# On Difference Schemes for Hyperbolic Equations with Discontinuous Initial Values 

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1. Introduction. In this paper we consider the hyperbolic equation

$$
\begin{align*}
\frac{\partial u}{\partial t} & =\rho(x, t) \frac{\partial u}{\partial x}+\phi(x, t) u \quad \text { with initial values }  \tag{1.1}\\
u(x, 0) & =f(x)
\end{align*}
$$

in the upper halfplane. It is well known that if $f(x)$ is a sufficiently smooth function (1.1) can be closely approximated by stable difference schemes, and very realistic error bounds are given in a number of papers (see, e.g., Lax [6]). But in applications there often arise initial functions, with simple discontinuities or discontinuities in the higher derivatives. A discontinuity in $f(x)$ propagates in the solution to (1.1) along a characteristic, but this propagation is disturbed in the solution of the corresponding difference equation. It is the aim of this paper to give rather sharp bounds for the error and to show how these bounds depend on the order of accuracy and dissipation of the difference scheme at hand.

We show that the error between the solutions of (1.1) and the corresponding difference approximation has two components $E_{1}$ and $E_{2}$, where $E_{1}=O\left(h^{p}\right), p$ is the order of accuracy. The other component $E_{2}=O\left(e^{-a}\right)$, where $a=N^{q} d(x, t)$ with $0<q<1$. Here $d(x, t)$ is the distance (cf. Section 2) from the characteristic to (1.1) through the jump-point. $h=N^{-1}$ is the mesh-size and $q$ depends on $p$ and the order of dissipation. Furthermore, for constant coefficients we can take $E_{1} \equiv 0$. In Section 2 we list a number of results. The first one shows that if $p$ is odd, then a dissipative scheme is necessarily dissipative of order $p+1$. But for even $p$, the dissipative schemes with largest $q$ are those with dissipation of order $p+2$. In Section 5, where we list the results of a number of numerical experiments, Experiment 9 shows that a nondissipative scheme can behave very badly; e.g., if $\rho \equiv 1$, $\phi \equiv 0$ and $N=200$, then the error for $t=1$, is bigger than 0.01 in an interval of length 2.

The main idea of this paper, due to H.-O. Kreiss, is to extend the real variable $\xi$ to the complex variable $\xi+i \alpha$ in the symbol $\hat{Q}(\xi)$ of the difference operator at hand. This technique is also used in a paper of H.-O. Kreiss and E. Lundqvist [5], where they consider a mixed problem for a hyperbolic equation.

The main part of our proofs, collected in Section 3, consists of algebraic manipulations of the transformed symbol $\hat{Q}(\xi+i \alpha)$, which lead to the definition of contractivity (Definition 2.4). Contraction can be considered as an algebraic condition on the difference scheme. Now, $\hat{Q}(\xi+i \alpha)$ is the symbol of a difference operator that originates from an equation obtained by a simple transformation of the original difference equation. This fact is used in Section 4, where the rest of the proofs are found. For the variable coefficients case, we also use a theorem of Lax-Nirenberg [7].

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In a forthcoming paper we intend to generalize our results to symmetric hyperbolic systems with variable coefficients.

There are some earlier results which ought to be mentioned. First, H.-O. Kreiss [4] has shown that the same relation between $p$ and the order of dissipation $2 s=p$ +1 or $2 s=p+2$ guarantees L.R.-stability for difference approximations to hyperbolic systems in any number of independent variables. (L.R.-stable means stable in the sense of Lax and Richtmyer.) B. Parlett has also discussed stability of dissipative schemes in [8], with the same relations between $p$ and $2 s$. We should also mention papers of V. Thomée [11], J. Peetré and V. Thomée [9], G. Hedström [3] and R. P. Fedorenko [2], which consider problems with discontinuous initial functions.

In [11] V. Thomée has proved a theorem like ours, but in his proof he uses a quite different technique containing Fourier-Stieltjes transforms of functions of bounded variation. His results are restricted to difference operators, which are stable in maximum-norm. But many of the difference schemes used in practice do not belong to this class; e.g., every scheme with even order of accuracy. Furthermore, nothing is said about equations with variable coefficients.

In [9], J. Peetré and V. Thomée give estimates for the rate of convergence of a very general class of difference schemes. They use for their proof the theory of interpolation of Banach spaces. In our special case they give an $L_{2}$-estimate of the error of order $h^{a}$, where $a=\frac{1}{2} p /(p+1)$. These estimates are better than the corresponding $L_{2}$-estimates using our results. However, an $L_{2}$-estimate in the problem of this paper gives quite misleading information, because the local error outside an interval of length $O\left(h^{q}|\log h|\right)$ is only of $O\left(h^{p}\right)$.

In [3] G. W. Hedström gives an estimate for the constant coefficient case. His results give $E_{1}=O\left(h^{p}\right)$ and $E_{2}=O\left(e^{-a}\right)$, where $a=$ const $N \cdot d(x, t)^{(p+1) / p}$ for $d(x, t) \geqq$ const $N^{-p /(p+1)}$. These results are therefore more precise* in this part of the $x$-axis, but are less general. In his proof he uses estimates of the norms of powers of absolutely convergent Fourier series. Generalizations to variable coefficients can hardly be treated with this method.

Finally, R. P. Fedorenko [2] describes a method that gives very good numerical results in a number of tested examples. The paper is of experimental character and contains no proofs. He starts with a high order method, say $p=3$, but when he comes to the discontinuity he changes the method to $p=1$ in order to avoid the parasitic waves, "the Gibbs phenomenon," and afterwards, when the danger is over, he again takes the first method with $p=3$. In order that the computer should recognize the discontinuities, he lets it compare the second and first differences in every mesh-point. This works well for simple examples like $\partial u / \partial t=\partial u / \partial x$, but to generalize to more complicated equations without this a priori knowledge about the characteristics seems to be hard.

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[^0]2. Preliminaries and Statement of Results. Consider the hyperbolic equation
\[

$$
\begin{equation*}
\frac{\partial u}{\partial t}(x, t)=\rho(x, t) \frac{\partial u}{\partial x}(x, t)+\phi(x, t) u(x, t), \tag{2.1}
\end{equation*}
$$

\]

where $\rho$ and $\phi$ are real valued functions, $-\infty<x<\infty$ and $t \geqq 0$. Its (generalized) solution is uniquely determined by initial values

$$
\begin{equation*}
u(x, 0)=f(x) \tag{2.2}
\end{equation*}
$$

We want to study the behavior of difference approximations to this problem. We therefore introduce a time-step $k>0$, a mesh-width $h$, and grid-points $x_{\nu}=\nu h$, $\nu=0, \pm 1, \pm 2, \cdots$. Furthermore let $k / h=\lambda, \lambda>0$ constant. Denoting by $v_{\nu}(t)$ $=v\left(x_{\nu}, t\right)$ a function defined for all $t=t_{n}=n k, n=0,1,2, \cdots$ we approximate (2.1), (2.2) either by

$$
\begin{equation*}
Q_{1} v_{\nu}(t+k)=Q_{2} v_{\nu}(t) \tag{2.3}
\end{equation*}
$$

or by

$$
\begin{equation*}
v_{\nu}(t+k)=Q_{0} v_{\nu}(t) \tag{2.3'}
\end{equation*}
$$

with initial conditions in either case

$$
\begin{equation*}
v_{\nu}(0)=f\left(x_{\nu}\right) . \tag{2.4}
\end{equation*}
$$

$Q_{j}, j=0,1,2$ are difference operators of the form

$$
\begin{aligned}
& Q_{1}(x, t, h)=\sum_{j=-l}^{m} b_{j}(x, t, h) E^{j} \\
& Q_{2}(x, t, h)=\sum_{j=-l}^{m} a_{j}(x, t, h) E^{j}
\end{aligned}
$$

and

$$
Q_{0}(x, t, h)=\sum_{j=-\infty}^{\infty} d_{j}(x, t, h) E^{j}
$$

respectively, where $E$ is the translation operator defined by

$$
E g_{\nu}=g_{\nu+1}
$$

We always assume that $a_{j}, b_{j}$ and $d_{j}$ are smooth functions of $h$. To get an algebraic description of the behavior of the solution of (2.3) or (2.3') we introduce the symbols

$$
\hat{Q}_{j}(x, t, \xi, h)=e^{-i \omega x} Q_{j} e^{i \omega x}, \quad j=0,1,2 .
$$

We always assume that for the implicit scheme (2.3)

$$
\left|\hat{Q}_{\mathbf{1}}(x, t, \xi, h)\right| \geqq \text { const }
$$

for all $x, t, \xi$ and for sufficiently small $h$.
To be complete, we have assumed $h$-dependent coefficients in (2.3) and (2.3'), but the essential features of a difference scheme come from the principal part of the difference operator, i.e. when $h=0$. From the assumption above, we see that $Q_{1}{ }^{-1}(x, t, 0)$ is a bounded operator in $L_{2}$ and hence

$$
\hat{Q}_{-1}(x, t, \xi, 0)=\hat{Q}_{1}^{-1}(x, t, \xi, 0) \hat{Q}_{2}(x, t, \xi, 0)
$$

is well defined.
Furthermore, we define the solution operator $S$ of (2.3)-(2.4) or (2.3')-(2.4) by

$$
v_{\nu}(t)=S(t, h) f\left(x_{\nu}\right) .
$$

We also introduce a discrete $L_{2}$-norm $\|\cdot\|_{h}$ by the scalar-product

$$
(f, g)_{h}=\sum_{\nu=-\infty}^{\infty} h f_{\nu} g_{\nu}, \quad\|f\|_{h}^{2}=(f, f)_{h}
$$

where $f$ and $g$ are mesh-functions.
Some definitions are now in order.
Definition 2.1. The difference scheme (2.3) or $\left(2.3^{\prime}\right)$ is stable if the operator $S(t, h)$ is uniformly bounded in $L_{2}$ for $0 \leqq t \leqq T$, independent of $h$.

Definition 2.2. The difference scheme (2.3) or (2.3') is accurate of order $p$, (consistent if $p \geqq 1$ ), if $p$ is the largest integer such that for all sufficiently differentiable solutions $u(x, t)$ of (2.1)

$$
u(x, t+k)=Q_{i}(x, t, h) u(x, t)+O\left(h^{p+1}\right), \quad h \rightarrow 0, i=0,-1
$$

It is well known that a scheme (2.3) or (2.3') which is accurate of order $p$ satisfies

$$
\begin{equation*}
\hat{Q}_{j}(x, t, \xi, 0)=e^{i \rho(x, t) \lambda \xi}+c_{0}(x, t) \xi^{p+1}(1+o(1)), \quad \xi \rightarrow 0 \tag{2.5}
\end{equation*}
$$

(cf. Lax [6]), $j=0,-1$.
Definition 2.3. A difference scheme defined by the operator $Q(x, t, h)$ is dissipative of order $2 s$ if $s$ is the least integer such that for all $x$ and $t$, and for all $\xi$ with $|\xi| \leqq \pi$,

$$
|\hat{Q}(x, t, \xi, 0)| \leqq 1-\delta|\xi|^{2 s},
$$

## $\delta>0$ some constant.

We now introduce another algebraic condition in terms of which we will state our main theorems.

Definition 2.4. Let $\alpha$ be a real constant. A difference scheme defined by the operator $Q(x, t, h)$ is contractive of order $\tau=\left(\tau_{-}, \tau_{+}\right)$if there is a uniformly bounded function $\sigma(x, t)$, independent of $\xi$, such that, for $\alpha$ sufficiently small,

$$
\hat{Q}(x, t, \xi+i \alpha, 0)=\exp \left(-\alpha \lambda \rho(x, t)+\sigma(x, t)|\alpha|^{\tau}\right) \cdot R(x, t, \xi), \quad|\xi| \leqq \pi
$$

with $\tau=\tau_{-}$for $\alpha \rho<0$ and $\tau=\tau_{+}$for $\alpha \rho>0$ and where $R(x, t, \xi)$ is such that $|R| \leqq 1$ for all $x$ and $t$, and for all $\xi$ with $|\xi| \leqq \pi$.

Let $r=2 s-p$; then it is easy to see that there are no dissipative difference schemes with $r \leqq 0$, and the following theorem says that for odd $p$ we can only have $r=1$.

Theorem 2.1. A dissipative difference scheme of odd order of accuracy $p$ is dissipative of order $p+1$.

But for even order of accuracy $p$, we can construct schemes with $r=2 j, j=1$, $2, \cdots$. (Section 5, Experiment 11 gives an example.)

There are connections between $\tau, p$ and $s$ for a given scheme, which we are going to describe in the next theorem.

Theorem 2.2. Suppose that the difference scheme (2.3) or (2.3') is
(i) accurate of order $p$,
(ii) dissipative of order $2 s$, and that
(iii) $\rho(x, t)$ is uniformly bounded and does not change sign; then, for $p=2 j-1$, $j=1,2, \cdots$, it is contractive of order $\tau=(p+1, p+1)$. Moreover, for $p=2 j$, $j=1,2, \cdots$, if
(iv) $i c_{0}(x, t) \rho(x, t)>0$
is satisfied, then it is contractive of order $\tau=(6 s /(r+2), 2 s / r) .\left(c_{0}(x, t)\right.$ was defined in (2.5).)

We believe condition (iv) is satisfied for all schemes used in practice. It is satisfied for explicit schemes (2.3) of maximum accuracy. (Cf. Strang [10] and our report [1].)

Theorem 2.3. Consider the difference scheme (2.3) with $Q_{1}=I$ (identity operator) and $Q_{2}=\sum_{j=-N}^{N} a_{j} E^{j}$ of maximum accuracy, i.e., $p=2 N$. For these difference schemes condition (iv) of Theorem 2.2 is satisfied, if $|\rho \lambda| \leqq 1$.

The proofs of these three theorems will be given in Section 3. They are purely algebraic in nature. The next three theorems have analytical proofs, which are presented in Section 4.

In this paper we only consider initial functions $f(x)$, which are step-functions

$$
\begin{align*}
f(x) & =0 & & x<0  \tag{2.6}\\
& =1 & & x \geqq 0 .
\end{align*}
$$

More general discontinuous initial functions could be divided into a sum of a smooth function and step-functions, which are handled in our original report [1]. That this works is essentially due to the linear character of our problem. For smooth initial functions we can refer to [6].

The error estimates in Theorems 2.4-2.5 show that the influence to the error from the discontinuity-jump is exponentially decreasing with the distance $d(x, t)$ from the characteristic through the origin.

Definition 2.5. The distance $d(x, t)$ from the characteristic through the origin is given by $d(x, t)=|g(x, t)|$, where $g(x, t)$ satisfies $\partial g / \partial t=\rho(x, t) \partial g / \partial x$ with $g(x, 0)$ $=x$.

We first study the constant coefficient case. Then $d(x, t)=|\rho t+x|$. It is no restriction to assume that $\phi=0$ in (2.1). Furthermore, let $T_{+1}$ denote the region in the ( $x, t$ )-plane, where $0 \leqq t \leqq T$ and where $(x, t)$ lies to the right of the characteristic of (2.1) through the origin. $T_{-1}$ is the corresponding region to the left of the same characteristic.

Theorem 2.4. Let $\rho \neq 0$ be a constant and let $\phi \equiv 0$. Consider the difference scheme (2.3) with coefficients independent of $x, t$ and $h$, and let $f(x)$ be defined by (2.6). Suppose that
(i) (2.3) is consistent with (2.1),
(ii) (2.3) is contractive of order $\tau=\left(\tau_{-}, \tau_{+}\right)$. Then for $h$ sufficiently small,

$$
\begin{equation*}
\left|u\left(x_{\nu}, t\right)-v_{\nu}(t)\right| \leqq c h^{(q-1) / 2} \exp \left(-h^{-q}\left|\rho t+x_{\nu}\right|\right), \tag{2.7}
\end{equation*}
$$

where $c$ is a constant. Here $q=\left(\tau_{\mp}-1\right) / \tau_{\mp}$, with minus sign in $T_{-\operatorname{sign}(\rho)}$ and plus sign in $T_{+ \text {sign }(\rho)}$.

We now turn to the variable coefficient case, and in order to be able to use a stability theorem of Lax and Nirenberg (see [7] and Section 4) we choose the scheme (2.3'). This theorem also needs some smoothness properties on the coefficients $d_{j}$ in
(2.3'). To simplify our proof, we assume that $\rho(x, t)$ and $\phi(x, t)$ have compact support. Here $\|\cdot\|$ means the ordinary $L_{2}$-norm.

Theorem 2.5. Consider the difference scheme (2.3') with initial values defined by (2.6). Suppose that
(i) $\left(2.3^{\prime}\right)$ is accurate of order $p$,
(ii) (2.3') is contractive of order $\tau=\left(\tau_{-}, \tau_{+}\right)$,
(iii) $\rho(x, t)$ and $\phi(x, t) \in c_{0}{ }^{\infty},{ }^{* *}$ and $\rho(x, t)$ does not change sign.

$$
\begin{align*}
& \sum_{j} \max \left\{\left\|d_{j}(\cdot, t, 0)\right\|, \quad\left\|\frac{\partial}{\partial x} d_{j}(\cdot, t, 0)\right\|, \quad\left\|\frac{\partial^{2}}{\partial x^{2}} d_{j}(\cdot, t, 0)\right\|\right\} \leqq \mathrm{const}  \tag{iv}\\
& \sum_{j}\left\|d_{j}(\cdot, t, 0)\right\|\left(1+j^{2}\right) \leqq \text { const for all } t
\end{align*}
$$

Then, for $h$ sufficiently small,

$$
\begin{align*}
\left|u\left(x_{\nu}, t\right)-v_{\nu}(t)\right| & \leqq c_{1} h^{p}+c_{2} D\left(x_{\nu}, t\right) & \text { in } & T_{1}  \tag{2.8}\\
& \leqq c_{3} D\left(x_{\nu}, t\right) & \text { in } & T_{-1},
\end{align*}
$$

where

$$
D\left(x_{\nu}, t\right)=h^{(q-1) / 2} \exp \left(-h^{-q} d\left(x_{\nu}, t\right)\right)
$$

and where $c_{i}$ are constants independent of $x, t$ and $h$. Here $q=\left(\tau_{\mp}-1\right) / \tau_{\mp}$, with minus sign in $T_{-\operatorname{sign}(\rho)}$ and plus sign in $T_{+\operatorname{sign}(\rho)}$.

Remarks. The assumption, that $\rho \in c_{0}^{\infty}$ simplifies the proof a lot, but it will of course work even for more general coefficients, which are sufficiently smooth. Furthermore, by modifying the proof of a stability theorem of Kreiss [4], it certainly will be possible to prove an analogous theorem for the scheme (2.3), with the following assumptions: (1) the coefficients in the scheme are Lipschitz continuous, and (2) $2 s=p+2-l, l=0,1$.

Note that the $q$ 's of Theorem 2.4-2.5 satisfy $0<q<1$. Furthermore the best $q$ 's for dissipative schemes with even order of accuracy are obtained for schemes with $r=2$. Then

$$
\begin{array}{rlll}
q & =p /(p+2) & \text { in } & T_{+\operatorname{sign}(\rho)}, \\
& =(3 p+2) /(3 p+6) & \text { in } & T_{-\operatorname{sign}(\rho)} .
\end{array}
$$

Example. The Lax-Wendroff scheme $Q_{1} \equiv I, Q_{2}=I+\rho k D_{0}+\sigma \rho^{2}\left(k^{2} / 2\right) D_{+} D_{-}$, where $2 h D_{0}=E-E^{-1}, h D_{+}=E-I$ and $h D_{-}=I-E^{-1}$, is accurate of order 2 and dissipative of order 4 if $\sigma=1$.

Hence, $q=1 / 2$ in $T_{+\operatorname{sign}(\rho)}$ and $q=2 / 3$ in $T_{-\operatorname{sign}(\rho)}$. If $\sigma \neq 1, p=1$ and $2 s=2$, i.e., $q=1 / 2$ on both sides of the characteristic. But the estimates (2.8) show that the local error outside an interval of length $O\left(h^{q}|\log h|\right)$ is only $O\left(h^{2}\right)$ for $\sigma=1$. Therefore the solution to the Lax-Wendroff scheme with $\sigma=1$, behaves fairly well even for discontinuous initial functions.
3. Proofs of Algebraic Theorems. We are going to prove Theorems 2.1-2.3 and begin by rewriting (2.5) in the form

[^1]$$
\hat{Q}(x, t, \xi, 0)=e^{i \rho(x, t) \lambda \xi+U(x, t, \xi)}
$$
where
$$
U(x, t, \xi)=\sum_{\nu=0}^{\infty} c_{\nu}(x, t) \xi^{p+\nu+1}
$$
which is analytic in $\xi$. In this section $\hat{\boldsymbol{Q}}$ stands for either $\hat{Q}_{-1}$ or $\hat{Q}_{0}$.
Proof of Theorem 2.1. Since (2.3) or (2.3') is dissipative of some order $2 s$, we have
\[

$$
\begin{equation*}
e^{\operatorname{Re}(U(x, t, \xi)\}}=|\hat{Q}(x, t, \xi, 0)| \leqq 1-\delta|\xi|^{2 s}, \quad \delta>0 \tag{3.1}
\end{equation*}
$$

\]

and thus $2 s \geqq p+1$. Write $G(i \xi)=\hat{Q}(x, t, \xi, 0)$, then

$$
1=|G(i \xi) / G(-i \xi)|=e^{2 \operatorname{Im}\left(c_{0}\right) \xi^{p+1}+\cdots}
$$

i.e., $\operatorname{Im}\left(c_{0}(x, t)\right)=0$ for odd $p$.

Moreover $c_{0}(x, t) \leqq 0$, since we otherwise could find a $\xi$, such that $|\hat{Q}|>1$. But $p$ is the order of accuracy and hence $c_{0}(x, t) \neq 0$ and $c_{0}(x, t)<0$. Thus we finally conclude that $2 s=p+1$.

Proof of Theorem 2.2. We first start with constant coefficients $\rho$ and $\phi$, where $\rho \neq 0$.

We are going to estimate $U(\xi+i \alpha)=\sum_{\nu=0}^{\infty} c_{\nu}(\xi+i \alpha)^{p+\nu+1}$ and therefore we state the following variant of Hölder's inequality.

Hölder's inequality. Let $A \geqq 0, B \geqq 0$ and $\epsilon>0$ be given. Then

$$
A B \leqq \eta A^{r}+\epsilon B^{q},
$$

where $r=q /(q-1), q>1$ and $\eta=\eta(\epsilon)$ is uniformly bounded for $\epsilon \geqq \epsilon_{0}>0$.
From (3.1) and assumption (ii) it follows that

$$
|\hat{Q}(\xi+i \alpha, 0)|=\exp \left(-\alpha \rho \lambda-\delta^{\prime} \xi^{2 s}+H(\alpha, \xi, p)\right)
$$

where $\delta^{\prime}>0$.
Here

$$
\begin{align*}
H(\alpha, \xi, p)= & \operatorname{Re}\left\{i \alpha U^{\prime}(\xi)+\frac{(i \alpha)^{2}}{2!} U^{\prime \prime}(\xi)+\cdots\right\} \\
= & \operatorname{Re}\left\{c _ { 0 } \left(d_{1}(i \alpha) \xi^{p}+d_{2}(i \alpha)^{2} \xi^{p-1}+\cdots+d_{p}(i \alpha)^{p} \xi+(i \alpha)^{p+1}\right.\right.  \tag{3.2}\\
& \left.+O\left(\sum_{j+k=p+2}|\alpha|^{j} \xi^{k}\right)\right\}
\end{align*}
$$

where

$$
d_{k}=\binom{p+1}{k}, \quad k=1,2, \cdots, p
$$

Now we have from Hölder's inequality the estimates

$$
|\alpha|^{j+1}|\xi|^{p-j} \leqq \eta_{j}|\alpha|^{\tau_{r, j}}+\epsilon \xi^{2 s}
$$

where for $r \geqq 1, \tau_{r, j}=2 s(j+1) /(r+j),(r=2 s-p), j=0,1, \cdots,(p-1)$.
Thus, we can find a $\tau_{r} \geqq 1$ such that

$$
H(\alpha, \xi, p) \leqq \nu_{1}|\alpha|^{\tau_{r}}(1+o(1))+\nu_{2} \xi^{2 s}(1+O(\alpha)+O(\xi))
$$

when $\alpha, \xi \rightarrow 0$ and where $\nu_{2}=O(\epsilon), \epsilon \rightarrow 0$.
It is easy to see that we have to take $\tau_{r}=\min _{j} \tau_{r, j}$ over the terms in (3.2) which are positive. Thus for odd $p$ we conclude that

$$
\tau_{r}^{(1)}=\tau_{r, 1}=\tau_{1,1}
$$

since $r=1$. (Theorem 2.1.) From this it follows that $\tau=(p+1, p+1)$ for odd $p$.
When $p=2 l, l=1,2, \cdots$, we find that

$$
\begin{equation*}
\tau_{r}{ }^{(0)}=\tau_{r, 0} \tag{3.3}
\end{equation*}
$$

when $\operatorname{sign}\left(i c_{0}\right) \operatorname{sign}(\alpha)=1$, and

$$
\boldsymbol{\tau}_{\boldsymbol{r}}^{(2)}=\boldsymbol{\tau}_{r, 2}
$$

when $\operatorname{sign}\left(i c_{0}\right) \operatorname{sign}(\alpha)=-1$.
Note that $\operatorname{Im}\left(i c_{0}\right)=0$. Summarizing, we have

$$
\begin{align*}
|\hat{Q}(\xi+i \alpha, 0)|=\exp \left(-\alpha \rho \lambda-\left(\delta^{\prime}-\nu_{2}\right) \xi^{2 s}+\nu_{1}|\alpha|^{\tau_{r}^{(m)}}+\cdots\right) &  \tag{3.4}\\
& m=0,1,2 .
\end{align*}
$$

Moreover, from assumption (iv) we know that

$$
\begin{equation*}
\operatorname{sign}\left(i c_{0}\right)=\operatorname{sign} \rho, \tag{3.5}
\end{equation*}
$$

and together with (3.3) this gives that (3.3) holds when

$$
\operatorname{sign}(\alpha)=\operatorname{sign}(\rho), \text { i.e. } \quad \alpha \rho>0
$$

Analogously, (3.3') holds for $\alpha \rho<0$, i.e., for even $p, \tau$ is given by

$$
\tau=\left(\tau_{r}^{(2)}, \tau_{r}{ }^{(0)}\right)=(6 s /(r+2), 2 s / r)
$$

Now it is easy to see that we can choose $\epsilon$ so small that (3.4) guarantees contractivity of order $\tau$. Indeed, there exists a $\xi_{0}$, such that for all $\xi,|\xi|<\xi_{0}$,

$$
\begin{equation*}
|\hat{Q}(\xi+i \alpha, 0)|=\exp \left(-\alpha \rho \lambda+\sigma|\alpha|^{\tau_{\tau}^{(m)}}\right) R(\xi), \quad m=0,1,2 \tag{3.6}
\end{equation*}
$$

where $|R(\xi)| \leqq 1$ for $|\xi|<\xi_{0}$. But for $|\xi| \geqq \xi_{0}$ there exists an $\alpha_{0}>0$ such that for all $|\alpha|<\alpha_{0}$

$$
|\hat{Q}(\xi+i \alpha, 0)| \leqq \exp \left(-\alpha \rho \lambda+\sigma|\alpha|^{\tau_{r}^{(m)}}-\delta^{\prime} / 2 \cdot \xi^{2 s}\right), \quad m=0,1,2 .
$$

Thus, Theorem 2.2 is proved for constant coefficients.
From the assumptions of Theorem 2.2 we have local contractivity of order $\tau$ for every $x$ and $t$, i.e.,

$$
\hat{Q}(x, t, \xi+i \alpha, 0)=\exp \left(-\alpha \lambda \rho(x, t)+\sigma(x, t)|\alpha|^{\tau}\right) \times R(x, t, \xi), \quad|\xi| \leqq \pi
$$

What remains to show is that $\sigma(x, t)$ is uniformly bounded. From Definition 2.2, we see that

$$
\hat{Q}(x, t, \xi, 0)=\exp \left(i \rho(x, t) \lambda \xi+c_{0}(x, t) \xi^{p+1}+\cdots\right) .
$$

This relation and assumptions (ii) and (iii) of Theorem 2.2 then show that $c_{0}(x, t)$ is uniformly bounded in the upper halfplane. Now $\nu_{2}=\nu_{2}(\epsilon)$ in (3.4) depends
essentially on $c_{0}(x, t)$ and thus $\nu_{2}=\epsilon \nu_{3}(x, t)$, where $\nu_{3}(x, t)$ is uniformly bounded for $t \geqq 0$.

Therefore it is sufficient to choose $\epsilon(x, t) \geqq \epsilon_{0}>0$ to obtain the relation (3.6) and from Hölder's inequality, we then see that $\nu_{1}(x, t, \epsilon)$ is uniformly bounded. Hence we can conclude that $\sigma(x, t)$ is uniformly bounded for $t \geqq 0$, and the proof of Theorem 2.2 is completed.

Finally we give the proof of Theorem 2.3. (Cf. Strang [10] and [1].)
Proof of Theorem 2.3. Put $e^{i \rho(x, t) \lambda \xi}=\hat{Q}_{2}(x, t, \xi, 0)+R(\xi)$, and take $y=\rho \lambda$; then $\hat{Q}_{2}(\xi, y)$ is the Lagrangian interpolation polynomial of degree $2 N$ through the points $\left\{\left(\nu, e^{i \xi \nu}\right)\right\}_{\nu=-N}^{N}$, and hence the remainder $R(\xi, y)$ is given by

$$
R(\xi, y)=\xi^{p+1} i^{p+1} e^{i \xi \eta} \cdot y \prod_{j=1}^{N}\left(y^{2}-j^{2}\right) /(2 N+1)!, \quad|\eta|<N
$$

Thus

$$
c_{0}(x, t)=-i \rho(x, t) \lambda \prod_{j=1}^{N}\left(j^{2}-\rho^{2} \lambda^{2}\right) /(2 N+1)!
$$

and since $|\rho(x, t) \lambda| \leqq 1$ we see that condition (iv) of Theorem 2.2 holds.
4. Proofs of Analytic Theorems. Proof of Theorem 2.4 (constant coefficients). The generalized solution of (2.1) with $\phi=0$ and with initial function $f(x)$ defined in (2.6) is

$$
\begin{array}{rlrl}
u(x, t)=f(x+\rho t) & =0 & x+\rho t<0  \tag{4.1}\\
& =1 & x+\rho t \geqq 0 .
\end{array}
$$

Now make the transformation

$$
w_{\nu}(t)=e^{\alpha\left(x_{\nu}+\rho t\right) / h} v_{\nu}(t) .
$$

If $v_{\nu}(t)$ satisfies (2.3) then $w_{\nu}(t)$ satisfies

$$
\begin{equation*}
Q_{1}{ }^{\prime} w_{\nu}(t+k)=Q_{2}{ }^{\prime} w_{\nu}(t), \tag{4.2}
\end{equation*}
$$

where

$$
Q_{1}{ }^{\prime} w_{\nu}=e^{-\alpha \rho \lambda} e^{\alpha \nu} Q_{1} e^{-\alpha \nu} w_{v}
$$

and

$$
Q_{2}{ }^{\prime} w_{\nu}=e^{\alpha \nu} Q_{2} e^{-\alpha \nu} w_{\nu} .
$$

Since $\left|\hat{Q}_{\mathbf{1}}(\xi, 0)\right| \geqq$ const $>0$ it follows that for $\alpha$ sufficiently small

$$
\left|\hat{Q}_{1}^{\prime}(\xi, 0)\right|=e^{-\alpha \rho \lambda}\left|\hat{Q}_{1}(\xi+i \alpha, 0)\right| \geqq \text { const }>0
$$

Thus

$$
\hat{Q}_{-1}^{\prime}(\xi+i \alpha, 0)=e^{\alpha \rho \lambda} \hat{Q}_{1}(\xi+i \alpha, 0)^{-1} Q_{2}(\xi+i \alpha, 0)
$$

is well defined.
From assumption (ii) it now follows that

$$
\hat{Q}_{-1}^{\prime}(\xi+i \alpha, 0)=\exp \left(\sigma|\alpha|^{\top \mp}\right) \times R(\xi)
$$

where $|R(\xi)| \leqq 1$ for $|\xi| \leqq \pi$. Then it is an easy matter to see that

$$
\begin{equation*}
\left|Q_{-1}^{\prime}(\xi+i \alpha, h)\right| \leqq 1+\nu_{0} h \quad \text { for } 0 \leqq h \leqq h_{1}, \tag{4.3}
\end{equation*}
$$

where $\nu_{0}$ is independent of $h$, if $\alpha$ is chosen such that

$$
\alpha=\alpha_{\mp}=\operatorname{sign}\left(\rho t+x_{\nu}\right) h^{1 / \tau_{\mp}} .
$$

It is well known (cf. Richtmyer [12]) that (4.3) implies stability of the scheme (4.2), i.e.,

$$
\left\|w_{\nu}(t)\right\|_{h} \leqq e^{\beta t}\left\|w_{\nu}(0)\right\|_{h}=e^{\beta t}\left\|e^{\alpha\left(x_{\nu}+\rho t\right) / h} f\left(x_{\nu}\right)\right\|_{h}
$$

for some constant $\beta$.
If $\alpha<0$, then

$$
\begin{align*}
\left\|w_{\nu}(0)\right\|_{h} & \leqq \max _{\nu} f\left(x_{\nu}\right)\left(\sum_{\nu=0}^{\infty} h e^{2 \alpha \nu}\right)^{1 / 2}=h^{1 / 2} /\left(1-e^{2 \alpha}\right)^{1 / 2}  \tag{4.4}\\
& =O\left(h^{1 / 2}|\alpha|^{-1 / 2}\right) .
\end{align*}
$$

Since $\left\|w_{\nu}(t)\right\|_{h} \geqq h^{1 / 2}\left|w_{\nu}(t)\right|$, we get

$$
\left|v_{\nu}(t)\right| \leqq \nu_{1} e^{\beta t} h^{(q-1) / 2} \exp \left(-h^{-q}\left|x_{\nu}+\rho t\right|\right)
$$

where $q=q_{\mp}=\left(\tau_{\mp}-1\right) / \tau_{\mp}$, and $\nu_{1}$ is a constant.
Note that this estimate only holds to the left of the characteristic $\rho t+x=0$.
Let

$$
\begin{aligned}
f_{1}(x)=f(-x) & =1, \quad x \leqq 0, \\
& =0, \quad x>0 .
\end{aligned}
$$

With $f_{1}(x)$ as initial function we get a corresponding estimate to the right of the characteristic $\rho t+x=0$, when $\alpha>0$.

Now $f_{2}(x)=f_{1}(x)+f_{0}(x) \equiv 1,-\infty<x<\infty,\left(f_{0}=f\right)$ as initial function gives rise to solutions of (2.1) and (2.3) that are identically equal to one. This follows from assumption (i). An elementary application of the triangle inequality then gives

$$
\left|v_{\nu}(t)-1\right| \leqq \nu_{2} e^{\beta t} h^{(q-1) / 2} \exp \left(-h^{-q}\left|x_{\nu}+\rho t\right|\right)
$$

where $\nu_{2}$ is a constant.
This estimate is used to the right of the characteristic $\rho t+x=0$. Thus the proof of Theorem 2.4 is completed.

Before we are able to prove Theorem 2.5 we are going to state the following stability theorem due to Lax-Nirenberg [7].

Theorem 4.1. Consider the difference scheme

$$
y_{\nu}(t+k)=S y_{\nu}(t),
$$

where $S$ is a difference operator of the form

$$
S=\sum d_{j}\left(x_{\nu}, t, 0\right) E^{j} .
$$

Suppose that for $S$ we have

$$
|\hat{S}(x, t, \xi, 0)| \leqq 1
$$

for all $t, x$ and $\xi$.
Furthermore assume that condition (iv) of Theorem 2.5 holds. Then the difference scheme is stable in the sense that

$$
\left\|y_{\nu}(t)\right\|_{h} \leqq e^{\beta t}\left\|y_{\nu}(0)\right\|_{h}
$$

where $\beta$ is a constant.
Proof of Theorem 2.5 (variable coefficients). In this case we have the scheme ( $2.3^{\prime}$ ) and we make the transformation

$$
w_{\nu}(t)=e^{\alpha \rho\left(x_{\nu}, t\right) / h} v_{\nu}(t)
$$

where $g(x, t)$ was defined in Section 2.
The new mesh-function $w_{\nu}(t)$ satisfies

$$
\begin{equation*}
w_{\nu}(t+k)=Q_{0}{ }^{\prime} w_{\nu}(t)=e^{\alpha \rho\left(x_{\nu}, t+k\right) / h} Q_{0} e^{-\alpha_{\rho}\left(x_{\nu}, t\right) / h} w_{\nu}(t) . \tag{4.5}
\end{equation*}
$$

From Definition 2.5 it follows that $\partial g / \partial t=\rho \partial g / \partial x$, and moreover, since $\rho(x, t) \in c_{0}^{\infty}$, that $\partial^{2} g / \partial t^{2}$ and $\partial^{2} g / \partial x^{2}$ is uniformly bounded for $t \geqq 0$.

Thus it is easy to conclude that

$$
\begin{aligned}
\hat{Q}_{0}^{\prime}(x, t, \xi, \alpha, h)= & e^{\alpha \lambda \rho(x, t) \theta_{x}(x, t)} \hat{Q}_{0}\left(x, t, \xi+i \alpha g_{x}(x, t), h\right) \\
& +|\alpha| h \hat{Q}(x, t, \xi), \quad x=x_{\nu},
\end{aligned}
$$

where $\hat{Q}$ is the symbol of a uniformly bounded operator $Q$.
Using assumption (ii) we have for $\alpha$ sufficiently small

$$
\hat{Q}_{0}{ }^{\prime}(x, t, \xi, \alpha, 0)=\exp \left(\sigma(x, t)\left|\alpha g_{x}(x, t)\right|^{\mp \mp}\right) \times R(x, t, \xi, \alpha)
$$

where $|R| \leqq 1$ for $|\xi| \leqq \pi$.
Now

$$
R(x, t, \xi, \alpha)=R(x, t, \xi, 0)+|\alpha|^{\mp \mp} \hat{Q}_{3}(x, t, \xi),
$$

and hence, since $R(x, t, \xi, 0)=\hat{Q}_{0}(x, t, \xi, 0)$,

$$
\hat{Q}_{0}{ }^{\prime}(x, t, \xi, \alpha, h)=\hat{Q}_{0}(x, t, \xi, 0)+\left(|\alpha|^{\mp \mp}+|\alpha| h+h\right) \hat{Q}_{4}(x, t, \xi, h) .
$$

Here, $\hat{Q}_{j}, j=3,4$ correspond to uniformly bounded operators $Q_{j}, j=3,4$.
Now the difference scheme $y_{\nu}(t+k)=Q_{0} y_{\nu}(t)$ satisfies the conditions of Theorem 4.1, and therefore it is stable. If we choose

$$
\alpha=\operatorname{sign}(g(x, t)) h^{1 / \tau} \mp
$$

then the scheme (4.5) is also stable, since $Q_{0}{ }^{\prime}=Q_{0}+h Q_{5}$, where $Q_{5}$ is a uniformly bounded operator. Summarizing, we can find a constant $\beta_{1}$ such that

$$
\left\|w_{\nu}(t)\right\|_{h} \leqq e^{\beta_{1} t}\left\|w_{\nu}(0)\right\|_{h}=e^{\beta_{1} t}\left\|e^{\alpha_{0}(\cdot, 0) / h} f(\cdot)\right\|_{h} .
$$

But now $g(x, 0)=x$, and therefore estimates corresponding to (4.4) hold even in the variable coefficient case. Since $u(x, t)=0$ for $g(x, t)<0$, the estimate (2.8) in $T_{-1}$ follows. Now call the solutions of (2.1) and (2.3') with initial function $f_{2}(x)$ $\equiv 1 u^{(2)}(x, t)$ and $v_{\nu}{ }^{(2)}(t)$ respectively. Then from assumption (i) we have
Table. Eleven numerical experiments for testing $\tau$.

| Experiment | $\rho$ | $Q_{2}\left(Q_{1} \equiv 1\right)$ | / 2 2s |  | $\sigma$ | $\epsilon$ | $n=1 / k$ | $D_{n} n^{-1}$ | $D_{n}+1$ | $D_{n}$ | $E_{h^{-1}}$ | $E_{n}{ }^{+1}$ | $E_{n}$ | Estimute Theoretical for $\tau \quad \tau$ |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1 | 0 | $I+\sigma h^{2} D_{+} D_{-}$ | $\begin{aligned} & 1 \\ & 1 \end{aligned}$ | 2 | $\begin{aligned} & 0.5 \\ & 0.5 \end{aligned}$ | $\begin{aligned} & 0.01 \\ & 0.001 \end{aligned}$ | $\begin{gathered} +00 \\ 500 \\ 500 \\ 1000 \\ 1000 \\ 400 \\ 5000 \\ 800 \\ 1000 \end{gathered}$ |  |  | 966 <br> 117 <br> 131 <br> 147 <br> 120 <br> 136 <br> 176 <br> 194 <br> 196 |  |  | $\left\|\begin{array}{l} 1.366 \\ 1.37 \\ 1.47 \\ 1.44 \end{array}\right\|$ | $\begin{aligned} & (2.2,2.2) \\ & (2.2,2.2) \\ & (1.5,1.8) \\ & (1.9,1.9) \end{aligned}$ | $\begin{aligned} & (2,2) \\ & (2,2) \\ & (2,2) \\ & (2,2) \\ & (2,2) \end{aligned}$ |
| 2 | 0 | $I+0.1 h^{3} D_{0} D_{+} D_{-}+\sigma h^{4} D_{+}^{2} D_{-}^{2}$ | 2 | 4 | -0.025 | 0.001 | $\begin{aligned} & 300 \\ & 330 \\ & 300 \\ & 6060 \\ & \hline 660 \end{aligned}$ | $\begin{aligned} & 14 \\ & 14 \\ & 18 \\ & 19 \\ & 19 \end{aligned}$ | $\begin{aligned} & 60 \\ & 62 \\ & 8.5 \\ & 82 \\ & 92 \end{aligned}$ |  | $\left.\begin{aligned} & 1.29 \\ & 1.36 \end{aligned} \right\rvert\,$ | $\begin{array}{\|l\|} 1.42 \\ 1.45 \end{array}$ |  | $\begin{array}{\|l\|l} (2.7,2.0) \\ (2.3,1.8) \\ \hline \end{array}$ | $\begin{aligned} & (3,2) \\ & (3,2) \end{aligned}$ |
| 3 | 0 | $I+\sigma h^{4} D_{+}^{2} D_{-}^{2}$ | 3 | 4 | -0.1 | 0.01 | $\begin{aligned} & 350 \\ & 100 \\ & 700 \\ & 800 \end{aligned}$ |  |  | 27 <br> 27 <br> 3.31 <br> 33 |  |  | $\begin{array}{\|c} 1.15 \\ 1.22 \\ 1 \end{array}$ | $\begin{aligned} & (5.0,5.0) \\ & (3,5,3,5) \end{aligned}$ | $\begin{gathered} (4,4) \\ (, 4,4) \end{gathered}$ |
| 4 | 1 | $\begin{aligned} & I+k D_{0}+\sigma \frac{k^{2}}{2} D_{+} D_{-} \\ & \text {with } \lambda=1 / 2 \end{aligned}$ |  | 2 | 1.2 <br> 1.6 | $\begin{aligned} & 0.01 \\ & 0.01 \end{aligned}$ | $\begin{array}{r} 3.50 \\ 400 \\ 7700 \\ 800 \\ 450 \\ 4.00 \\ .000 \\ 9000 \\ 1000 \end{array}$ |  |  | $\begin{aligned} & 27 \\ & 27 \\ & 37 \\ & 34 \\ & 36 \\ & 49 \\ & 30 \\ & 68 \\ & 71 \\ & \hline \end{aligned}$ |  |  | $\begin{array}{\|l\|} \hline 1.36 \\ 1.33 \\ 1.39 \\ 1.39 \\ 1.37 \\ \hline \end{array}$ | $\begin{aligned} & (2.3,2,3) \\ & (2.4,2.4) \\ & (2.1,2,1) \\ & (2.2,2.2) \\ & (2.2,2) \end{aligned}$ | $\begin{aligned} & (2,2) \\ & (2,2) \\ & (2,2) \\ & (2,2) \\ & (2,2) \\ & \hline \end{aligned}$ |
| ј | 1 | $\begin{aligned} & I+k D_{0}+\sigma \frac{k^{2}}{2} D_{+} D_{-} \\ & \text {with } \lambda=1 / 2 \end{aligned}$ | 2 | 4 | 1 | 0.01 | $\begin{array}{r} 75 \\ 100 \\ 150 \\ 150 \\ \hline 200 \\ \hline \end{array}$ | $\begin{aligned} & 4 \\ & 5 \\ & 5 \\ & 6 \\ & \hline \end{aligned}$ | $\begin{aligned} & 17 \\ & 20 \\ & 23 \\ & 23 \\ & \hline 28 \\ & \hline \end{aligned}$ |  | $\begin{gathered} 1.25 \\ 1.20 \\ 1.20 \end{gathered}$ | $\begin{array}{\|l\|l} 1.36 \\ 1.40 \\ \hline \end{array}$ |  | $\begin{array}{\|l\|l} (3.1,2 \cdot 2) \\ (3.8,2,2) \\ \hline \end{array}$ | $\begin{gathered} (3,2,2) \\ (3,2) \\ \hline \end{gathered}$ |
| 6 | 1 | $\begin{aligned} & I+k D_{0}+\frac{k^{2}}{2} D_{+} D_{-}+\left(k^{3}-k h^{2}\right) / 6 D_{0} D_{+} D_{-} \\ & \quad+\sigma\left(k^{4}-k^{2} k^{2}\right) / 24 D_{+}^{2} D_{-}^{2} \text { with } \lambda=0.9 \end{aligned}$ | 3 | 4 | $\begin{aligned} & 0.75 \\ & 1.2 \end{aligned}$ | 0.001 |  |  |  | $\begin{aligned} & 17 \\ & 17 \\ & 19 \\ & 120 \\ & 13 \\ & 14 \\ & 14 \\ & 16 \\ & 16 \end{aligned}$ |  |  | $\left\|\begin{array}{l} 1.12 \\ 1.15 \\ 1.23 \\ 1.14 \end{array}\right\|$ | $\begin{aligned} & (6.1,6.1) \\ & (4.2,4.2) \\ & (3.3,3.3) \\ & (3.3,3.3) \end{aligned}$ | $(4,4)$ $(4,4)$ $\left(\begin{array}{l}(4,4) \\ (4,4)\end{array}\right.$ |


| 7 | 1 | the same as in exp. 6 | 4 | 6 | 1 | 0.001 | $\begin{aligned} & 280 \\ & 300 \\ & 660 \\ & 600 \end{aligned}$ | 7 7 8 8 | 16 17 21 22 22 |  | $\begin{aligned} & 1.14 \\ & 1.14 \end{aligned}$ | $\begin{aligned} & 1.31 \\ & 1.29 \end{aligned}$ |  | $\begin{aligned} & (5.3,2.6) \\ & (5.3,2.7) \end{aligned}$ | $\begin{aligned} & (4 . \pi, 3) \\ & (4 . \pi, 3) \end{aligned}$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 8 | 1 | $\begin{gathered} \text { Implicit methods } Q=Q_{1}^{-1} Q_{2} \\ \left(I-k / 2 D_{0}\right)^{-1}\left(9 E+32 I-E^{-1}\right) / 40 \quad \text { with } \quad \lambda=1 / 2 \end{gathered}$ | 1 | 2 | - | 0.01 | $\begin{aligned} & 100 \\ & 1.50 \\ & 200 \\ & 300 \\ & \hline \end{aligned}$ |  |  | 20 <br> 20 <br> 24 <br> 28 <br> 34 |  |  | $\begin{aligned} & 1.40 \\ & 1.42 \end{aligned}$ | $\begin{aligned} & (2.1,2.1) \\ & (2.0,2.0) \\ & \hline \end{aligned}$ | $\begin{aligned} & (2,2,2) \\ & (2,2) \end{aligned}$ |
| 9 | 1 | $\left.\left(I-k / 2 D_{0}\right)^{-1}\left(I+\frac{k}{2} D_{0}\right)\right)$ with $\lambda=1 / 2$ | 2 | 0 | - | 0.01 | $\begin{array}{r} 50 \\ 100 \\ 150 \\ 200 \end{array}$ | $\begin{aligned} & 5 \\ & 6 \\ & 7 \\ & 8 \end{aligned}$ | $\begin{array}{r} 5 t \\ 95 \\ 155 \\ 206 \end{array}$ | $\begin{array}{r} 59 \\ 104 \\ 162 \\ 214 \end{array}$ |  |  |  | $\begin{aligned} & (1,1) \\ & (1,1) \end{aligned}$ |  |
| 10 | 1 | $\begin{aligned} & \left(I-k / 2 D_{0}\right)^{-1}\left(I+\frac{k}{2} D_{0}+\sigma h^{4} D_{+}^{2} D_{-}^{2}\right) \\ & \text { with } \lambda=1 / 2 \end{aligned}$ | 2 | 4 | -0.01 | 0.01 | $\begin{array}{r} 170 \\ 180 \\ 340 \\ 360 \\ \hline \end{array}$ | $\begin{array}{r} 8 \\ 8 \\ 10 \\ 10 \\ \hline \end{array}$ | $\begin{aligned} & 48 \\ & 49 \\ & 69 \\ & 70 \\ & \hline \end{aligned}$ |  | $\begin{aligned} & 1.25 \\ & 1.25 \end{aligned}$ | $\begin{aligned} & 1.44 \\ & 1.45 \end{aligned}$ |  | $\begin{aligned} & (3.1,1.9) \\ & (3.1,1.9) \\ & \hline \end{aligned}$ | $\begin{aligned} & (3,2) \\ & (3,2) \\ & \hline \end{aligned}$ |
| 11 | 1 | $\begin{aligned} & \left(I-k / 2 D_{0}\right)^{-1}\left(I+\frac{k}{2} D_{0}+\sigma h^{6} D_{+}^{3} D_{-}^{3}\right) \\ & \text { with } \lambda=1 / 2 \end{aligned}$ | 2 | 6 | 0.01 | 0.01 | $\begin{aligned} & 170 \\ & 180 \\ & 340 \\ & 360 \end{aligned}$ | 8 8 10 10 | $\begin{aligned} & 46 \\ & 49 \\ & 73 \\ & 77 \end{aligned}$ |  | $\begin{aligned} & 1.25 \\ & 1.25 \end{aligned}$ | $\begin{aligned} & 1.59 \\ & 1.57 \end{aligned}$ |  | $\begin{aligned} & (3.1,1.5) \\ & (3.1,1.5) \end{aligned}$ | $\begin{aligned} & (3,1.5) \\ & (3,1.5) \end{aligned}$ |

$$
\left|u^{(2)}\left(x_{\nu}, k\right)-v_{\nu}{ }^{(2)}(k)\right|=\left|u^{(2)}\left(x_{\nu}, k\right)-Q_{0} u^{(2)}\left(x_{\nu}, 0\right)\right|=\left(h^{p+1}\right) .
$$

A standard calculation then shows that

$$
\left|u^{(2)}\left(x_{\nu}, t\right)-v_{\nu}^{(2)}(t)\right|=O\left(h^{p}\right),
$$

and hence the estimate (2.8) in $T_{+1}$ follows as in the proof of Theorem 2.4.
5. Numerical Experiments. In this section our theoretical results of Sections $2-4$ will be tested by a series of numerical experiments performed on the CDC 3600 computer at the University of Uppsala.

In the main theorem of Section 2, we have estimates of the form

$$
\begin{equation*}
\left|u\left(x_{\nu}, t\right)-v_{\nu}(t)\right| \leqq c(t) \exp \left(-h^{-q}\left|\rho t+x_{\nu}\right|\right) . \tag{5.1}
\end{equation*}
$$

We are now going to study the interval $I_{\epsilon}(t)$ outside which this error is smaller than $\epsilon, \epsilon>0$. Therefore, we let our computer program count the number of points in this interval. We call the number to the left of the characteristic line, $\rho t+x=0$, $D_{h}{ }^{-1}$ and the number to the right $D_{h}{ }^{+1}$, and the total number of points within $I_{\epsilon}(t)$, $D_{h}=D_{h}^{-1}+D_{h}^{+1}$. The width of $I_{\epsilon}(t)=m\left(I_{\epsilon}(t)\right)$, is therefore approximately $h \cdot D_{h}$. We can also find $m\left(I_{\epsilon}(t)\right)$ by solving

$$
\begin{equation*}
c(t) \exp \left(-h^{-q}|\rho t+x|\right)=\epsilon . \tag{5.2}
\end{equation*}
$$

This last equation has two roots $x_{1}, x_{2}, x_{1}<x_{2}$, and hence $m\left(I_{\epsilon}(t)\right)=x_{2}-x_{1}$.
If we also introduce the double step-size $2 h$ and count $D_{2 h}^{-1}, D_{2 h}^{+1}$ and $D_{2 h}$ respectively, we can define the following quantities

$$
\begin{aligned}
E_{h}{ }^{i} & =D_{h}{ }^{i} / D_{2 h}^{i} \doteqdot 2^{-q+1}=\sqrt{ } 2, \quad i=-1,+1 \\
E_{h} & =D_{h} / D_{2 h} \doteqdot \sqrt[\tau]{ } 2, \quad \tau=\tau_{\mp} \text { defined in Section } 2
\end{aligned}
$$

Therefore we could use these quantities to control our theoretical value of $\tau_{\mp}$.
In the table on the following pages we describe 11 numerical experiments. In column 8 we have the number of iterations, called $n$, which we can relate to $h$ and $k$ in the following way. Set $n k=1$ and $h=k / \lambda$. We have thus normalised our results to the time $t=1 . \phi \equiv 0$ in all of the experiments. Observe that for odd $\rho$, we have $\tau_{-}=\tau_{+}$, and therefore we are only considering $E_{h}$ for these experiments.

We are now going to comment on the results of the experiments. A general remark must be that $E_{h}{ }^{i}, i=-1,+1$ and $E_{h}$ are rational step functions of $h$, and when $D_{h}{ }^{i}, i=-1,+1$, and $D_{h}$ respectively are rather small natural numbers, the risk of obtaining crude estimates of $\tau_{\mp}$ increases seriously.

In Experiment 5 we have the Lax-Wendroff scheme mentioned in Section 2, and hence $\tau=(3,2)$.

The difference operator in Experiment 7 is the most accurate one, when $N=2$ in Theorem 2.3. Hence $\tau=(4.5,3)$. The more precise results of Hedström [3] gives for this special case $\tau=(5,3)$, which agrees better with the experiment.

The Crank-Nicolson scheme of Experiment 9 was introduced to see that some sort of dissipation is necessary to get a meaningful result. (Cf. Section 1.)

The results of Experiment 10 verify the theory closely, and finally in the last experiment we see the effect of increasing $r$ from 2 to 4 , i.e., to $2 s=6$.

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[^0]:    * Compare Experiment 7 in Section 5.

[^1]:    ${ }^{* *} h(x, t) \in c_{0}{ }^{\infty}$ if $h$ has partial derivatives of all order and $h \equiv 0$ outside a bounded region in the ( $x, t$ )-plane.

