# Weak and Strong Extensions of First-Order Differential Operators in $R^{m}$ 

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## 1. Introduction

Let $L$ be the differential operator

$$
\begin{equation*}
L=\sum_{k=1}^{m} A_{k k}(x) D^{k}+B(x) \tag{1.1}
\end{equation*}
$$

where $x=\left(x_{1}, \ldots, x_{m}\right), D^{k}=\delta / \partial x_{k}, A_{1}, \ldots, A_{m}$ are $\nu \times \nu$ continuously differentiable matrices in $R^{m}$, and $B$ is a continuous $v \times v$ matrix; naturally, $L$ acts on $\nu$-dimensional vector functions. (We assume that all entities involved are real valued, although the results hold equally well for complex-valued functions and matrices.) The minimal operator $A_{0}$ associated with $L$ is defined by $A_{0} u=L u$ with domain $D\left(A_{0}\right)=\mathscr{D}^{\prime \prime}(\mathscr{D}=$ the space of Schwartz test functions in $R^{m}$ ), while the maximal operator $A$ associated with $L$ (or the weak extension of $A_{0}$ in $\left.L^{p}\left(R^{m}\right)^{v}=L^{p}\left(R^{m}, R^{v}\right)(1<p<\infty)\right)$ is the adjoint $\left(A_{0}^{\prime}\right)^{*}$ of $A \mathrm{t}$, the formal adjoint of $A_{0}$, in $L^{p^{\prime}}\left(R^{m}\right)^{\nu}, p^{\prime-1}+p^{-1}=1\left(A_{0}^{\prime}\right.$ is defined as the minimal operator associated with the formal adjoint $L^{\prime} u=-\sum D^{k}\left(A_{k}^{*} u\right)+$ $B^{*} u$, where ${ }^{*}$ indicates matrix adjoint). Finally, the strong extension of $A_{0}$ in $L^{p}\left(R^{m}\right)^{\nu}$ is $C l\left(A_{0}\right)$. A classical result of Friedrichs [1] establishes that the weak and strong extensions of $L$ coincide, i.e.,

$$
\begin{equation*}
C l\left(A_{0}\right)=A \tag{1.2}
\end{equation*}
$$

for $p=2$ assuming that $A_{k}, D^{j} A_{k}$, and $B$ are all bounded in $R^{m}$, and the proof is easily seen to work just as well for $1<p<\infty$. Although this result has been extended in many directions (see, for example [3]) I have been unable to find in the literature any results on whether (1.2) holds under weaker assumptions on $A_{k}, B$ at infinity. The object of this note is to show that the boundedness

[^0]conditions on the $D^{j} A_{k}$ and $B$ can be entirely discarded; as for the $A_{k}$ we need only assume that
\[

$$
\begin{equation*}
\left|A_{k}(x)\right| \leqslant \rho(|x|) \quad\left(x \in R^{i n}\right) \tag{1.3}
\end{equation*}
$$

\]

for $k=1,2, \ldots, m$, where $\rho$ is a positive continuous function growing so slowly at infinity that

$$
\begin{equation*}
\int^{\infty} \frac{d r}{\rho(r)}=\infty \tag{1.4}
\end{equation*}
$$

The proof of this result (Theorem 2.1) is no more than a simple modification of Friedrichs' method of mollifiers via a "reduction to compact support' that has been employed many times (see, for instance, [6]). What makes condition (1.3)(1.4) interesting, however, is that it is best possible in the following sense: Given a positive nondecreasing continuous $\rho$ which does not satisfy (1.4) there exists an $L$ (with $m=2, v=1$ ) satisfying (1.3) but not (1.2). It is not by chance, incidentally, that (1.4) is Wintner's condition in [9] for existence in the large of solutions of the ordinary differential system

$$
\begin{equation*}
X^{\prime}(t)=a(X(t)) \tag{1.5}
\end{equation*}
$$

where $a=\left(a_{1}, \ldots, a_{m}\right)$, is a suitably smooth map from $R^{m}$ into itself satisfying $|a(x)| \leqslant \rho(|x|)$, since there exists a close relation among (1.5) and the hyperbolic equation

$$
D_{t} u=a_{1} D^{1} u+\cdots+a_{m} D^{m} u
$$

namely, solutions of (1.6) are constant along curves ( $X(t), t)$ in ( $m+1$ )dimensional space. This relation is in fact used in the construction of the counterexample (Theorem 3.1).

We point in Section 4 an application of Theorem 2.1 to symmetric hyperbolic systems.

## 2. The Main Result

Theorem 2.1. Let $L$ be as in Section 1, and assume the coefficients $A_{1}, \ldots, A_{m}$ satisfy (1.3)-(1.4). Then (1.2) holds, i.e., Cl $A_{0}=A$ for $1<p<\infty$.

Proof. We recall briefly the definition and some properties of Friedrichs' mollifiers [1,2,4]. Let $\psi$ be a nonnegative function in $\mathscr{D}$ with integral 1 and vanishing (say) in $|x| \geqslant 1$; the operators $J_{n}(n \geqslant 1)$ are defined by $J_{n}=\psi_{n} *$, where $\psi_{n}(x)=n^{m} \psi(n x)$. Each $J_{n}$ is a bounded operator in $L^{p}\left(R^{m}\right)^{\nu}(1 \leqslant p \leqslant \infty)$ with norm $\leqslant 1$; if $1 \leqslant p<\infty, J_{n} u \rightarrow u$ in the $L^{p}$-norm as $n \rightarrow \infty$ (see [4] for proofs). In the following result the only hypotheses on $A_{1}, \ldots, A_{m}, B$ are those at the beginning of Section 1; i.e., no growth conditions at infinity are prescribed.

Lemma 2.2. Let $u \in D(A)$ have compact support. Then

$$
\begin{equation*}
\left\|J_{n} A u-A J_{n} u\right\| \rightarrow 0 \quad(n \rightarrow \infty) \tag{2.1}
\end{equation*}
$$

Proof. Lemma 2.2 is a well-known result of Friedrichs [1] (see also [4]) when $A_{k}, D^{i} A_{k}$, and $B$ are bounded in $R^{m}$, with no restrictions on the support of $u$. We only have to apply this theorem to the operator $\tilde{L}=\sum \widetilde{A}_{k} D^{k}+\widetilde{B}$, where the $\tilde{A}_{k}$ (resp. $\tilde{B}$ ) are, say, continuously differentiable (resp. continuous) matrix functions defined in all of $R^{n t}$, having compact support and coinciding with $A_{k}($ resp. $B)$ in the set of all those $x \in R^{m}$ with $\operatorname{dist}(x, K) \geqslant 1$; in fact, if $\tilde{A}$ is the maximal operator of $\tilde{L}$ we obviously have $A J_{n} u=\tilde{A} J_{n} u, J_{n} A u=J_{n} \tilde{A} u$.

Proof of Theorem 2.1. We may obviously assume that $\rho$ is infinitely differentiable. Let $\varphi$ be another infinitely differentiable function of $r$ with $\varphi(r)=0$ if $r \leqslant 0, \varphi(r)=1$ if $r \geqslant 1$. Given $0<r<\infty$ define

$$
\begin{array}{rlrl}
f_{1}(s) & =1-\int_{r}^{s} \frac{d \sigma}{\rho(\sigma)} & &  \tag{2.2}\\
& \text { if } & 0 \leqslant s \leqslant r & r \leqslant s \leqslant s_{r} \\
& =0 & \text { if } \quad s_{r} \leqslant s
\end{array}
$$

where $s_{r}$ is such that $\int_{r}^{s_{r}} \rho(s)^{-1} d s=1$. The function

$$
\begin{equation*}
\chi_{r}(x)=\varphi\left(f_{r}(|x|)\right) \tag{2.3}
\end{equation*}
$$

belongs to $\mathscr{D}$, its support being contained in $|x| \leqslant S_{r}$; its first partials have support in the ring $r \leqslant|x| \leqslant s_{r}$ and

$$
\begin{equation*}
\left|D^{k} \chi_{r}(x)\right| \leqslant C / \rho(|x|) \quad\left(x \in R^{n n}\right) \tag{2.4}
\end{equation*}
$$

If $u \in L^{p}\left(R^{m}\right)^{v}$ it is clear that

$$
\begin{equation*}
\chi_{r} u \rightarrow u \tag{2.5}
\end{equation*}
$$

as $r \rightarrow \infty$. On the other hand, if $u \in D(A)$ a simple computation with adjoints shows that

$$
\begin{equation*}
A\left(\chi_{r} u\right)=\left(\sum_{k=1}^{m} A_{k} D^{k} \chi_{r}\right) u+\chi_{r} A u \tag{2.6}
\end{equation*}
$$

In view of (1.3) and (2.4) we obtain, taking into account that the $D^{k} \chi_{r}$ vanish for $|x| \leqslant r$,

$$
\|\left(A_{k} D^{k} \chi_{r}\right) u_{\|}^{| |} \leqslant C \int_{|x| \geqslant r}|u|^{p} d x
$$

which tends to zero as $r \rightarrow \infty$; accordingly,

$$
A\left(x_{r} u\right)->A u
$$

which, combined with (2.5), shows that $C l\left(A_{C}\right)=A, A_{C}$ the restriction of $A$ to elements of $D(A)$ having compact support. Hence we only have to prove that, given $u \in D\left(A_{C}\right)$ there exists a sequence $\left\{u_{n}\right\}$ in $D\left(A_{0}\right)$ with $u_{n} \rightarrow u, A_{0} u_{n} \rightarrow A u$. In view of (2.2), $\left\{J_{n} u\right\}$ is such a sequence. This ends the proof.

## 3. A Counterexample

Throughout this section $\rho$ is a positive, nondecreasing differentiable function in $r \geqslant 0$ with $\rho^{\prime}(0)=0$, and
and

$$
\begin{equation*}
\int_{0}^{x} \frac{d r}{\rho(r)}<\infty \tag{3.1}
\end{equation*}
$$

$$
\begin{equation*}
L=\rho\left(x_{2}\right) D^{1}+\rho\left(x_{1}\right) D^{2} \tag{3.2}
\end{equation*}
$$

where $\rho$ is extended to $r<0$ by setting $\rho(r)=\rho(0)$. The operator $L$ is a particular case of that in (1.1) (here $m=2, \nu=1$ ).

Theorem 3.1. Let $A_{0}$ (resp. A) be the minimal (resp. maximal) operator associated with L. Then, for $p=2$

$$
\begin{equation*}
C l\left(A_{0}\right) \neq A \tag{3.3}
\end{equation*}
$$

Proof. Since $D^{1} \rho\left(x_{2}\right)+D^{2} \rho\left(x_{1}\right)=0$, the operator $A_{0}^{\prime}=-A_{0}$ is dissipative (in fact, conservative, since $\left(A_{0}^{\prime} w, w\right)=0$ for $w \in D\left(A_{0}^{\prime}\right)$ ) thus it follows from standard duality arguments ([4]) that

$$
(I \pm A) D(A)=L^{2}\left(R^{2}\right)
$$

If $C l\left(A_{0}\right)=A$ then $A$ itself is conservative; hence it generates a group $S(t)$, $-\infty<t<\infty$ of isometric operators in $L^{2}\left(R^{2}\right)$. Consider now the ordinary differential system

$$
\begin{equation*}
X_{1}^{\prime}(t)=\rho\left(X_{2}(t)\right), \quad X_{2}^{\prime}(t)=\rho\left(X_{1}(t)\right) \tag{3.4}
\end{equation*}
$$

Its trajectories can be described as follows. Solutions starting (say, at $t=0$ ) in the negative quadrant $\pi_{-}\left(x_{1}, x_{2}<0\right)$ are straight lines

$$
X_{1}(t)=X_{1}(0)+t \rho(0), \quad X_{2}(t)=X_{2}(0)+t \rho(0)
$$

until they leave $\pi_{--}$. If $\left(X_{1}(0), X_{2}(0)\right) \in \pi_{+-}\left(x_{1} \geqslant 0, x_{2}<0\right)$, the corresponding trajectory is

$$
\begin{aligned}
& X_{1}(t)=X_{1}(0)+t \rho(0) \\
& X_{2}(t)=X_{2}(0)+\int_{0}^{t} \rho\left(X_{1}(0)+s \rho(0)\right) d s
\end{aligned}
$$

until it leaves $\pi_{+-}$(which it will eventually do, since its slope is $>1$ ). An entirely symmetric expression holds in $\pi_{-+}$. Finally, if ( $\left.X_{1}(0), X_{2}(0)\right)$ belongs to the positive quadrant $\pi_{++}, X_{1}(t)$ is obtained inverting the function

$$
\begin{equation*}
t\left(X_{1}\right)=\int_{0}^{X_{1}} \frac{d r}{\rho\left(R^{-1}\left(R(r)+C_{1}\right)\right.}+C_{2} \tag{3.5}
\end{equation*}
$$

and setting

$$
X_{2}(t)-R^{-1}\left(R\left(X_{1}(t)\right)+C_{1}\right)
$$

where

$$
R(r)=\int_{0}^{r} \rho(s) d s \quad(s \geqslant 0)
$$

$R^{-1}$ is the functional inverse of $R$, and $C_{1}, C_{2}$ are constants determined by the initial conditions

$$
\begin{aligned}
& C_{1}=R\left(X_{2}(0)\right)-R\left(X_{1}(0)\right) \\
& C_{2}=-\int_{0}^{X_{1}(0)} \frac{d r}{\rho\left(R^{-1}\left(R(r)+C_{1}\right)\right)} .
\end{aligned}
$$

If $C_{1} \geqslant 0, R^{-1}\left(R(r)+C_{1}\right) \geqslant r$, and thus $\rho\left(R^{-1}\left(R(r)+C_{1}\right)\right) \geqslant \rho(r)$. Hence, $t$ tends to a finite limit as $X_{1}$ (thus $X_{2}$ ) tends to infinity; consequently, solutions with $X_{1}(0) \geqslant X_{2}(0)$ cease to exist at a time $t=t\left(X_{1}(0), X_{2}(0)\right)$ given by

$$
\begin{equation*}
t=\int_{X_{1}(0)}^{\infty} \frac{d r}{\rho\left(R^{-1}\left(R(r)+C_{1}\right)\right)} \leqslant \int_{X_{1}(0)}^{\infty} \frac{d r}{\rho(r)} . \tag{3.6}
\end{equation*}
$$

(Since we can argue in the same way interchanging $X_{1}$ and $X_{2}$, all solutions starting in $\pi_{++}$have a finite life expectancy.) Putting together all these observations we conclude that an arbitrary trajectory of (2.4) will eventually enter $\pi_{++}$ and it will blow itself up in finite time; on the hand, if time is run backward, the trajectory will enter $\pi_{-}$and will then exist forever (see Fig. 1).

Let $T>0, u_{0}$ be a vector function in $\mathscr{P}^{\nu}$. Given $\left(x_{1}, x_{2}\right) \in R^{2}$ denote by $X_{1}\left(x_{1}, x_{2}, t\right), X_{2}\left(x_{1}, x_{2}, t\right)$ the solution of (3.5) with initial conditions

$$
\begin{equation*}
X_{1}\left(x_{1}, x_{2}, 0\right)=x_{1}, \quad X_{2}\left(x_{1}, x_{2}, t\right)=x_{2} \tag{3.7}
\end{equation*}
$$

and let $u\left(x_{1}, x_{2}, t\right)$ be defined by

$$
u\left(x_{1}, x_{2}, t\right)=u_{0}\left(X_{1}\left(x_{1}, X_{2}, T-t\right), X_{2}\left(x_{1}, x_{2}, T-t\right)\right)
$$

for $\left(x_{1}, x_{2}\right) \in R^{2}, 0 \leqslant t \leqslant T\left(u\left(x_{1}, x_{2}, t\right)=0\right.$ if $\left(X_{1}\left(x_{1}, x_{2}, s\right), X\left(x_{1}, x_{2}, s\right)\right)$ blows up before $T-t$ ). It follows from standard theorems on dependence on initial data (see [5]) and from our previous observations on the trajectories of


Figure 1
(3.4) that $u$ is continuously differentiable in $R^{2} \times[0, T]$ and has compact support there. It is obvious that $u$ is constant on trajectories of (3.4), thus $D_{t} u=\rho\left(x_{2}\right) D^{1} u+\rho\left(x_{1}\right) D^{2} u$, and it is easy to see that the $L^{2}\left(R^{2}\right)$-valued function $u(t)=u(\cdot, \cdot, t)$ is a solution of $u^{\prime}(t)=A u(t)$ in $0 \leqslant t \leqslant T$ with $u(T)=u_{0}$; thus

$$
\begin{equation*}
u(0)=S(-T) u_{0} \tag{3.8}
\end{equation*}
$$

It follows from (3.6) and the comments preceding and following this inequality that

$$
\begin{equation*}
u\left(x_{1}, x_{2}, 0\right)=0 \quad\left(\left(x_{1}, x_{2}\right) \in \Omega\right) \tag{3.9}
\end{equation*}
$$

where $\Omega$ is the quadrant $x_{1}, x_{2} \geqslant \omega$ with $\omega$ so large that

$$
\int_{\omega}^{\infty} \frac{d r}{\rho(r)} \leqslant T
$$

We end the proof as follows. Let $u$ be an arbitrary element of $L^{2}\left(R^{2}\right)$; choose a sequence $\left\{u_{n}\right\}$ in $\mathscr{D}^{v}$ with $u_{n} \rightarrow u$, write (3.8) for $u_{n}$, and take limits. We obtain using (3.9) that $S(-T) u=0$ a.e. in $\Omega$, which contradicts the fact that $S(\cdot)$ is a group. We must then conclude that (3.3) holds.

## 4. An Application

Consider the symmetric hyperbolic system

$$
\begin{equation*}
D_{t} u=\sum_{k=1}^{m} A_{k} D^{k} u+B u \tag{4.1}
\end{equation*}
$$

where $A_{1}, \ldots, A_{m}, B$ satisfy the smoothness assumptions at the beginning of Section 2; moreover, each $A_{k}$ is symmetric. Under the customary assumption that $A_{k}, D^{k} A_{k}$, and $B$ are bounded in $R^{m}$ it is known that the Cauchy problem for (4.1) is properly posed in $L^{2}\left(R^{m}\right)^{\nu}$; precisely, $A-\omega I$ is maximal dissipative there for some $\omega$. This is proved by showing that $A_{0}-\omega I$ is dissipative (by integration by parts) and extending this property to $A-\omega I$ using (1.2). We can then obtain the same conclusion under weaker hypotheses, namely,

$$
\begin{equation*}
B(x)-\frac{1}{2} \sum_{k=1}^{m} D^{k} A_{k}(x) \leqslant \omega I \quad\left(x \in R^{m}\right) \tag{4.2}
\end{equation*}
$$

which suffices to ensure dissipativity of $A_{0}-\omega I$ and (1.3)-(1.4), which yields (1.2) via Theorem 2.1.

The example in Theorem 3.1 has some intriguing features. It is not difficult to see that the operator $A$ there, although not a group generator, is maximal dissipative, and thus generates a contraction semigroup $S(t)(t \geqslant 0)$ in $L^{2}\left(R^{2}\right)$ which is explicitly given by the formula

$$
S(t) u\left(x_{1}, x_{2}\right)=u\left(X_{1}\left(x_{1}, x_{2},-t\right), X_{2}\left(x_{1}, x_{2},-t\right)\right), \quad\left(x_{1}, x_{2}\right) \in R^{2}
$$

However, the equation $u^{\prime}(t)=A u(t)$ does not possess the "finite-domain-ofdependence" properties usually expected of hyperbolic equations; in other words, perturbations arbitrarily far away in space at time $t=0$ may reach a given point all at (or near) the same time $T>0$. To see this, let $\left(X_{1}(t), X_{2}(t)\right)$ be an arbitrary trajectory of (3.4), $T$ its escape time (we may assume, translating time if necessary that $T>0$ ), and $\varphi$ a function in $\mathscr{D}$ with $\varphi(0) \neq 0$. If

$$
u_{n}\left(x_{1}, x_{2}\right)=\varphi\left(x_{1}-X_{1}(T-1 / n), x_{2}-X_{2}(T-1 / n)\right)
$$

the supports of the $u_{n}$ wander off to infinity; however, $S(t) u_{n}\left(x_{1}, x_{2}\right)=$ $\varphi\left(x_{1}-X_{1}(-1 / n), x_{2}-X_{2}(-1 / n)\right)$, whose support lies arbitrarily near that of $\varphi$.

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