# Atti Accademia Nazionale dei Lincei <br> Classe Scienze Fisiche Matematiche Naturali RENDICONTI 

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# Periodic or almost-periodic solutions of a non linear functional equation. Nota I 

Atti della Accademia Nazionale dei Lincei. Classe di Scienze Fisiche, Matematiche e Naturali. Rendiconti, Serie 8, Vol. 43 (1967), n.3-4, p. 161-167.

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#### Abstract

Analisi matematica. - Periodic or almost-periodic solutions of a non linear functional equation ${ }^{(*)}$. Nota $I^{(*)}$ di Giovanni Prouse, presentata dal Corrisp. L. Amerio.


Sunto. - Si studiano alcuni problemi di propagazione per un'equazione funzionale non lineare. Dopo un'analisi preliminare del problema di Cauchy, si danno, in particolare, teoremi di esistenza ed unicità delle soluzioni periodiche e quasi-periodiche.
i. In the present paper we shall study some evolution problems for the non linear functional equation

$$
\begin{equation*}
u^{\prime}(\eta)+\left(\mathrm{A}_{1}+\mathrm{A}_{3}\right) u(\eta)+\mathrm{BA}_{2} u(\eta)=f(\eta) \tag{I.I}
\end{equation*}
$$

where $A_{1}$ and $A_{2}$ are unbounded positive linear operators, $B$ is a "dissipative " non linear operator and $\mathrm{A}_{3}$ is a linear operator which is "smaller" (in an appropriate sense) than $A_{1}$. Moreover, $A_{1}$ and $A_{2}$ are permutable, while $\eta \in[0, T]$ or ( $-\infty,+\infty$ ).

Under the hypotheses we shall give later on in this paragraph, we shall prove an existence and uniqueness theorem of the solution (in the large) of the Cauchy problem, essential introduction to the study of periodic and almostperiodic (a.p.) solutions, which is the main aim of this paper.

An example of a type of equation for which the results we shall give hold is the functional equation

$$
\frac{\partial u(x, \eta)}{\partial \eta}-\mathrm{A} u(x, \eta)+\mathrm{BA}^{\sigma} u(x, \eta)=f(x, \eta) \quad\left(x=\left\{x_{1}, \cdots, x_{m}\right\}\right)
$$

where A is a self adjoint second order elliptic operator and $0 \leq \sigma<\mathrm{I}$; this example will be considered in detail and generalized in $\S 5$.

We wish to recall that Lions and Strauss [r] have proved an existence and uniqueness theorem of the solution of the Cauchy problem for the equation

$$
\begin{equation*}
u^{\prime}(\eta)+\mathrm{A}(\eta) u(\eta)+\mathrm{B}(\eta) u(\eta)=f(\eta), \tag{1.2}
\end{equation*}
$$

which, formally, contains (I.I). If $\mathrm{A}_{2}=\mathrm{I}$, then equation (I.I) is obviously a particular case of (I.2) and the existence and uniqueness theorem we shall prove is included in that of Lions and Strauss.

In the general case $\left(A_{2} \neq I\right)$, the proof of the existence and uniqueness of the solution of the Cauchy problem is analogous to that given in [I], but differs from it in the part that refers to the non linear operator $B$. One of the
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(**) Pervenuta all'Accademia il 23 settembre 1967.
fundamental hypotheses we shall make on B (analogous to the corresponding hypothesis made by Lions and Strauss) is, in fact, that of being " asymptotically monotone '"; while it may be observed that, in equation (I.2), it is always possible, by means of the transformation $u(\eta)=e^{k \eta} v(\eta)$, to assume that B is "strictly monotone ", this may not be done in the case of equation (I.I) (at least in the case we shall consider). The procedure followed by Lions and Strauss, based on the "strict monotonicity" of B will be substituted by a compactness method, direct consequence of Lemma 3, which we shall prove in § 2.

We observe that it is possible to write (I.I) in an equivalent form. Multiplying (I.I) by $\mathrm{A}_{2}$, setting $\mathrm{A}_{2} u(\eta)=v(\eta), \mathrm{A}_{2} f(\eta)=g(\eta)$ and bearing in mind that the operators $\mathrm{A}_{1}$ and $\mathrm{A}_{2}$ are permutable, we obtain

$$
\begin{equation*}
v^{\prime}(\eta)+\mathrm{A}_{1} v(\eta)+\mathrm{A}_{2} \mathrm{~A}_{3} \mathrm{G}_{2} v(\eta)+\mathrm{A}_{2} \mathrm{~B} v(\eta)=g(\eta), \tag{I.3}
\end{equation*}
$$

where $G_{2}$ is Green's operator corresponding to $A_{2}$.
Equation (I.3) is formally contained in (I.2) and the observations concerning (I.I) hold also for this equation. In what follows, we shall however always consider equation (I.I); the results we shall obtain could anyway be proved substituting (I.I) with (I.3).

Let us now give some fundamental definitions and hypotheses.
If $\mathscr{H}_{1}$ and $\mathscr{H}_{2}$ are any two Hilbert spaces, with $\mathscr{H}_{1} \subset \mathscr{H}_{2}{ }^{(1)}$ and dense in $\mathscr{H}_{2}$, we shall indicate with $\left[\mathcal{H}_{1}, \mathscr{H}_{2}\right]_{\theta}(\mathrm{O}<\theta<\mathrm{I})$ the Hilbert space $\mathscr{H}_{1}^{1-\theta} \mathscr{H}_{2}^{\theta}$, intermediate between $\mathscr{H}_{1}$ and $\mathscr{H}_{2}$, defined according to the procedure given by Lions in [2]. If K is a Banach space, we shall indicate with $\mathrm{L}^{p}$ (o, $\mathrm{T} ; \mathrm{K}$ ) ( $p \geq \mathrm{I}$ ) the space of functions $v(\eta)$ which are $\mathrm{L}^{p}$ over ( $\mathrm{O}, \mathrm{T}$ ) with values in $K$, i.e. such that

$$
\begin{equation*}
\|v\|_{\mathrm{L}^{p}(0, \mathrm{~T} ; \mathrm{K})}=\left\{\int_{0}^{\mathrm{T}}\|v(\eta)\|_{\mathrm{K}}^{p} d \eta\right\}^{1 / p}<+\infty . \tag{I.4}
\end{equation*}
$$

Let $\mathrm{V}_{1}, \mathrm{~V}_{2}, \mathrm{H}$ be three real Hilbert spaces, with $\mathrm{V}_{1} \subset \mathrm{~V}_{2} \subseteq \mathrm{H}$, the embedding of $\mathrm{V}_{1}$ in $\mathrm{V}_{2}$ being completely continuous and each space dense in the following one ${ }^{(2)}$. Moreover, let $a_{1}(u, v), a_{2}(u, v)$ be two bilinear hermitian forms such that there exist two positive constants, $c_{1}$ and $c_{2}$ for which

$$
\begin{align*}
\left|a_{i}(u, v)\right| & \leq c_{1}\|u\|_{\mathrm{V}_{i}}\|v\|_{\mathrm{V}_{i}} & & \forall u, v \in \mathrm{~V}_{i},  \tag{i.5}\\
a_{i}(v, v) & \geq c_{2}\|v\|_{\mathrm{v}_{i}}^{2} & & \forall v \in \mathrm{~V}_{i} \quad(i=\mathrm{I}, 2) .
\end{align*}
$$

(I) By the notation $\mathscr{H}_{1} \subset \mathscr{H}_{2}$ we mean that $\mathscr{H}_{1}$ is contained in $\mathscr{H}_{2}$ and that the inclusion mapping of $\mathscr{H}_{1}$ into $\mathscr{H}_{2}$ is continuous.
(2) The complete continuity of $\mathrm{V}_{1}$ in $\mathrm{V}_{2}$ is not essential for the existence and uniqueness of the solution of the Cauchy problem when $\mathrm{V}_{2}=\mathrm{H}$ (which corresponds to the case $\mathrm{A}_{2}=\mathrm{I}$ ). This hypothesis is, in fact, not necessary whenever the non linear operator B is "strictly monotone" as, in such a case, the procedure outlined by Lions and Strauss can be followed.

It is then possible to define in H two linear, closed, unbounded operators $\mathrm{A}_{1}$ and $\mathrm{A}_{2}$, with domain $\mathrm{D}\left(\mathrm{A}_{1}\right)$ and $\mathrm{D}\left(\mathrm{A}_{2}\right)$ respectively, in the following way. An element $u \in V_{i}$ belongs to $\mathrm{D}\left(\mathrm{A}_{i}\right)$ if the linear functional $v \rightarrow a_{i}(u, v)$ is continuous on $\mathrm{V}_{i}$ with the topology induced by H . The operators $\mathrm{A}_{i}$ are therefore defined by the relations

$$
\begin{equation*}
a_{i}(u, v)=\left(\mathrm{A}_{i} u, v\right)_{\mathrm{H}} \quad \forall v \in \mathrm{~V}_{i} \quad(i=\mathrm{I}, 2) \tag{I.7}
\end{equation*}
$$

and, for the hypotheses made, are linear, self-adjoint and positive. Defining in $D\left(A_{i}\right)$ the scalar product by the relation

$$
\begin{equation*}
(u, v)_{\mathrm{D}\left(\mathrm{~A}_{i}\right)}=\left(\mathrm{A}_{i} u, \mathrm{~A}_{i} v\right)_{\mathrm{H}}{ }^{(3)} \tag{1.8}
\end{equation*}
$$

$\mathrm{D}\left(\mathrm{A}_{i}\right)$ is a Hilbert space, dense in H .
It is then possible to introduce the operators $\mathrm{A}_{i}^{v}, \nu$ real $>0$, and it results

$$
\begin{equation*}
\left[\mathrm{D}\left(\mathrm{~A}_{i}^{r}\right), \mathrm{D}\left(\mathrm{~A}_{i}^{s}\right)\right]_{\theta}=\mathrm{D}\left(\mathrm{~A}_{i}^{r(1-\theta)+s \theta}\right) \tag{I.9}
\end{equation*}
$$

Moreover (Lions [3])

$$
\begin{equation*}
\mathrm{V}_{i}=\mathrm{D}\left(\mathrm{~A}_{i}^{1 / 2}\right)=\left[\mathrm{D}\left(\mathrm{~A}_{i}\right), \mathrm{H}\right]_{1 / 2} \tag{1.10}
\end{equation*}
$$

We shall therefore set, in accordance with (i.8) and (I.9),

$$
\begin{equation*}
(u, v)_{\mathrm{D}\left(\mathrm{~A}_{i}^{v}\right)}=\left(\mathrm{A}_{i}^{v} u, \mathrm{~A}_{i}^{v} v\right)_{\mathrm{H}} . \tag{I.II}
\end{equation*}
$$

By a density theorem of Lions [4], from (1.9) follows that the spaces $\mathrm{D}\left(\mathrm{A}_{i}^{v}\right)$ are dense in H .

Let us observe that the equations $\mathrm{A} u=f, \mathrm{~A} v=g$ have, if the hypotheses made on the forms $a_{i}(u, v)$ are verified, one and only one solution $\forall f, g \in \mathrm{H}$ given respectively by

$$
\begin{equation*}
u=\mathrm{G}_{1} f \in \mathrm{D}\left(\mathrm{~A}_{1}\right) \quad, \quad v=\mathrm{G}_{2} g \in \mathrm{D}\left(\mathrm{~A}_{2}\right), \tag{1:12}
\end{equation*}
$$

where $G_{1}$ and $G_{2}$ are Green's operators corresponding to $A_{1}$ and $A_{2}$ respectively. Furthermore, $\mathrm{D}\left(\mathrm{A}_{1}\right)=\mathrm{G}_{1} \mathrm{H}, \mathrm{D}\left(\mathrm{A}_{2}\right)=\mathrm{G}_{2} \mathrm{H}$.

Let $h$ be an arbitrary element $\in H$ and consider the equation

$$
\begin{equation*}
\mathrm{A}_{2} \mathrm{~A}_{1} u=h \tag{1.13}
\end{equation*}
$$

Setting $\tilde{u}=\mathrm{A}_{1} u$, it results, by the second of (I.I2),

$$
\tilde{u}=\mathrm{G}_{2} h
$$

and, consequently, by the first of (I.12),

$$
\begin{equation*}
u=\mathrm{G}_{1} \mathrm{G}_{2} h \tag{1.14}
\end{equation*}
$$

(3) The norm corresponding to ( I .8 ) is, under the hypotheses made, obviously equivalent to

$$
\|v\|_{\mathrm{H}}^{2}+\left\|\mathrm{A}_{i} v\right\|_{\mathrm{H}}^{2} .
$$

It is, therefore

$$
\begin{equation*}
\mathrm{D}\left(\mathrm{~A}_{2} \mathrm{~A}_{1}\right)=\mathrm{G}_{1} \mathrm{G}_{2} \mathrm{H} \tag{1.15}
\end{equation*}
$$

Correspondingly, if we consider the equation $\mathrm{A}_{1} \mathrm{~A}_{2} v=h$, we find that

$$
\begin{equation*}
v=\mathrm{G}_{2} \mathrm{G}_{1} h \quad, \quad \mathrm{D}\left(\mathrm{~A}_{1} \mathrm{~A}_{2}\right)=\mathrm{G}_{2} \mathrm{G}_{1} \mathrm{H} \tag{1.16}
\end{equation*}
$$

From (I.I4), (I.I5), (I.I6) it follows that, if the operators $G_{1}$ and $G_{2}$ are permutable, also the operators $A_{1}$ and $A_{2}$ are permutable, being

$$
\begin{equation*}
\mathrm{D}\left(\mathrm{~A}_{1} \mathrm{~A}_{2}\right)=\mathrm{D}\left(\mathrm{~A}_{2} \mathrm{~A}_{1}\right) \quad, \quad \mathrm{A}_{1} \mathrm{~A}_{2} u=\mathrm{A}_{2} \mathrm{~A}_{1} u \quad \forall u \in \mathrm{D}\left(\mathrm{~A}_{1} \mathrm{~A}_{2}\right) . \tag{I.17}
\end{equation*}
$$

If, on the other hand, relations (I.I7) hold, then, by (I.I4), (I.16),

$$
\begin{equation*}
\mathrm{G}_{1} \mathrm{G}_{2} h=u=\mathrm{G}_{2} \mathrm{G}_{1} h \quad \forall h \in \mathrm{H} \tag{1.18}
\end{equation*}
$$

i.e. the operators $G_{1}$ and $G_{2}$ are permutable.

All that has been said above applies also to the operators $A_{i}^{v}$.
Let us now assume that the operators $\mathrm{A}_{1}$ and $\mathrm{A}_{2}$ satisfy the following assumptions:
I) There exists a number $\sigma$, with $0 \leq \sigma<\mathrm{I}$, such that $\mathrm{D}\left(\mathrm{A}_{1}^{\sigma}\right) \subseteq \mathrm{D}\left(\mathrm{A}_{2}\right)$, being moreover $\mathrm{D}\left(\mathrm{A}_{1}^{\sigma}\right)$ dense in $\mathrm{D}\left(\mathrm{A}_{2}\right)$;
II) The operators $\mathrm{A}_{1}$ and $\mathrm{A}_{2}$ are permutable, i.e. relations (I.17) hold. Conditions I) and II) are, for instance, satisfied, if $\mathrm{A}_{2}=\mathrm{A}_{1}^{\sigma}$.

Let $\mathrm{A}_{3}$ be a continuous, linear operator from $\mathrm{D}\left(\mathrm{A}_{1}\right)$ to H ; we shall assume that:
III) There exists a positive number $\alpha$ such that, setting $\mathrm{W}=\mathrm{D}\left(\left(\mathrm{A}_{1} \mathrm{~A}_{2}\right)^{1 / 2}\right)$, it results

$$
\begin{equation*}
\left(\left(\mathrm{A}_{1}+\mathrm{A}_{3}\right) v, \mathrm{~A}_{2} v\right)_{\mathrm{H}} \geq \alpha\|v\|_{\mathrm{W}}^{2} \quad \forall v \in \mathrm{D}\left(\mathrm{~A}_{1} \mathrm{~A}_{2}\right) ; \tag{I.19}
\end{equation*}
$$

IV) $A_{3} W \subseteq A_{1} W$.

In what follows, if K is a Banach space CH and dense in H , we shall denote by $\mathrm{K}^{\prime}$ its dual space, while the space $\mathrm{H}^{\prime}$ will be identified with H .

Let B be a non linear operator from a reflexive Banach space YCH, dense in H , to $\mathrm{Y}^{\prime}$, with $\mathrm{BO}=\mathrm{o}$; we assume that the following hypotheses hold:
$\mathrm{V}) \mathrm{B}$ is weakly continuous from finite-dimensional subspaces of Y to $\mathrm{Y}^{\prime}$;
VI) If $v(\eta) \in \mathrm{L}^{p}(\mathrm{o}, \mathrm{T} ; \mathrm{Y})(p>2), \quad \mathrm{B} v(\eta) \in \mathrm{L}^{p^{\prime}}\left(\mathrm{o}, \mathrm{T} ; \mathrm{Y}^{\prime}\right)\left(\mathrm{I} / p+\mathrm{I} / p^{\prime}=\mathrm{I}\right)$; moreover, the map $v(\eta) \rightarrow \mathrm{B} v(\eta)$ sends bounded sets of $\mathrm{L}^{p}(\mathrm{o}, \mathrm{T} ; \mathrm{Y})$ into bounded sets of $\mathrm{L}^{p^{\prime}}\left(\mathrm{o}, \mathrm{T} ; \mathrm{Y}^{\prime}\right)$;
VII) If $\left\{v_{n}\right\}$ is an $\mathrm{L}^{p}(\mathrm{o}, \mathrm{T} ; \mathrm{Y})$-bounded sequence, such that

$$
\begin{equation*}
\lim _{n \rightarrow \infty} v_{\left.n L_{L^{2}(0, \mathrm{~T}} ; \mathrm{H}\right)}^{=} v, \tag{1.20}
\end{equation*}
$$

it is possible to extract a subsequence $\left\{v_{n^{\prime}}\right\} \subset\left\{v_{n}\right\}$ such that

$$
\begin{equation*}
\lim _{n^{\prime} \rightarrow \infty}^{*} \mathrm{~B} v_{n^{\prime}} \underset{\mathrm{L}^{p^{\prime}}\left(0, \mathrm{~T} ; \mathrm{Y}^{\prime}\right)}{ }=\mathrm{B} v^{(4)} . \tag{I.2I}
\end{equation*}
$$

(4) We shall denote by lim* the limit in the weak topology.
VIII) There exist three positive constants, $c_{3}, c_{4}, c_{5}$ such that

$$
\begin{array}{cl}
\langle\mathrm{B} v, v\rangle \geq c_{3}\|v\|_{\mathrm{Y}}^{p}-c_{4} & \forall v \in \mathrm{Y}, \\
\left\langle\mathrm{~B} v_{1}-\mathrm{B} v_{2}, v_{1}-v_{2}\right\rangle \geq-c_{5}\left\|v_{1}-v_{2}\right\|_{\mathrm{H}}^{2} & \forall v_{1}, v_{2} \in \mathrm{Y} . \tag{I.23}
\end{array}
$$

Finally, we shall assume that the spaces introduced satisfy the following conditions:
IX) The space $\left[\mathrm{D}\left(\mathrm{A}_{1}\right), \mathrm{D}\left(\mathrm{A}_{2}\right)^{\prime}\right]_{1 / 2} \cap \mathrm{Y}$ is separable and dense in $\left[\mathrm{D}\left(\mathrm{A}_{1}\right), \mathrm{D}\left(\mathrm{A}_{2}\right)^{\prime}\right]_{1 / 2}$ and in Y ;

X ) There exists a number $s>0$ such that $\mathrm{Y}^{\prime} \subset \mathrm{D}\left(\mathrm{A}_{1}^{s}\right)^{\prime}$.
We observe that hypotheses III), IV), IX) correspond to those given by Lions and Strauss and also hypotheses V), VI), VIII) are analogous to the conditions imposed by these Authors on the non linear operator B. Hypothesis II) (which is automatically verified if $\mathrm{A}_{2}=\mathrm{I}$ ) is, on the other hand, essential for the validity of all the lemmas and theorems we shall give. Finally, hypotheses VII) and X) will be used in the proof of lemma 3 which, as we have already mentioned, states a compactness property of the solutions. If, therefore, in relation (1.23), it were $c_{5}=0$, (i.e. if B were " strictly monotone '") it would be possible, in the proof of the existence and uniqueness theorem of the solution of the Cauchy problem, to utilize the procedure given by Lions and Strauss, based on the " strict monotonicity " of B (see footnote (2)); hypotheses VII) and X) would then no longer be necessary.

Let us now consider the equation

$$
\begin{equation*}
u^{\prime}(\eta)+\left(\mathrm{A}_{1}+\mathrm{A}_{3}\right) u(\eta)+\mathrm{BA}_{2} u(\eta)=f(\eta) \quad(0 \leq \eta \leq \mathrm{T}) \tag{..24}
\end{equation*}
$$

where $f$ takes its values in $\mathrm{Y}^{\prime}$. Setting $\mathrm{Z}=\left[\mathrm{D}\left(\mathrm{A}_{2}\right), \mathrm{D}\left(\mathrm{A}_{1}\right)^{\prime}\right]_{1 / 2}$, we shall say that $u(\eta)$ is a solution in $[\mathrm{O}, \mathrm{T}]$ of (1.24) if:
a) $u \in \mathrm{~L}^{\infty}\left(\mathrm{o}, \mathrm{T} ; \mathrm{V}_{2}\right) \cap \mathrm{L}^{2}(\mathrm{o}, \mathrm{T} ; \mathrm{W}), \mathrm{A}_{2} u \in \mathrm{~L}^{p}(\mathrm{o}, \mathrm{T} ; \mathrm{Y})(p>2)$;
b) $u(\eta)$ satisfies (I.24) almost everywhere in ( $\mathrm{O}, \mathrm{T}$ ), in the sense of distributions with values in $\mathrm{Z}+\mathrm{Y}^{\prime}$.

As will be seen, the hypotheses made above are sufficient, if $f \in \mathrm{~L}^{p^{\prime}}\left(\mathrm{o}, \mathrm{T} ; \mathrm{Y}^{\prime}\right)$, to guarantee the existence and uniqueness of the solution of the Cauchy problem and the existence of a periodic solution (if $f(\boldsymbol{\eta})$ is periodic). In order to prove the uniqueness of the periodic solution and the existence and uniqueness of an a.p. solution (assuming that $f(\eta)$ is a.p.) it will be necessary to make the following further hypothesis:
XI) If $\gamma$ is the embedding constant of W in $\mathrm{D}\left(\mathrm{A}_{2}\right)$, it results $c_{5}<\frac{\alpha}{\gamma^{2}}$, where' $\alpha$ and $\mathrm{c}_{5}$ are the constants appearing in relations (1.19) and (1.23).

Let us enounce the theorems which, together with some auxiliary lemmas, will be proved in detail in the following paragraphs.

Theorem i: Assume that all the hypotheses made above (with the exception of XI$)$ ) are verified and that $f \in \mathrm{~L}^{p^{\prime}}\left(\mathrm{o}, \mathrm{T} ; \mathrm{Y}^{\prime}\right)$. There exists then in $[\mathrm{o}, \mathrm{T}]$ one, and only one, solution of (1.24) satisfying the initial condition

$$
\begin{equation*}
u(0)=u_{0} \tag{1.25}
\end{equation*}
$$

with $u_{0}$ arbitrary element of $\mathrm{V}_{2}$.
Theorem 2: If the hypotheses of theorem I are satisfied and $f(\eta)$ is periodic with period T , there exists at least one solution which is periodic with period T .

Theorem 3: Assume that all the hypotheses made (except XI)) are verified and that ${ }^{(5)}$

$$
\begin{align*}
\operatorname{Sup}_{t \in \mathrm{~J}}\left\{\int_{0}^{1}\|f(t+\eta)\|_{\mathrm{y}^{\prime}}^{p^{\prime}} d \eta\right\}^{1 / p^{\prime}}=\operatorname{Sup}_{t \in \mathrm{~J}}\|f(t)\|_{\mathrm{L}^{p^{\prime}}\left(0,1 ; \mathrm{v}^{\prime}\right)}= & \mathrm{M}_{0}<+\infty  \tag{I.26}\\
& (\mathrm{J}=(-\infty,+\infty)) .
\end{align*}
$$

There exists then in J at least one solution $\tilde{u}(\eta)$ such that

$$
\begin{array}{r}
\operatorname{Sup}_{\eta \in \mathrm{J}}\|\tilde{u}(\eta)\|_{\mathrm{V}_{2}}=\mathrm{M}_{1}<+\infty \quad, \quad \operatorname{Sup}_{t \in \mathrm{~J}}\|\tilde{u}(t)\|_{\mathrm{L}^{2}(0,1 ; \mathrm{W})}=\mathrm{M}_{2}<+\infty,  \tag{1.27}\\
\operatorname{Sup}_{t \in \mathrm{~J}}\left\|\mathrm{~A}_{2} \tilde{u}(t)\right\|_{\mathrm{L}^{p}(0,1 ; \mathrm{Y})}=\mathrm{M}_{3}<+\infty, \\
\operatorname{Sup}_{t \in \mathrm{~J}}\left\|\mathrm{~A}_{2} \tilde{u}(t)\right\|_{\mathrm{H}^{\varepsilon}\left(0,1 ; \mathrm{D}\left(\mathrm{~A}_{1}^{\mathrm{\varepsilon}}\right)\right)}=\mathrm{M}_{4}<+\infty \\
(\varepsilon>0) .
\end{array}
$$

Moreover, if also condition XI) holds, then $\tilde{u}(t)$ is the only $\mathrm{L}^{2}\left(\mathrm{O}, \mathrm{I} ; \mathrm{V}_{2}\right)-$ bounded solution in J and, if $v(\eta)$ is any other solution corresponding to the same known term, it results

$$
\begin{equation*}
\lim _{\eta \rightarrow+\infty}\|\tilde{u}(\eta)-v(\eta)\| \mathrm{v}_{2}=0 . \tag{I.28}
\end{equation*}
$$

ThEOREM 4: Assume that all hypotheses made hold and that $f(t)$ is $\mathrm{L}^{p^{\prime}}\left(\mathrm{o}, \mathrm{I} ; \mathrm{Y}^{\prime}\right)$-weakly a.p. Then the bounded solution $\tilde{u}(t)$ (which, by theorem 3, exists and is unique) is $\mathrm{L}^{2}\left(\mathrm{o}, \mathrm{I} ; \mathrm{V}_{2}\right)-a . p$.

Theorem 5: If the assumptions of theorem 4 are verified and $f(t)$ is $\mathrm{L}^{p^{\prime}}\left(\mathrm{O}, \mathrm{I} ; \mathrm{Y}^{\prime}\right)-a . p$. , then $\tilde{u}(\eta)$ is $\mathrm{V}_{2}-a . p$. and $\tilde{u}(t)$ is $\mathrm{L}^{2}(\mathrm{O}, \mathrm{I} ; \mathrm{W})-a . p$.

As will be shown in $\S 5$, the theorems given above can, in particular, be applied to the following example.

Let $\Omega$ be an open, bounded set of the Euclidean space $\mathrm{S}^{m}(x=$ $\left.=\left\{x_{1}, x_{2}, \cdots, x_{m}\right\}\right)$, with boundary $\Gamma$. Using the notations introduced at the beginning of this $\S$, we set $H=L^{2}(\Omega)=L^{2}, V_{1}=H_{0}^{1}(\Omega)=H_{0}^{1}$ and

$$
\begin{array}{lll}
\mathrm{A}_{1}=\sum_{j, k=1}^{m} a_{j k} \frac{\partial^{2}}{\partial x_{j} \partial x_{k}} & , \quad \mathrm{~A}_{3}=\sum_{j=1}^{m} a_{j}(x) \frac{\partial}{\partial x_{j}}+a_{0}(x),  \tag{1.29}\\
\mathrm{E}=\sum_{j, k=1}^{m} e_{j k} \frac{\partial^{2}}{\partial x_{j} \partial x_{k}}+e_{0} & , \quad \mathrm{~A}_{2}=\mathrm{E}^{\sigma} \quad & (0 \leq \sigma<\mathrm{I}),
\end{array}
$$

(5) In what follows, we shall set $v(t)=\{v(t+\eta) ; \eta \in[0, \mathrm{I}]\}$ and, consequently,

$$
\|v(t)\|_{L^{q}(0,1 ; K)}=\left\{\int_{0}^{1}\|v(t+\eta)\|_{\mathrm{K}}^{q} d \eta\right\}^{1 / q} .
$$

with $a_{j k}=a_{k j}, e_{j k}=e_{k j}, a_{j}(x) \in \mathrm{L}^{\infty}(\Omega)$. We assume moreover that, $\forall v \in \mathrm{H}_{0}^{1}$

$$
\begin{equation*}
\left\langle\mathrm{A}_{1} v, v\right\rangle \geq \mu_{1}\|v\|_{\mathrm{H}_{0}^{1}}^{2} \quad, \quad\langle\mathrm{E} v, v\rangle \geq \mu_{2}\|v\|_{\mathrm{H}_{0}^{1}}^{2} \quad\left(\mu_{1}, \mu_{2}>0\right) . \tag{I.30}
\end{equation*}
$$

Let $\beta(\xi)$ be a continuous function, satisfying locally a Lipschitz condition, of the real variable $\xi$, defined in $J$ and such that $\beta(0)=0$. We shall suppose that there exists a number $p>2$ and three positive constants $c_{6}, c_{7}, c_{8}$ such that

$$
\begin{array}{ll}
c_{6}|\xi|^{p} \leq \xi \beta(\xi) \leq c_{7}|\xi|^{p} & \text { when }|\xi| \geq \bar{\xi} \\
\frac{\beta\left(\xi_{1}\right)-\beta\left(\xi_{2}\right)}{\xi_{1}-\xi_{2}} \geq-c_{8} & \forall \text { couple of points } \xi_{1}, \xi_{2} \in J \tag{I.32}
\end{array}
$$

We consider the equation

$$
\begin{equation*}
u^{\prime}(\eta)+\left(\mathrm{A}_{1}+\mathrm{A}_{3}\right) u(\eta)+\beta\left(\mathrm{E}^{\sigma} u(\eta)\right)=f(\eta) \tag{I.33}
\end{equation*}
$$

having set $u(\eta)=\{u(x, \eta) ; x \in \Omega\}, u^{\prime}(\eta)=\left\{\frac{\partial u(x, \eta)}{\partial \eta} ; x \in \Omega\right\}$, $\mathrm{A} u(\eta)=$ $=\{\mathrm{A} u(x, \eta) ; x \in \Omega\}, \quad \beta\left(\mathrm{E}^{\sigma} u(\eta)\right)=\left\{\beta\left(\mathrm{E}^{\sigma} u(x, \eta)\right) ; x \in \Omega\right\}, \quad f(\eta)=$ $=\{f(x, \eta) ; x \in \Omega)\}$.

This equation is obviously a special case of (1.24) and hypotheses I) . . X X are verified (see §5) provided the coefficients $a_{j}(x)(j=0, \mathrm{I}, \cdots, m)$ are "sufficiently small" and $\Gamma$ is of class $\mathrm{C}^{2 s}$, with $s=\left[\frac{m(p-2)}{4 p}\right]+\mathrm{I}$.

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