# DECAY OF SOLUTIONS OF SYMMETRIC HYPERBOLIC SYSTEMS OF PARTIAL DIFFERENTIAL EQUATIONS 

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#### Abstract

We consider systems of the form $u_{t}+\sum_{j=1}^{n} A_{j} u_{x_{j}}=0$, where the $A_{j}^{\prime}$ 's are constant $k \times k$ (hermitian) symmetric matrices, and $u$ is a column vector of $k$ components. We use Fourier transform to prove that non-static solutions decay in time at every point $x$. As a consequence, it follows that the energy of any such solution decays locally. More generally, we show that if $B(t)$ is a set which does not increase "too" fast, the energy in $B(t)$ of any non-static solution also decays.


1. Introduction. We consider systems of the form

$$
\begin{equation*}
\frac{\partial u}{\partial t}+\sum_{j=1}^{n} A_{j} \frac{\partial u}{\partial x_{j}}=0, \tag{1}
\end{equation*}
$$

where the $A_{j}$ 's are constant $k \times k$ (hermitian) symmetric matrices, and $u$ is a column vector of $k$ components. These are functions of the independent variables $t \in \mathbf{R}$ and $x=\left(x_{1}, \cdots, x_{n}\right) \in \mathbf{R}^{n}$. Systems of this type are the general form of a large number of equations of mathematical physics, such as Maxwell's equation, the equations of transmission lines, acoustics, elasticity (see Appendix in [7]), and even the equations of magnetogasdynamics (see [1]).

It is customary to discuss the above systems under additional assumptions on the matrices $A_{j}$. One such assumption is that the roots $\lambda=$ $\lambda(p)$ of the characteristic equation

$$
\begin{equation*}
P(\lambda, p)=\operatorname{det}\left(\lambda I-\sum_{j=1}^{n} p_{j} A_{j}\right)=0 \tag{2}
\end{equation*}
$$

are all different from zero for $p \neq 0$, that is, the operator $\sum_{j=1}^{n} A_{j} \partial / \partial x_{j}$ is elliptic ([4], p. 178); or a fixed number of them never vanish for $p \neq$ 0 ([3]); or the assumption contained in the definition of uniformity propagative systems of Wilcox ([7]). In our treatment we impose no restrictions on the $A_{j}{ }^{\prime}$ s other than those stated in the previous paragraph. This is important because there are systems, such as those of magnetogasdynamics, which possess roots $\lambda(p)$ that vanish for certain $p \neq 0$, but not identically. It has been shown that if a characteristic root $\lambda(p)$ is not identically zero then the set of those $p$ where $\lambda(p)=0$ is of measure zero ([1]). Since the $\lambda(p)$, for $|p|=1$, are speeds of propagation of
plane waves in the direction $p$, it is thus reasonable to expect that for a given solution $u$, the part of it associated with such a $\lambda(p)$ should decay in time. Our objective in this paper is to prove this result and derive from it the local decay of the energy of $u$, provided (of course) $u$ is non-static. It also follows that the energy of any solution in a "slowly increasing" set is asymptotically constant.
2. The solution of the Cauchy Problem. We use the Fourier transform to derive an explicit formula for the solution of equation (1) with initial value $u(x, 0)=f(x)$. We assume initially that $f \in \mathscr{\rho}^{k}$, where $\mathscr{\rho}$ is the Schwartz space of rapidly decreasing functions. The Fourier transform of $g \in \rho^{k}$ is defined by

$$
\hat{g}(p)=(2 \pi)^{-n / 2} \int_{\mathbf{R}^{n}} e^{-i p \cdot x} g(x) d x
$$

and the inverse formula

$$
g(x)=(2 \pi)^{-n / 2} \int_{\mathbf{R}^{n}} e^{i p \cdot x} \hat{g}(p) d p
$$

holds. Since, under Fourier transformation, differentiation with respect to $x_{j}$ changes into multiplication by $i p_{j}$, equation (1) transforms into

$$
\begin{equation*}
\hat{u}_{t}+i A(p) \hat{u}=0 \tag{3}
\end{equation*}
$$

where $A(p)=\sum_{j=1}^{n} p_{j} A_{j}$.
We want to solve equation (3) with the initial condition $\hat{u}(p, 0)=$ $\hat{f}(p)$. Let $\left\{e_{j}(p)\right\}$ be a complete set of orthonormal eigenvectors of $A(p)$ with corresponding eigenvalues $\lambda_{j}(p)$. Setting

$$
\hat{u}_{j}(p, t)=\hat{u}(p, t) \cdot e_{j}(p)
$$

and

$$
\hat{f}(p)=\hat{f}(p) \cdot e_{j}(p)
$$

equation (3) gives us, upon scalar multiplication by $e_{j}(p)$,

$$
\frac{\partial \hat{u}_{j}}{\partial t}=-i \lambda_{j}(p) \hat{u}_{j}(p)
$$

The solution of this equation that satisfies the initial condition $\hat{u}_{j}(p, 0)$ $=\hat{f}_{j}(p)$ is

$$
\hat{u}_{j}(p, t)=\hat{f}_{j}(p) e^{-i \lambda j(p) t}
$$

Therefore,

$$
\hat{u}(p, t)=\sum_{j=1}^{k} \hat{f}_{j}(p) e^{-i \lambda_{j}(p) t} e_{j}(p)
$$

At this point we note that the eigenvectors $e_{j}(p)$ can be chosen in such way that they are measurable functions of $p$. That this is possible was proved by C. H. Wilcox (see [9], Theorem 2), and we assume that this choice has been made. We note also that the roots $\lambda_{j}(p)$ are continuous functions of $p$ ([9], Theorem 1). Therefore, we are justified in taking the inverse Fourier transform of the above expression for $\hat{u}$. We obtain

$$
\begin{equation*}
u(x, t)=(2 \pi)^{-n / 2} \sum_{j=1}^{k} \int_{\mathbf{R}^{n}} \hat{f}_{j}(p) e^{i\left[p \cdot x-\lambda_{j}(p) t\right]} e_{j}(p) d p \tag{4}
\end{equation*}
$$

This formula shows that $u$ is the superposition of $k$ waves: each of these waves, in turn, is the superposition-given by the integration over $p$-space-of the plane waves

$$
(2 \pi)^{-n / 2} \hat{f}_{j}(p) e^{i\left[p \cdot x-\lambda_{j}(p) t\right]} e_{j}(p) .
$$

This wave is a signal that propagates in the direction $p$ with speed $v_{j}(p)$ $=\lambda_{j}(p) /|p|$. Each one of the $k$ terms of the sum in (4) is referred to as a normal mode of propagation. We see that each mode is associated with an eigenvalue (counting multiplicity) of $A(p)$, that is, with a speed of plane wave propagation. A $j^{\text {th }}$ mode is excited or not, according to whether the initial value $f$ is such that $\hat{f}_{j}(p)$ is different from zero or not, respectively. If $\lambda_{j}(p) \equiv 0$, then the associated mode does not depend on the time: it is static. Thus, a smooth solution $u$ is the superposition of its static and non-static parts,

$$
\begin{equation*}
u(x, t)=u_{s}(x, t)+u_{n s}(x, t) \tag{5}
\end{equation*}
$$

and it is clear what a non-static solution means in this case. Considering the Hilbert space $H=L^{2}\left(\mathbf{R}^{n}\right)^{k}$ with its usual norm

$$
\|f\|^{2}=\int_{\mathbf{R}^{n}}|f(x)|^{2} d x=\sum_{j=1}^{k} \int_{\mathbf{R}^{n}}\left|f_{j}(x)\right|^{2} d x,
$$

we notice that for any $f \in H$, (4) defines a function $u(\cdot, t)=U(t) f \in$ $H$ for each $t \in \mathbf{R}: u(\cdot, t)=U(t) f$ is called a solution with finite energy (clearly, $u(\cdot, t)$ is not, in general, a classical solution of (1)). It can be shown that $U(t)=\exp (-t A), t \in R$, the group of unitary operators generated by the skew selfadjoint operator

$$
A=\sum_{j=1}^{n} A_{j} \frac{\partial}{\partial x_{j}}, \quad D(A)=\{f \in H \mid A(p) \hat{f} \in H\}
$$

A solution with finite energy $u(\cdot, t)=U(t) f$ is said to be non-static if $f$ belongs to the orthogonal complement $N(A)^{\perp}$ of the null space of $A$.
3. Decay of solution. Let us enumerate, in decreasing order of magnitude, the solutions of equation (2) that are not identically zero, say, a total of $r$ :

$$
\begin{equation*}
\lambda_{1}(p) \geqq \cdots \geqq \lambda_{r}(p) \tag{6}
\end{equation*}
$$

The remaining $k-r$ roots are then identically zero:

$$
\begin{equation*}
\lambda_{r+1}(p)=\cdots=\lambda_{k}(p) \equiv 0 \tag{7}
\end{equation*}
$$

It follows that $u(\cdot, t)=U(t) f$ is a non-static solution if and only if

$$
\begin{equation*}
\hat{f_{j}}=0 \text { for } j=r+1, \cdots, k \tag{8}
\end{equation*}
$$

Theorem 1. For each $m=1,2, \cdots$ and $f$ in a dense subset $S$ of $N(A)^{\perp}$, there exists a constant $C_{m, f}>0$ such that

$$
\begin{gathered}
|u(x, t)| \leqq C_{m, f}\left(1+|x|+\cdots+|x|^{m}\right)|t|^{-m} \\
\text { for all } t \neq 0, \text { for all } x \in \mathbf{R}^{n}
\end{gathered}
$$

Proof. Let us introduce polar coordinates $\rho$, $\omega$ with $\rho>0,|\omega|=1$, through the relation $p=\rho \omega$. Making use of the (easy to check) facts that $\lambda_{j}(\rho \omega)=\rho \lambda_{j}(\omega)$ and $e_{j}(\rho \omega)=e_{j}(\omega)$, a non-static solution with initial value $f \in \mathscr{\rho}^{k}$ is given by (see (4) and (8))

$$
u(x, t)=
$$

$$
\begin{equation*}
(2 \pi)^{-n / 2} \sum_{j=1}^{r} \int_{|\omega|=1}\left\{\int_{0}^{\infty} e^{i \rho \omega \cdot x} e^{-i t \rho_{j}(\omega)} \hat{f}_{j}(\rho \omega) e_{j}(\omega) \rho^{n-1} d \rho\right\} d S_{\omega} \tag{9}
\end{equation*}
$$

We further restrict $f$ in the following way: each $\hat{f}_{j}(\rho \omega)$ is assumed to be of the form

$$
\hat{f}_{j}(\rho \omega)=\varphi_{j}(\rho) \psi_{j}(\omega)
$$

where $\varphi_{j} \in C_{0}^{\infty}(0, \infty)$ and $\psi_{j}$ is taken from the set of $C^{\infty}$ functions on the unit sphere $S^{n-1}$ which vanish on some neighborhood of

$$
N_{j}=\left\{\omega \in S^{n-1} \mid \lambda_{j}(\omega)=0\right\}
$$

We denote by $S$ this set of data $f$. It was proved in [1] that if a root $\lambda_{j}(\omega)$ does not vanish identically then the corresponding set $N_{j}$ is of measure zero in $S^{n-1}$. Also, $N_{j}$ is closed, for $\lambda_{j}(p)$ is a continuous function. It thus follows easily that $S$ is dense in $N(A)^{\perp}$.

We are now in a position to estimate the expression in (9). Each term $u_{j}(x, t)$ there can be written

$$
u_{j}(x, t)=(2 \pi)^{-n / 2} \int_{S^{n-1} \mid V_{j}}\left(\int_{\alpha}^{\beta} e^{-i t \lambda_{j} \lambda_{j}(\omega)} F_{j} d \rho\right) d S_{\omega},
$$

where

$$
F_{j}=e^{i \rho \omega \cdot x} \varphi_{j}(\rho) \psi_{j}(\omega) e_{j}(\omega) \rho^{n-1}
$$

$\varphi_{j} \in C_{0}{ }^{\infty}(\alpha, \beta)$ with $0<\alpha<\beta<\infty$, and $V_{j}$ is some open neighborhood of $N_{j}$. If we integrate by parts once, we obtain

$$
u_{j}(x, t)=(2 \pi)^{-n / 2} \int_{S^{n-1} \backslash V_{j}}\left(\int_{\alpha}^{\beta} \frac{e^{-i t \rho_{\rho}(\omega)} F_{j}^{\prime}}{i t \lambda_{j}(\omega)} d \rho\right) d S_{\omega},
$$

where $F_{j}{ }^{\prime}=\partial F_{j} / \partial \rho$. We can repeat this procedure any number of times, say $m$ times, thus obtaining

$$
u_{j}(x, t)=(2 \pi)^{-n / 2} \int_{S^{n-1} \backslash V_{j}}\left(\int_{\alpha}^{\beta} \frac{e^{-i t \rho \lambda_{j}(\omega)} F_{j}^{(m)}}{\left[i t \lambda_{j}(\omega)\right]^{m}} d \rho\right) d S_{\omega}
$$

Now, since each $\lambda_{j}(\omega)$ is continuous, there exists $a_{j}>0$ such that $\left|\lambda_{j}(\omega)\right|$ $\geqq a_{j}>0$ on the set $S^{n-1} \backslash V_{j}$. Thus we can estimate $u_{j}(x, t)$ above by

$$
\frac{C\left(1+|x|+\cdots+|x|^{m}\right)}{|t|^{m}}
$$

where $C$ is a constant that depends on $\varphi_{j}, \psi_{j}, V_{j}, \alpha, \beta$ and $m$. The proof is complete.

Now, if $K \subset \mathbf{R}^{n}$ is any compact set and $u(\cdot, 0)=f \in S$, we obtain the fact that the energy of $u(\cdot, t)$ over $K$ decays faster than any power of $1 / t$ :

Corollary. For $u(\cdot, 0)=f \in S$ and $m=1,2, \cdots$, we have

$$
\|u(\cdot, t)\|_{K}^{2}=\int_{K}|u(x, t)|^{2} d x \leqq C_{m, f}(K)|t|^{-2 m}, \text { for all } t \neq 0
$$

The next result shows that if $B(t)$ is a set that does not increase "too" fast as $t \rightarrow \infty$, then the energy in $B(t)$ of non-static solution decays to zero.

Theorem 2. Let $\{B(t) \mid t>0\}$ be a family of bounded measurable sets in $\mathbf{R}^{n}$ such that, for some $\alpha<1, \Theta(t)=0\left(t^{\alpha}\right)$ as $t \rightarrow \infty$, where $\Theta(t)=\sup \{|x| \mid x \in B(t)\}$. Then,

$$
\lim _{t \rightarrow \infty}\|v(\cdot, t)\|_{B(t)}^{2}=\lim _{t \rightarrow \infty} \int_{B(t)}|v(x, t)|^{2} d x=0
$$

for any non-static solution $v(\cdot, t)$ with finite energy.

Proof. First, let us assume that $v(\cdot, t)$ is a non-static solution with initial value in $S$. Also, since $\alpha<1$, we can choose a positive integer $m$ such that

$$
\begin{equation*}
\alpha(2 m+n)<2 m \tag{10}
\end{equation*}
$$

Therefore, by Theorem 1, for some constant $C>0$, we have that

$$
|v(x, t)|^{2} \leqq C t^{-2 m}\left(1+|x|^{2 m}\right) \leqq C t^{-2 m}\left(1+\Theta(t)^{2 m}\right)
$$

for any $x \in B(t)$. It follows that

$$
\begin{aligned}
\int_{B(t)} \quad|v(x, t)|^{2} d x & \leqq C t^{-2 m}\left(1+\Theta(t)^{2 m}\right) \operatorname{meas}(B(t)) \\
& \leqq C^{\prime} t^{-2 m}\left(1+\Theta(t)^{2 m}\right) \Theta(t)^{n}
\end{aligned}
$$

and, since $\Theta(t)=0\left(t^{\alpha}\right)$, we obtain

$$
\int_{B(t)}|v(x, t)|^{2} d x \leqq C^{\prime \prime} t^{-2 m}\left(t^{\alpha n}+t^{\alpha(2 m+n)}\right)
$$

for $t$ large. Hence, in view of (10), the theorem is proved in this case.
Now, let $v(\cdot, t)$ be an arbitrary non-static solution with finite energy, that is, $v(\cdot, 0)=g \in N(A)^{\perp}$. Since $S$ is dense in $N(A)^{\perp}$, given any $\epsilon>$ 0 , there exists $g_{\epsilon} \in S$ such that $\left\|g-g_{\epsilon}\right\|<\epsilon$. And, since $\|U(t)\|=1$, we obtain

$$
\begin{aligned}
\|v(\cdot, t)\|_{B(t)} & =\|U(t) g\|_{B(t)} \leqq \| U(t)\left(g-g_{\epsilon}\left\|_{B(t)}+\right\| U(t) g_{\epsilon} \|_{B(t)}\right. \\
& \leqq\left\|U(t)\left(g-g_{\epsilon}\right)\right\|+\left\|U(t) g_{\epsilon}\right\|_{B(t)} \\
& =\left\|g-g_{\epsilon}\right\|+\left\|U(t) g_{\epsilon}\right\|_{B(t)}<\epsilon+\left\|U(t) g_{\epsilon}\right\|_{B(t)} .
\end{aligned}
$$

But, by what we just proved, $\lim _{t \rightarrow \infty}\left\|U(t) g_{\epsilon}\right\|_{B(t)}=0$. Hence, lim $\sup _{t \rightarrow \infty}\|v(\cdot, t)\|_{B(t)} \leqq \epsilon$, and the proof is complete since $\epsilon>0$ is arbitrary.

If we take $B(t)$ to be a fixed (bounded, measurable) set $B$ for all $t>$ 0 , Theorem 2 yields the usual local energy decay, that is, the decay of the energy in a fixed bounded measurable set:

Corollary. Given a bounded measurable set $B \subset \mathbf{R}^{n}$,

$$
\lim _{t \rightarrow \infty}\|v(\cdot, t)\|_{B}^{2}=\lim _{t \rightarrow \infty} \int_{B}|v(x, t)|^{2} d x=0
$$

for any non-static solution $v(\cdot, t)$ with finite energy.
Now, let us say that a family $\{B(t) \mid t>0\}$ of sets in $R^{n}$ converges to a set $B, B(t) \rightarrow B$, if the characteristic function $\chi_{B(t)}$ of $B(t)$ converges almost everywhere to the characteristic function $\chi_{B}$ of $B$, as $t \rightarrow \infty$. In
this case, it follows that the energy in $B(t)$ of any solution $u(\cdot, t)$ is asymptotically constant.
(We remark that the requirement that $\chi_{B(t)}(x) \rightarrow \chi_{B}(x)$ for all $x \in \mathbf{R}^{n}$ is equivalent to the set-theoretical notion of limit $\cup_{s>0} \cap_{t \geqq s} B(t)=$ $\cap_{s>0} \cup_{t \geqq s} B(t)=B$.)

Theorem 3. In addition to the hypotheses of Theorem 2, assume that $B(t) \rightarrow B$, where $B$ is any measurable set, not necessarily bounded. Then, for any solution $u(\cdot, t)=U(t) f$ with finite energy,

$$
\lim _{t \rightarrow \infty}\|u(\cdot, t)\|_{B(t)}^{2}=E_{B}^{o}[f]
$$

exists.
Proof. Let $u(\cdot, t)=U(t) f$ be a solution with finite energy. We can decompose $f \in H$ uniquely as a sum $f=h+g$, with $h \in N(A)$ and $g$ $\in N(A)^{\perp}$. We observe that if $u(\cdot, t)$ has no static part (that is, $h=0$ ), the proof is that of Theorem 2 and does not depend on the additional assumption that $B(t) \rightarrow B$. In any case, $U(t) h=h$ for all $t$ so that

$$
U(t) f=h+U(t) g .
$$

And, since $U(t) f-h$ is non-static, Theorem 2 gives $\lim _{t \rightarrow \infty} \| U(t) f-$ $h \|_{B(t)}=0$, hence

$$
\lim _{t \rightarrow \infty}\left(\|u(\cdot, t)\|_{B(t)}-\|h\|_{B(t)}\right)=0
$$

Now, the assumption $B(t) \rightarrow B$ implies (by the Lebesgue dominated convergence theorem) that

$$
\lim _{t \rightarrow \infty}\|h\|_{B(t)}=\|h\|_{B} \equiv\left(E_{B}^{o}[f]\right)^{1 / 2}
$$

The proof is complete.
4. Final remarks. 1) As already mentioned in the introduction, symmetric hyperbolic systems (1) have been studied by other authors under additional conditions on the matrices $A_{j}$. Results on decay of solutions can be found in [4], [7], for example. In [8], the case of the wave equation is thoroughly investigated. Uniform (over all $x$-space) decay results can be found in [2], where a further assumption of convexity on the connected sheets of the "wave surface" $P(1, p)=0$ is imposed (see also [6], where the wave equation and the Klein-Gordon equation are considered).

One approach to studying decay of solutions is clearly through the so called Riemann matrix of (1), that is, the distribution matrix-valued solution $R(x, t)$ of (1) with initial value $R(x, 0)=\delta(x) I_{n}$ : any smooth solu-
tion $u(x, t)$ with initial value $f(x)$ can then be written as $u(x, t)=[R(\cdot$, $t) * f](x)$. Since $R(x, t)$ is homogeneous of degree $-n, R(x, t)=$ $t^{-n} R(x / t, 1)$, and it could be thought that any (non-static) smooth solution would decay at every point $x$ at least like $t^{-n}$. This is certainly not the case, due to the singularities of $R(x, t)$. In fact, concrete examples in two and three dimensions have been given where solutions decay only as $t^{-1 / 2}([5])$. We conjecture that there is no rate of decay for the general symmetric hyperbolic systems considered in this paper. In other words, we conjecture that there does not exist a function $\varphi(t) \rightarrow 0$ as $t$ $\rightarrow \infty$, with the property that for each non-static solution $u$, the inequality $\sup _{x}|u(x, t)| \leqq C \varphi(t), t>0$, holds with some constant $C=$ $C(u)$.
2) Given a symmetric hyperbolic system for which the null space of $A(p)$ has constant dimension for all $p \neq 0$, we can naturally associate a minimum positive speed of propagation, $\lambda_{\min }=\inf \left\{\left|\lambda_{j}(\omega)\right| \mid \omega \in S^{n-1}\right.$, $j=1, \cdots, r\}($ see $(6),(7))$. In this case, it is not hard to show that the energy of any non-static solution $u(\cdot, t)$ inside a cone $|x|=C t$ tends to zero as $t \rightarrow \infty$, provided that $C<\lambda_{\min }$. On the other hand, if $A(\omega)$ possesses characteristic roots that vanish for certain $\omega$; but not identically, then there is no minimum positive speed of propagation and the above result does not apply. However, since any paraboloid $|x|=C t^{\alpha}$, $\alpha<1$, has slope $d|x| / d t$ that approaches zero as $t \rightarrow \infty$, we see that Theorem 2 is natural counterpart of the above mentioned result on energy decay in a cone. We contend that we cannot allow $\alpha=1$.
3) Analogous results can be easily obtained for systems which are "perturbations" of "free" systems (1), as long as their solutions behave asymptotically as "free" solutions. For example, in [1] systems of the form $E(x) u_{t}+\sum_{j=1}^{n} A_{j} u_{x_{j}}=0$ are considered. It is shown there that if $\int\left(1+|x|^{2}\right)^{2}|E(x)-I|^{2} d x<\infty$ then, for any solution $u(\cdot, t)$ in a "certain class", there exists a (free) solution $u_{+}(\cdot, t)$ of $u_{t}+\sum_{j=1}^{n} A_{j} u_{x_{j}}$ $=0$ such that

$$
\lim _{t \rightarrow \infty}\left\|u(\cdot, t)-u_{+}(\cdot, t)\right\|=0
$$

For any such $u$ and for $B(t) \rightarrow B$, it clearly follows that

$$
\lim _{t \rightarrow \infty}\|u(\cdot, t)\|_{B(t)}=\left\|h_{+}\right\|_{B}
$$

where $h_{+}$is the orthogonal projection of $f_{+}=u_{+}(\cdot, 0)$ on $N(A)$.

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