that

$$
\mathrm{P}(t) \leqslant\left(\mathrm{C}^{\prime} / \mu\right) \operatorname{coth}\left(\mathrm{C}^{\prime} t / 2\right)
$$

It follows that for an arbitrary unit vector $z \in \mathrm{E}^{n},|z|=1$, we have

$$
\sum_{i, j} v_{x_{i} x_{j}} z_{i} z_{j} \leqslant\left(\mathrm{C}^{\prime} / \mu\right) \operatorname{coth}\left(\mathrm{C}^{\prime} t / 2\right) \quad \text { in } \quad \mathrm{Z}(\mathrm{~T})
$$

and thus

$$
\begin{aligned}
v(x & +\rho z, t)+v(x-\rho z, t)+2 v(x, t) \\
& =\int_{0}^{1}\left[\int_{-1}^{1} \rho^{2} z_{i} z_{j} v_{x_{i} x_{j}}\left(x+\theta \theta^{\prime} \rho z\right) d \theta^{\prime}\right] \theta d \theta \\
& \leqslant\left(\mathrm{C}^{\prime} / \mu\right) \rho^{2} \operatorname{coth}\left(\mathrm{C}^{\prime} t / 2\right) \text { for }(x, t) \in \mathrm{Z}(\mathrm{~T}), z \in \mathrm{E}^{n},|z|=1 .(7)
\end{aligned}
$$

## 4. Precompactness of stratified solutions.

We shall demonstrate the following result, similar in its nature and its proof to the precompactness theorem of Section 4, Chapter 2. Again $v_{\boldsymbol{h}}$ denotes the stratified solution with layer thickness $h$.

Precompactness theorem. A bounded, Lipschitz continuous function $u$ in $Z(T)$ and a null sequence $\left\{h_{k}\right\}$ exist such that for any $t$ in the interval $0 \leqslant t<\mathrm{T}$

$$
\lim _{k \rightarrow \infty} v_{h_{k}}(\cdot, t)=u(\cdot, t)
$$

in $\mathrm{E}^{n}$, uniformly in any bounded region of $\mathrm{E}^{n}$.
Again the proof is in several steps. First, we have by methods of the earlier section :

Proposition 1. - A null sequence $\left\{h_{k}\right\}$ exists such that for any bounded, measurable function $\phi$ with compact support in $\mathrm{E}^{n}$,

$$
\lim _{k \rightarrow \infty} \int_{\mathrm{E}^{n}} v_{h_{k}}(x, t) \phi(x) d x
$$

exists for $0 \leqslant t<\mathrm{T}$.
Secondly, we prove :

Proposition 2. - Let $t$ be a fixed value in ( $0, \mathrm{~T}$ ), and let $\left\{h_{k}\right\}$ be an arbitrary subsequence of the sequence $\left\{h_{k}\right\}$ in Proposition 1. Correspondingly, a function $u^{\prime}(\cdot, t)$ continuous in $\mathrm{E}^{n}$ and a subsubsequence $\left\{h_{k^{\prime \prime}}\right\}$-a subsequence of $\left\{h_{k^{\prime}}\right\}$-exist such that

$$
\lim _{k^{\prime \prime} \rightarrow \infty} v_{h_{k^{\prime \prime}}}(\cdot, t)=u^{\prime}(\cdot, t)
$$

in $\mathrm{E}^{n}$, uniformly in every bounded region of $\mathrm{E}^{n}$.
The proof is immediate from the equiboundedness and equicontinuity of the $\nu_{h}(\cdot, t)$ (inequalities (5), Section 3).

Propositions 1 and 2 imply that for any bounded measurable $\phi$ with compact support in $E^{n}$,
$\lim _{k \rightarrow \infty} \int_{\mathrm{E}^{n}} \nu_{h_{k}}(x, t) \phi(x) d x$

$$
\lim _{k^{\prime \prime} \rightarrow \infty} \int_{\mathrm{E}^{n^{\prime}}} v_{h_{k^{\prime \prime}}}(x, t) \phi(x) d x=\int_{\mathrm{E}^{n}} u^{\prime}(x, t) \phi(x) d x
$$

which in turn implies that $u^{\prime}(\cdot, t)$ is independent of the subsequence $\left\{h_{k^{\prime}}\right\}$. This and Proposition 2 give us:

Proposition 3. - For each fixed value of $t$ in ( $0, \mathrm{~T}$ ), a function $u(\cdot, t)$ continuous in $\mathrm{E}^{n}$ exists such that

$$
\lim _{k \rightarrow \infty} v_{h_{k}}(\cdot, t)=u(\cdot, t)
$$

in $\mathrm{E}^{n}$, uniformly in any bounded region of $\mathrm{E}^{n}$.
Next, we shall prove that $u$ is Lipschitz continuous in $\mathrm{Z}(\mathrm{T})$ and begin by discussing its continuity with respect to $t$. For this purpose, we return to the functions

$$
\mathrm{V}(t) \equiv \mathrm{V}\left(t ; v_{h}\right)=\int v_{h}(x, t) \phi(x) d x
$$

with $\phi \in \mathrm{C}_{0}^{2}\left(\mathrm{E}^{n}\right)$ and $\phi \geqslant 0$. The domain of integration is $\mathrm{E}^{n}$. Since the $v_{h}$ are stratified solutions, the one-sided limits $\mathrm{V}(k h \pm 0)$ exist, and for $(k-1) h<t<k h$ the derivative

$$
\mathrm{V}^{\prime}(t)=\int v_{h, t}(x, t) \phi(x) d x
$$

exists, with $k=1,2, \ldots, \mathrm{~T} / h$. If $0 \leqslant t_{1} \leqslant t_{2}<\mathrm{T}$ and, more specifically,

$$
m_{1} h \leqslant t_{1}<\left(m_{1}+1\right) h \leqslant m_{2} h \leqslant t_{2}<\left(m_{2}+1\right) h
$$

we thus have

$$
\begin{aligned}
\mathrm{V}\left(t_{2}\right)-\mathrm{V}\left(t_{1}\right) & =\left\{\int_{t_{1}}^{\left(m_{1}+1\right) h}+\sum_{k=m_{1}+1}^{m_{2}-1} \int_{k h}^{(k+1) h}+\int_{m_{2} h}^{t_{2}}\right\} \mathrm{V}^{\prime}(t) d t \\
& +\sum_{k=m_{1}+1}^{m_{2}}[\mathrm{~V}(k h-0)-\mathrm{V}(k h+0)]
\end{aligned}
$$

Property (4), Section 3, shows the second summation on the right, in absolute value, to be $\leqslant \mathrm{C}_{\phi}\left(t_{2}-t_{1}\right) h^{-1} y(\mathrm{~T}) a(\varepsilon)$ and thus to
tend towards zero with $h$. The remaining part of the right hand side, in absolute value, is $\leqslant \mathrm{F}_{0}\left(t_{2}-t_{1}\right) \int \phi(x) d x$, where $\mathrm{F}_{0}$ is an upper bound for $|f(x, t, u, p)|$ for $x \in \mathrm{E}^{n}, 0 \leqslant t \leqslant \mathrm{~T},|u| \leqslant y(\mathrm{~T})$, $|p| \leqslant y(\mathrm{~T})$. Consequently, we have

$$
\begin{aligned}
\lim _{k \rightarrow \infty}\left|\mathrm{~V}\left(t_{2} ; v_{h_{k}}\right)-\mathrm{V}\left(t_{1} ; v_{n_{k}}\right)\right| & =\left|\int\left(u\left(x, t_{2}\right)-u\left(x, t_{1}\right)\right) \phi(x) d x\right| \\
& \leqslant \mathrm{F}_{0}\left(t_{2}-t_{1}\right) \int \phi(x) d x
\end{aligned}
$$

This implies that

$$
\int \phi(x)\left\{\mathrm{F}_{0}\left(t_{2}-t_{1}\right) \pm\left(u\left(x, t_{2}\right)-u\left(x, t_{1}\right)\right)\right\} d x \geqslant 0
$$

and since $u$ is continuous with respect to $x$ and $\phi$ is arbitrary,

$$
\mathrm{F}_{0}\left(t_{2}-t_{1}\right) \pm\left(u\left(x, t_{2}\right)-u\left(x, t_{1}\right)\right) \geqslant 0
$$

This says that $u$ is Lipschitz continuous with respect to $t$ with Lipschitz constant $\mathrm{F}_{0}$. It follows that $u$ is Lipschitz continuous with respect to $(x, t)$.

Properties (5) and (7) of stratified solutions, Section 3, carry over to $u$, and we have

$$
|u(x, t)| \leqslant y(\mathrm{~T}),|u(x+z, t)-u(x, t)| \leqslant y(\mathrm{~T})|z|,
$$

$u(x+z, t)+u(x-z, t)-2 u(x, t) \leqslant\left(\mathrm{C}^{\prime} / \mu\right)|z|^{2} \quad \operatorname{coth}\left(\mathrm{C}^{\prime} t / 2\right)(2)$ for $(x, t) \in Z(0, T)$ and $z \in \mathrm{E}^{n}$.

## 5. Generalized solutions of (F), (G).

The Lipschitz continuous function $u$ obtained in the previous section obviously fulfills the initial condition (G). Now we shall show that $u$ satisfies the differential equation (F) almost everywhere in $Z(T)$.

By the previous section,

$$
\lim _{m \rightarrow \infty} u^{(m)}=u \quad \text { in } \quad \mathrm{Z}(\mathrm{~T})
$$

where $u^{(m)} \equiv u_{h_{m}}, m=1,2, \ldots$, is a suitable sequence of stratified solutions as described. Since $u^{(m)}$ is a stratified solution, we have for $0<t<\mathrm{T}$ and, say, $k h_{m} \leqslant t<(k+1) h_{m}$,

$$
\begin{aligned}
& u^{(m)}(x, t)-u^{(m)}\left(x, k h_{m}+0\right) \\
&=\int_{k h_{m}}^{t} f\left(x, s, u^{(m)}(x, s), \operatorname{grad} u^{(m)}(x, s)\right) d s
\end{aligned}
$$

also,
$u^{(m)}\left(x, j h_{m}-0\right)-u^{(m)}\left(x,(j-1) h_{m}+0\right)$

$$
=\int_{(j-1) h_{m}}^{i h_{m}} f\left(x, s, u^{(m)}(x, s), \operatorname{grad} u^{(m)}(x, s)\right) d s
$$

for $j=1, \ldots, k$.
Adding these $k+1$ relations, multiplying by an arbitrary test function $\phi \in \mathrm{C}_{\mathbf{0}}^{2}(\mathrm{Z}(\mathrm{T})$ ), and integrating over $\mathrm{Z}(\mathrm{T})$ gives us

$$
\begin{align*}
\int_{Z(\mathrm{~T})} \phi(x, t) & \left\{u^{(m)}(x, t)-u^{(m)}(x, 0)\right. \\
& -\int_{0}^{t} f\left(x, s, u^{(m)}(x, s), \operatorname{grad} u^{(m)}(x, s) d s\right\} d x d t  \tag{1}\\
& =\sum_{j=1}^{k} \int \phi\left(x, j h_{m}\right)\left[u^{(m)}\left(x, j h_{m}+0\right)-u^{(m)}\left(x, j h_{m}-0\right)\right] d x
\end{align*}
$$

We shall wish to let $m \rightarrow \infty$ in this formula and therefore now consider the convergence of the first derivatives

$$
p_{i}^{(m)} \equiv \partial u^{(m)} / \partial x_{i}
$$

of the stratified solutions of the sequence to the respective first derivatives

$$
p_{i} \equiv \partial u / \partial x_{i}
$$

of their limit.
By Lipschitz continuity, the $p_{i}$ exist on a set $G \subset Z(T)$ such that $\mathrm{Z}(\mathrm{T}) \backslash \mathrm{G}$ has measure zero. We shall now argue that the $p_{i}$ are the respective limits almost everywhere of the $p_{i}^{(m)}$ as $m \rightarrow \infty$.

Fix the index $i$. By Fubini's theorem, $G$ includes almost all points of almost all lines in $\mathrm{Z}(\mathrm{T})$ that are parallel to the $x_{i}$-axis.

On such a line L, almost all points of which belong to $G$, we regard $x_{i}$ as the only variable in the functions

$$
q^{(m)}\left(x_{i}\right) \equiv p_{i}^{(m)}(x, t), q\left(x_{i}\right) \equiv p_{i}(x, t)
$$

and consider the other $x_{j}$ and $t$ to be fixed parameters. For $t>0$, it follows from the last inequality (2), Section 4, that $q$ satisfies the one-sided Lipschitz condition

$$
\begin{equation*}
\frac{q\left(x_{i}\right)-q\left(x_{i}^{\prime}\right)}{x_{i}-x_{i}^{\prime}} \leqslant \kappa \quad \text { for } \quad x_{i} \neq x_{i}^{\prime} \tag{2}
\end{equation*}
$$

with $\kappa \equiv \kappa(t)=\left(\mathrm{C}^{\prime} / \mu\right) \operatorname{coth}\left(\mathrm{C}^{\prime} t / 2\right)$; hence, $q$ is continuous at almost all points on L. (See a remark on "semi-monotonic" functions in [6], page 78). Inequality (7), Section 3, implies that the $q^{(m)}$ for $m=1,2, \ldots$ also satisfy a one-sided Lipschitz condition of the form (2) and all with the same constant $\kappa$ independent of $m$. Consequently, from any subsequence $\left\{m^{\prime}\right\}$ of indices, a sub-subsequence $\left\{m^{\prime \prime}\right\}$ can be selected such that $q^{\left(m^{\prime \prime}\right)}$ converges almost everywhere on L. (See [6]). The limit function $q^{\prime}$ again is subject to the one-sided Lipschitz condition (2) and therefore is continuous at almost all points of L . Furthermore, $q^{\left(m^{\prime \prime}\right)} \rightarrow q^{\prime}$ at every point at which $q^{\prime}$ is continuous.

The next step is to identify $q^{\prime}$ with $q$ at the points at which $q$ is continuous. Let $L^{\prime}$ consist of the points of $L$ at which $q$ is continuous and $L^{\prime \prime}$ consist of the $i$ - th coordinates of these points. Setting

$$
\mathrm{U}^{(m)}\left(x_{i}\right) \equiv u^{(m)}(x), \mathrm{U}\left(x_{i}\right) \equiv u(x) \quad \text { along } \mathrm{L}
$$

we have for $x_{i}, a_{i} \in \mathrm{~L}^{\prime \prime}$

$$
\mathrm{U}^{(m)}\left(x_{i}\right)=\mathrm{U}^{(m)}\left(a_{i}\right)+\int_{a_{i}}^{x_{i}} q^{(m)}(\xi) d \xi
$$

Letting $m=m^{\prime \prime} \rightarrow \infty$ gives

$$
\mathrm{U}\left(x_{i}\right)=\mathrm{U}\left(a_{i}\right)+\int_{a_{i}}^{x_{i}} q^{\prime}(\xi) d \xi
$$

while the same relation holds with $q$ in place of $q^{\prime}$. Therefore, $q^{\prime}=q$ almost everywhere on $L$ and, in particular, at every point of continuity of the two functions. This means that $\lim _{m^{\prime \prime} \rightarrow \infty} q^{\left(m^{\prime \prime}\right)}=q$ at every
point of continuity of $q$ on $L$, while the limit $q$ is independent of the arbitrarily selected subsequence $\left\{m^{\prime}\right\}$. With respect to the original sequence, we conclude, therefore, that

$$
\lim _{m \rightarrow \infty} q^{(m)}=q
$$

at every point at which $q$ is continuous on $L$. It follows that

$$
\lim _{m \rightarrow \infty} p_{i}^{(m)}=p_{i}
$$

almost everywhere in $\mathrm{Z}(\mathrm{T})$. Let $p_{i}^{*}$ and $p_{i^{*}}$ denote the upper and lower limits of the $p_{i}^{(m)}$, respectively, as $m \rightarrow \infty$. The set of points S at which $p_{i}^{*}-p_{i^{*}}>0$ is measurable and is included in the set of discontinuities of $p_{i}$ with respect to $x_{i}$. Hence, for almost every line $L$ parallel to the $x_{i}$-axis the intersection $\mathrm{S} \cap \mathrm{L}$ has one-dimensional measure 0 . Therefore, the $(n+1)$-dimensional measure of $S$ is 0 , i.e., $\lim _{m \rightarrow \infty} p_{i}^{(m)}$ exists almost everywhere in $Z(T)$. Similar reasoning, in which $p^{*}-p_{i}$ is considered, will show that $\lim _{m \rightarrow \infty} p_{i}^{(m)}=p_{i}$ almost everywhere.

It is now possible to let $m \rightarrow \infty$ in (1). The limit of the right hand side is 0 in view of property (4), Section 3, and we obtain :

$$
\begin{aligned}
\int_{\mathrm{Z}(\mathrm{~T})} \phi(x, t)\left\{\begin{array}{l}
u(x, t)-u_{0}(x) \\
\\
\end{array} \quad-\int_{0}^{t} f(x, s, u(x, s), \operatorname{grad} u(x, s)) d s\right\} d x d t=0
\end{aligned}
$$

Since $\phi$ is arbitrary, we can infer that the quantity in curly brackets is 0 almost everywhere in $\mathrm{Z}(\mathrm{T})$ and therefore, in particular, at almost all points of almost all line segments $x=$ const., $0 \leqslant t<\mathrm{T}$, parallel to the $t$-axis. If M is such a segment, by differentiation with respect to $t$ we have $u_{t}=f$ almost everywhere on M. Since $u_{t}$ and $f$ are measurable, this relation, which is (F), holds almost everywhere in $\mathrm{Z}(\mathrm{T})$, as was to be proved.

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$$

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