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## The energy method for a class of hyperbolic equations

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Equazioni a derivate parziali. - The energy method for a class of hyperbolic equations. Nota ${ }^{(*)}$ di Enrico Jannelli, presentata dal Corrisp. E. De Giorgi.

Riassunto. - In questa nota viene introdotto un nuovo metodo per ottenere espressioni esplicite dell'energia della soluzione dell'equazione iperbolica

$$
\begin{equation*}
\left(\frac{\partial}{\partial t}\right)^{m} u+\sum_{\substack{|v|+j \leq m \\ j \leq m-1}} a_{\nu, j}(t)\left(\frac{\partial}{\partial x}\right)^{\nu}\left(\frac{\partial}{\partial t}\right)^{j} u=0 \tag{*}
\end{equation*}
$$

Stimando opportunamente queste espressioni si ottengono nuovi risultati di buona positura negli spazi di Gevrey per l'equazione ${ }^{(*)}$ quando questa è debolmente iperbolica

## § 1. Introduction

Let us consider the following Cauchy problem

$$
\left\{\begin{array}{l}
L u=\left(\frac{\partial}{\partial t}\right)^{m} u+\sum_{\substack{|v| j \leq m \\
j \leq m-1}} a_{v, j}(t)\left(\frac{\partial}{\partial x}\right)^{\nu}\left(\frac{\partial}{\partial t}\right)^{j} u=0 \quad \text { on } \quad \mathbf{R}_{\mathrm{x}}^{\mathrm{n}} \times[0, T] \\
u(0, x)=\varphi_{1}(x)  \tag{1}\\
\quad: \\
\left(\frac{\partial}{\partial t}\right)^{m-1} u(0, x)=\varphi_{m}(x) .
\end{array}\right.
$$

We assume that equation (1) is hyperbolic, in the sense that the principal symbol

$$
\begin{equation*}
\mathrm{P}_{m}(t ; \tau, \xi)=\tau^{m}+\sum_{\substack{|v|+j=m \\ j \leq m-1}} a(t) \xi^{v} \tau^{j} \quad(t ; \xi) \in[0, \mathrm{~T}] \times \mathbf{R}_{\xi}^{n} \tag{2}
\end{equation*}
$$

(*) Pervenuta all'Accademia il 5 ottobre 1985.
has only real roots $\tau_{h}(t ; \xi)$. More precisely, if we define

$$
\begin{equation*}
\lambda=\inf _{\substack{|\xi|=1 \\ t \in[0 ; T] \\ i \neq j}}\left|\tau_{i}(t ; \xi)-\tau_{j}(t ; \xi)\right| \tag{3}
\end{equation*}
$$

then we shall say that equation (1) is strictly hyperbolic if $\lambda>0$ and weakly hyperbolic if $\lambda=0$.

Throughout this work we shall suppose that $a_{\nu, j}(t)$ belong to $\mathrm{L}^{\infty}([0, \mathrm{~T}])$ if $|\nu|+j \leq m-1$; as regards the coefficients $a_{\nu, j}(t)$ of the principal part $(|\nu|+j=m)$, we suppose for the moment that they belong to $\mathrm{C}^{1}([0, \mathrm{~T}])$.

We want to obtain suitable «a priori» estimates on $u(t, x)$, depending on the coefficients $a_{v, j}(t)$ and on the initial data $\varphi_{1}(x) \ldots \varphi_{m}(x)$.

To this aim, we introduce the Fourier transform of $u(t, x)$ (with respect to $x)$ and of $\varphi_{h}(x)(h=1 \ldots m)$

$$
\begin{equation*}
v(t ; \xi)=\int_{\mathbf{R}_{\mathbf{x}}^{n}} u(t ; x) e^{-i(\xi, x\rangle} \mathrm{d} x \tag{4}
\end{equation*}
$$

$$
\begin{equation*}
\widehat{\varphi}_{h}(\xi)=\int_{\mathbf{R}_{\mathrm{x}}^{n}} \varphi_{h}(x) e^{-i(\xi, x\rangle} \mathrm{d} x \quad h=1 \ldots m \tag{5}
\end{equation*}
$$

Then we define
(6)

$$
\begin{aligned}
& \mathrm{V}(t ; \xi)=\left|\begin{array}{c}
(i|\xi|)^{m-1} v \\
(i|\xi|)^{m-2} \frac{\partial}{\partial t} v \\
\vdots \\
i|\xi|\left(\frac{\partial}{\partial t}\right)^{m-2} v \\
\left(\frac{\partial}{\partial t}\right)^{m-1} v
\end{array}\right| \quad m-\text { vector ; } \\
& \mathrm{H}_{m-j}(t ; \xi)=-\sum_{|\nu|=m-j} a_{\nu, j}(t) \frac{(i \xi)^{\nu}}{|\xi|^{m-j}} ; \\
& \left.\mathrm{A}(t ; \xi)=\left\lvert\, \begin{array}{ccccc}
0 & 1 & 0 & \ldots & \ldots \\
0 & 0 & 1 & 0 & \ldots
\end{array}\right.\right)
\end{aligned}
$$

$$
\Phi(\xi)=\left|\begin{array}{c}
\widehat{\varphi_{1}}(\xi)  \tag{9}\\
\cdot \\
\cdot \\
\cdot \\
\cdot \\
\widehat{\varphi}_{m}(\xi)
\end{array}\right| \quad m \text {-vector }
$$

Now problem (1) may be transformed into a family of first order ordinary differential systems, parametrized by $\xi$, as follows (here and in the following, if F is any function, vector or matrix, we shall denote $\frac{\partial}{\partial t} \mathrm{~F}$ by $\mathrm{F}^{\prime}$ ):

$$
\begin{cases}\mathrm{V}^{\prime}(t ; \xi)=i|\xi| \mathrm{A}(t ; \xi) \mathrm{V}(t ; \xi)+\mathrm{B}(t ; \xi) \mathrm{V}(t ; \xi)  \tag{10}\\ \mathrm{V}(0 ; \xi)=\Phi(\xi) & |\xi| \geq 1 ; t \in[0, \mathrm{~T}]\end{cases}
$$

where $\mathrm{B}\left(t ; \xi_{)}\right)$is a bounded $m \times m$ matrix (we remark that we confine ourselves to consider only $|\xi| \geq 1$ ).

In view of the Paley-Wiener theorem, we are interested in estimating the growth of $|\mathrm{V}(t ; \xi)|$ as $|\xi| \rightarrow+\infty$; now, as the eigenvalues of $\mathrm{A}(t ; \xi)$ are the roots $\tau_{h}\left(t ; \frac{\xi}{|\xi|}\right)$ of the principal symbol $\mathrm{P}_{m}\left(t ; \frac{\xi}{|\xi|}\right)$, we see that, in the strictly hyperbolic case, the matrix $\mathrm{A}(t ; \xi)$ may be uniformly diagonalized; this is the classical method used to obtain estimates on $|\mathrm{V}(t ; \xi)|$ and, therefore, on $u(t, x)$ (see, for instance, [6]).

On the other hand, in the weakly hyperbolic case $\mathrm{A}(t ; \xi)$ is no longer diagonalizable; moreover, the roots $\tau_{h}\left(t ; \frac{\xi}{|\xi|}\right)$ may not be regular in $t$.

Nevertheless, we shall perform an energy method which will be useful both in strictly and weakly hyperbolic cases, and will lead us to obtain explicit expressions for the energy of $\mathrm{V}(t ; \xi)$, which will be written only in terms of $v$, its derivatives, the dual variable $\xi$ and the coefficients $a_{v, j}(t)$ of the principal part of equation (1) (|v|+j=m). This energy method, as we shall see later on, will allow us to get some new results of Gevrey well-posedness for equation (1) in the weakly hyperbolic case; these results will appear in [5].

## § 2. The energy method

Let us consider problem (10). Let $\mathrm{Q}(t ; \xi)$ be any symmetric non negative real-valued $m \times m$ matrix of $\mathrm{C}^{1}$-class in $t$.

We set (here and in the following $\langle$,$\rangle denotes the usual bracket in \mathrm{C}^{m}$ )

$$
\begin{equation*}
\mathrm{E}(t ; \xi)=\langle\mathrm{Q}(t ; \xi) \mathrm{V}(t ; \xi), \mathrm{V}(t ; \xi)\rangle \tag{11}
\end{equation*}
$$

and we call $\mathrm{E}(t ; \xi)$ an energy for $\mathrm{V}(t ; \xi)$. Obviously, we must carefully
choose the matrix Q . In fact, deriving (11) with respect to $t$ and taking into account (10), we easily get

$$
\begin{equation*}
\mathrm{E}^{\prime}=\left\langle\mathrm{Q}^{\prime} \mathrm{V}, \mathrm{~V}\right\rangle+i|\xi|\left\langle\left(\mathrm{QA}-\mathrm{A}^{*} \mathrm{Q}\right) \mathrm{V}, \mathrm{~V}\right\rangle+\langle(\mathrm{QB}+\mathrm{B} * \mathrm{Q}) \mathrm{V}, \mathrm{~V}\rangle \tag{12}
\end{equation*}
$$

where $A^{*}$ is the transposed of $A$ (which is real-valued) and $B^{*}$ is the transposed conjugated of $B$.

Now, in view of the Paley-Wiener theorem, we require that $\mathrm{E}(t ; \xi)$ has, if possible, the same behaviour as $\mathrm{E}(0 ; \xi)$ with respect to $\xi$ when $|\xi| \rightarrow+\infty$; if we want to obtain this fact, it is clear from (12) that we must try to find a nonzero matrix $\mathrm{Q}(t ; \xi)$ such that $\mathrm{QA}=\mathrm{A} * \mathrm{Q}$. This is just what we have obtained; more precisely, the following theorem holds:

Theorem 1. Given the hyperbolic operator L with principal symbol $\mathrm{P}_{m}$, there exists a symmetric real-valued $m \times m$ matrix $\mathrm{Q}(t ; \xi)$ having the following properties:
(13) the entries of $\mathrm{Q}(t ; \xi)$ are polynamials of $m$ variables, whose coefficients depend only on $m$, calculated in the m-tuple $\left(\mathrm{H}_{1} \ldots \mathrm{H}_{m}\right)$ defined by (7);
(14) $\quad \mathrm{Q}(t ; \xi)$ is weakly positive defined if L is weakly hyperbolic, and it is strictly positive defined if L is strictly hyperbolic;

$$
\begin{equation*}
\mathrm{Q}(t ; \xi) \mathrm{A}(t ; \xi)=\mathrm{A}^{*}(t ; \xi) \mathrm{Q}(t ; \xi) \quad \text { for any } \quad(t ; \xi) . \tag{15}
\end{equation*}
$$

Before we go on, let us consider a few examples.

Example 1. (The case $m=2$ ).
If $L$ is a second order operator, then

$$
\mathrm{A}(t ; \xi)=\left(\begin{array}{cc}
0 & 1  \tag{16}\\
\mathrm{H}_{2} & \mathrm{H}_{1}
\end{array}\right) .
$$

In this case, the matrix Q given by Theorem 1 is

$$
\mathrm{Q}(t ; \xi)=\left(\begin{array}{cc}
2 \mathrm{H}_{2}+\mathrm{H}_{1}^{2} & -\mathrm{H}_{1}  \tag{17}\\
-\mathrm{H}_{1} & 2
\end{array}\right)
$$

For instance, if we consider the equation

$$
\begin{equation*}
\mathrm{L} u=\frac{\partial^{2}}{\partial t^{2}} u-a(t) \frac{\partial^{2}}{\partial x^{2}} u-b(t) \frac{\partial^{2}}{\partial t \partial x} u=0 \tag{18}
\end{equation*}
$$

then we have (see (7))

$$
\begin{equation*}
\mathrm{H}_{1}(t ; \xi)=b(t) \frac{\xi}{|\xi|} \quad ; \quad \mathrm{H}_{2}(t ; \xi)=a(t) \tag{19}
\end{equation*}
$$

Hence

$$
\mathrm{Q}(t ; \xi)=\left|\begin{array}{cc}
2 a(t)+b^{2}(t) & b(t) \frac{\xi}{|\xi|}  \tag{20}\\
b(t) \frac{\xi}{|\xi|} & 2
\end{array}\right|
$$

and

$$
\begin{align*}
\mathrm{E}(t ; \xi)=\langle\mathrm{QV}, \mathrm{~V}\rangle & =\left(2 a(t)+b^{2}(t)\right) \xi^{2}|v|^{2}+2\left|v^{\prime}\right|^{2}-  \tag{21}\\
& -2 \xi b(t) \operatorname{Im}\left(v \bar{v}^{\prime}\right)
\end{align*}
$$

We remark that, when $b(t) \equiv 0$, equation (18) reduces to a wave-type equation and the energy $\mathrm{E}(t ; \xi)$ reduces (save a multiplicative constant) to the well-known energy $\mathrm{E}(t ; \xi)=a(t) \xi^{2}|v|^{2}+\left|v^{\prime}\right|^{2}$.

Example 2 (the case $m=3$ ).
If $L$ is a third order operator, then
(22) $\mathrm{A}(t ; \xi)=$

$$
\left|\begin{array}{lll}
0 & 1 & 0 \\
0 & 0 & 1 \\
\mathrm{H}_{3} & \mathrm{H}_{2} & \mathrm{H}_{1}
\end{array}\right| .
$$

In this case, the matrix Q given by Theorem 1 is

$$
\mathrm{Q}(t ; \xi)=\left|\begin{array}{llc}
\mathrm{H}_{2}^{2}-2 \mathrm{H}_{1} \mathrm{H}_{3} & \mathrm{H}_{1} \mathrm{H}_{2}+3 \mathrm{H}_{3} & -\mathrm{H}_{2}  \tag{23}\\
\mathrm{H}_{1} \mathrm{H}_{2}+3 \mathrm{H}_{3} & 2 \mathrm{H}_{1}^{2}+2 \mathrm{H}_{2} & -2 \mathrm{H}_{1} \\
-\mathrm{H}_{2} & -2 \mathrm{H}_{1} & 3
\end{array}\right| .
$$

For instance, if we consider the equation

$$
\begin{equation*}
\mathrm{L} u=\frac{\partial^{3}}{\partial t^{3}} u-a(t) \frac{\partial^{3}}{\partial^{2} t \partial x} u-b(t) \frac{\partial^{3}}{\partial t \partial^{2} x} u-c(t) \frac{\partial^{3}}{\partial x^{3}} u=0 \tag{24}
\end{equation*}
$$

then we have

$$
\begin{equation*}
\mathrm{H}_{1}(t ; \xi)=a(t) \frac{\xi}{|\xi|} ; \mathrm{H}_{2}(t ; \xi)=b(t) ; \mathrm{H}_{3}(t ; \xi)=c(t) \frac{\xi}{|\xi|} . \tag{25}
\end{equation*}
$$

Hence

$$
\begin{align*}
\mathrm{E}(t ; \xi)= & \langle\mathrm{QV}, \mathrm{~V}\rangle=  \tag{26}\\
= & \left(b^{2}(t)-2 a(t) c(t)\right) \xi^{4}|v|^{2}+ \\
& +2\left(a^{2}(t)+b(t)\right) \xi^{2}\left|v^{\prime}\right|^{2}+3\left|v^{\prime \prime}\right|^{2}- \\
& -2(a(t) b(t)+3 c(t)) \xi^{3} \operatorname{Im}\left(\bar{v} \bar{v}^{\prime}\right)+ \\
& +2 b(t) \xi^{2} \operatorname{Re}\left(v \bar{v}^{\prime \prime}\right)+4 a(t) \xi \operatorname{Im}\left(v^{\prime} \bar{v}^{\prime \prime}\right) .
\end{align*}
$$

It is clear that in the strictly hyperbolic case, when Q is strictly positive defined, any estimate regarding $\mathrm{E}(t ; \xi)$ obviously corresponds to an estimate on $\mathrm{V}(t ; \xi)$.

The situation is much more complicated in the weakly hyperbolic case, due to possible degenerations of $\mathrm{Q}(t ; \xi)$; however, we shall take advantage, in this case, of the fact that $\mathrm{Q}(t ; \xi)$ has the same regularity (in $t$ ) of the coefficients $a_{\nu, j}(t)$ of the principal part: in fact, if $a_{\nu, j}(t) \in \mathrm{C}^{p, \alpha}([0, \mathrm{~T}])$, so it is for $\mathrm{H}_{k}(t$; $\xi)(k=1 \ldots m)$ and then for $\mathrm{Q}(t ; \xi)$. In this sense, we see that the loss of regularity of the characteristic roots $\tau_{h}\left(t ; \frac{\xi}{|\xi|}\right)$ is not an intrinsic difficulty.

When L is weakly hyperbolic, we shall consider some perturbed energies $\mathrm{E}_{\varepsilon}(t ; \xi)=\left\langle\left(\mathrm{Q}+\Gamma_{\varepsilon}\right) \mathrm{V}, \mathrm{V}\right\rangle$, where $\Gamma_{\varepsilon}$ is a suitable diagonal matrix with constant coefficient such that $Q+\Gamma_{\varepsilon}$ is strictly positive defined, and a keypoint of our construction will be made up by the evaluation (as a function of $\varepsilon$ ) of the ratio $\left\langle\mathrm{Q}^{\prime} \mathrm{V}, \mathrm{V}\right\rangle /\left\langle\left(\mathrm{Q}+\Gamma_{\varepsilon}\right) \mathrm{V}, \mathrm{V}\right\rangle$, to estimate which we shall use the regularity in $t$ of $\mathrm{Q}(t ; \xi)$ together with some consequences of the following

Lemma 1 ([5]). Let $f(t):[0, \mathrm{~T}] \rightarrow \mathbf{R}$ be a $\mathrm{C}^{\mathbf{1}, \alpha}$ function, with $f(t) \geq 0$. Then there exists a constant $\mathrm{C}=\mathrm{C}\left(\|f\|_{1, \alpha}\right)$ such that

$$
\begin{equation*}
\left|f^{\prime}(t)\right| \leq \mathrm{C}\left[\frac{f(t)}{t(\mathrm{~T}-t)}\right]^{\alpha / 1+\alpha} \quad \forall t \in(0, \mathrm{~T}) \tag{27}
\end{equation*}
$$

On the other hand, if the coefficients $a_{v, j}(t)$ of the principal part of (1) are only hölder continuous, we shall consider some suitable perturbations $\mathrm{Q}_{\varepsilon}(t ; \xi)$ of $\mathrm{Q}(t ; \xi)$ in such a way that $\mathrm{Q}+\Gamma_{\varepsilon}$ is regular in $t$ and strictly positive defined; obviously, we shall estimate the error terms so introduced.

By means of the techniques exposed here, we get the following

Theorem 2 ([5]). Let us consider the operator L in (1). We suppose that
L is weakly hyperbolic, i.e. $\lambda=0$, where $\lambda$ is defined by (3).

Let $k$ be the greatest multiplicity of the roots $\tau_{h}(t ; \xi)$ af the principal symbol $\mathrm{P}_{m}(t ; \tau, \xi)$ defined by (2). Obviously, $k$ is an integer such that $2 \leq k \leq m$. Then
i) if the coefficients $a_{y, j}(t)$ of the principal part of $\mathrm{L}(|\nu|+j=m)$ belong to $\mathrm{C}^{0, \alpha}([0, \mathrm{~T}])$, the Cauchy problem (1) is well-posed in the Gevrey spaces $\gamma_{l o c}^{(s)}$ for

$$
\begin{equation*}
1 \leq s<1+\frac{\alpha}{(\alpha+1)(k-1)+1-\alpha} ; \tag{29}
\end{equation*}
$$

ii) if the coefficients $a_{v, j}(t)$ of the principal part of $\mathrm{L}(|\vee|+j=m)$ belong to $\mathrm{C}^{1, \alpha}([0, \mathrm{~T}])$, the Cauchy problem (1) is well-pased in the Gevrey spaces $\gamma_{l o c}^{(s)}$ for

$$
\begin{equation*}
1 \leq s<1+\frac{1+\alpha}{2(k-1)} \tag{30}
\end{equation*}
$$

We recall that a function $f(x): \mathbf{R}_{x}^{n} \rightarrow \mathbf{R}$ belongs to $\gamma_{l o c}^{(s)}$ if for any K compact subset of $\mathbf{R}$ there exist $\Lambda_{k}$ and $\mathrm{A}_{k}$ such that

$$
\begin{equation*}
\left|\mathrm{D}^{\alpha} f(x)\right| \leq \Lambda_{k} \mathrm{~A}_{k}^{|\alpha|}(|\alpha|!)^{s} \quad x \in \mathrm{~K} \tag{31}
\end{equation*}
$$

We point out that the results of Theorem 2, in the cases $a_{v, j}(t) \in \mathrm{C}^{0,1}$ ([0, $\mathrm{T}])$ or $a_{\nu, j}(t) \in \mathrm{C}^{1,1}([0, \mathrm{~T}])(|\vee|+j=m)$, have been already obtained, using quite different techniques, by T. Nishitani, who has considered in [8] the case of coefficients depending also on $x$.

Moreover, the results of Theorem 2 generalize some previous results regarding second order equations, which have been obtained in [2].

We remark that, if we suppose in Theorem 2 that $k=1$ (i.e. strict hyperbolicity), then by (29) we get that problem (1) is well-posed in $\gamma_{l o c}^{(s)}$ for $1 \leq s<$ $<1 /(1-\alpha)$; this result has already been proved in the general case of regularly hyperbolic systems with coefficients depending on $x$ and $t$ (see [3] and [4]), while the same result had been previously obtained in some particular cases at first in [1] and then in [7].

Let us also observe that we deduce by (30) that, if the coefficients $a_{v, j}(t)$ of the principal part belong to $\mathrm{C}^{1,1}([0, \mathrm{~T}])$, then problem (1) is well-posed in $\gamma_{l o c}^{(s)}$ for $1 \leq s<k /(k-1)$, and this result is not improvable, because, if $s>$ $>k /(k-1)$, problem (1) is not well-posed in $\gamma_{l o c}^{(s)}$ even in the case of constant coefficients.

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