# Exact Controllability and Stabilization. The Multiplier Method 

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## PREFACE

This book grew out of a series of lectures given over the past four years in France, Hungary and the USA.

In the first part exact boundary controllability problems are studied by the Hilbert Uniqueness Method. This approach, introduced by Lions [3] in 1986, is based on uniqueness theorems leading to the construction of suitable Hilbert spaces of the controllable spaces. It is closely related to a duality theory of Dolecki and Russell [1]. Following Ho [1] and Lions [2], these spaces may often be characterized by using the multiplier method. In chapters 2 to 4 we reproduce some results of Lions [4], [5] with certain changes :

- some compactness-uniqueness arguments are replaced by constructive proofs;
- equations containing lower-order terms are also considered;
- more general boundary conditions are used which in fact simplify the theory.

The results of chapters 5 and 6 were obtained after the publication of Lions' monography. In chapter 5 we develop a general and constructive approach to improve the usual estimates of the exact controllability time. It was inspired by a new estimation method of Haraux [3]. Using this approach, in chapter 6 we improve most of the results obtained in chapters 3 and 4 . We also give elementary and constructive proofs for certain results of Zuazua [1], obtained earlier by indirect arguments.

The second part of the book is devoted to stabilizability. In chapters 8 and 9 strong and uniform boundary stabilization theorems are proved. Our method is a modified and simplified version of a Liapunov type approach introduced in Komornik and Zuazua [1]. We also present here a classical principle of Russell [2] connecting the exact controllability to the stabilizability, and some recent results of Conrad and Rao [1].

For the sake of brevity in the first nine chapters we consider only the wave equation, Maxwell's equations and very simple plate models. In the last chapter we consider the internal stabilization of the Korteweg-de Vries equation : we prove a special case of a theorem in Komornik, Russell and Zhang [2].

The multiplier method, applied systematically in this book, is remarkably elementary and efficient. In the bibliography we have included some references
which use other approaches : see in particular the work of Bardos et al., Littmann, Russell, Joó and their references. We have also included some material concerning other equations.

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## 0. Introduction. Vibrating strings

Let $I=(a, b)$ be a bounded interval, $T$ a positive number and consider the following problem, modelling among other things the small transversal vibrations of a string :

$$
\begin{gather*}
\quad\left(u_{t t}-u_{x x}\right)(x, t)=0, \quad(x, t) \in I \times(0, T)  \tag{1}\\
u(a, t)=v_{a}(t) \quad \text { and } \quad u(b, t)=v_{b}(t), \quad t \in[0, T]  \tag{2}\\
u(x, 0)=u^{0}(x) \quad \text { and } \quad u_{t}(x, 0)=u^{1}(x), \quad x \in I . \tag{3}
\end{gather*}
$$

The problem (1)-(3) is said to be exactly controllable if for "arbitrarily" given initial "state" $\left(u^{0}, u^{1}\right)$ there exist suitable "control" functions $v_{a}$ and $v_{b}$ such that the solution of (1)-(3) satisfies

$$
\begin{equation*}
u(x, T)=u_{t}(x, T)=0, \quad x \in I . \tag{4}
\end{equation*}
$$

We say that the controls $v_{a}$ and $v_{b}$ drive the system to rest in time $T$.
Naturally, we have to specify the functional spaces of the initial states and of the controls; the results depend on these choices.

The solution of (1)-(3) is by definition a function

$$
\begin{equation*}
u \in C\left([0, T] ; H^{1}(I)\right) \cap C^{1}\left([0, T] ; L^{2}(I)\right) \tag{5}
\end{equation*}
$$

satisfying (1) in the distributional sense, the equalities (2) pointwise, and the equalities (3) almost everywhere. (As for the usual properties of the Sobolev spaces applied in this book we refer e.g. to Lions and Magenes [1].)

We have the following result :
Theorem 0.1. - Let $T=b-a$ and let $\left(u^{0}, u^{1}\right) \in H^{1}(I) \times L^{2}(I)$ be such that

$$
\begin{equation*}
u^{0}(a)+u^{0}(b)+\int_{a}^{b} u^{1}(s) \mathrm{ds}=0 \tag{6}
\end{equation*}
$$

Then there is a unique choice of functions

$$
\begin{equation*}
v_{a}, v_{b} \in H^{1}(0, T) \tag{7}
\end{equation*}
$$

such that the solution of (1)-(3) satisfies (4).
Moreover, $v_{a}$ and $v_{b}$ are given by the formulae

$$
\begin{equation*}
2 v_{a}(t)=u^{0}(a+t)+u^{0}(a)+\int_{a}^{a+t} u^{1}(s) \mathrm{ds} \tag{8}
\end{equation*}
$$

$$
\begin{equation*}
2 v_{b}(t)=u^{0}(b-t)+u^{0}(b)+\int_{b-t}^{b} u^{1}(s) \mathrm{ds} \tag{9}
\end{equation*}
$$

and the solution $u$ has the following supplementary property:

$$
\begin{equation*}
u(a, t)+u(b, t)+\int_{a}^{b} u_{t}(s, t) \mathrm{ds}=0, \quad \forall t \in[0, T] \tag{10}
\end{equation*}
$$

Proof. - Applying d'Alembert's formula the solutions of (1) may be written in the form

$$
\begin{equation*}
u(x, t) \equiv f(x+t)+g(x-t) \tag{11}
\end{equation*}
$$

with suitable functions $f:(a, b+T) \rightarrow \mathbb{R}$ and $g:(a-T, b) \rightarrow \mathbb{R}$. Using (2) and (3) we obtain that

$$
\begin{gather*}
f(a+t)+g(a-t)=v_{a}(t), \quad t \in(0, T),  \tag{12}\\
f(b+t)+g(b-t)=v_{b}(t), \quad t \in(0, T),  \tag{13}\\
f(x)+g(x)=u^{0}(x), \quad x \in I,  \tag{14}\\
f^{\prime}(x)-g^{\prime}(x)=u^{1}(x), \quad x \in I . \tag{15}
\end{gather*}
$$

We deduce from (14), (15) that

$$
\begin{equation*}
2 f(x)=u^{0}(x)+U^{1}(x), \quad x \in I \tag{16}
\end{equation*}
$$

and

$$
\begin{equation*}
2 g(x)=u^{0}(x)-U^{1}(x), \quad x \in I \tag{17}
\end{equation*}
$$

where $U^{1}$ is a suitable primitive of $u^{1}$.
We conclude from (12) and (16) that

$$
2 g(a-t)=2 v_{a}(t)-u^{0}(a+t)-U^{1}(a+t), \quad 0<t<b-a
$$

similarly, we conclude from (13) and (17) that

$$
2 f(b+t)=2 v_{b}(t)-u^{0}(b-t)+U^{1}(b-t), \quad 0<t<b-a ;
$$

we can rewrite these relations in the following form :

$$
\begin{array}{ll}
2 g(s)=2 v_{a}(a-s)-\left(U^{1}+u^{0}\right)(2 a-s), & 2 a-b<s<a, \\
2 f(s)=2 v_{b}(s-b)+\left(U^{1}-u^{0}\right)(2 b-s), & b<s<2 b-a . \tag{19}
\end{array}
$$

We deduce from (11), (18) and (19) that for

$$
\begin{equation*}
a<x<b \quad \text { and } \quad \max \{x-a, b-x\}<t<b-a \tag{20}
\end{equation*}
$$

we have

$$
\begin{gather*}
2 u(x, t)=2 v_{a}(a+t-x)+2 v_{b}(x+t-b) \\
-u^{0}(2 a-x+t)-u^{0}(2 b-x-t)+\int_{2 a-x+t}^{2 b-x-t} u^{1}(s) \mathrm{ds} \tag{21}
\end{gather*}
$$

We have in particular

$$
u(x, b-a)=v_{a}(b-x)+v_{b}(x-a)-u^{0}(a+b-x), \quad x \in I
$$

and

$$
u_{t}(x, b-a)=v_{a}^{\prime}(b-x)+v_{b}^{\prime}(x-a)-u^{1}(a+b-x), \quad x \in I
$$

Thus the final conditions (4) are equivalent to

$$
v_{a}(b-x)+v_{b}(x-a)=u^{0}(a+b-x), \quad x \in I
$$

and

$$
-v_{a}(b-x)+v_{b}(x-a)=-V^{1}(a+b-x), \quad x \in I
$$

with an arbitrary primitive $V^{1}$ of $u^{1}$ whence

$$
\begin{array}{ll}
2 v_{a}(b-x)=\left(u^{0}+V^{1}\right)(a+b-x), & x \in I \\
2 v_{b}(x-a)=\left(u^{0}-V^{1}\right)(a+b-x), & x \in I
\end{array}
$$

or

$$
\begin{array}{ll}
2 v_{a}(t)=\left(u^{0}+V^{1}\right)(a+t), & 0<t<b-a, \\
2 v_{b}(t)=\left(u^{0}-V^{1}\right)(b-t), & 0<t<b-a . \tag{23}
\end{array}
$$

Let us introduce the subsets $U, R, D, L$ of $I \times(0, T)$ ( the letters mean "Up", "Right", "Down", "Left") defined by

$$
\begin{aligned}
U:= & \{(x, t) \in I \times(0, T): t>x-a \quad \text { et } t>b-x\}, \\
& R:=\{(x, t) \in I \times(0, T): b-x<t<x-a\}, \\
D:= & \{(x, t) \in I \times(0, T): t<x-a \text { et } t<b-x\}, \\
& L:=\{(x, t) \in I \times(0, T): x-a<t<b-x\} .
\end{aligned}
$$

One can easily deduce from (11), (16)-(19), (22) and (23) the following formulae :

$$
2 u(x, t)= \begin{cases}0, & \text { if }(x, t) \in U  \tag{24}\\ \left(u^{0}-V^{1}\right)(x-t), & \text { if }(x, t) \in R \\ u^{0}(x+t)+u^{0}(x-t)+\int_{x-t}^{x+t} u^{1}(s) \mathrm{ds}, & \text { if }(x, t) \in D \\ \left(u^{0}+V^{1}\right)(x+t), & \text { if }(x, t) \in L\end{cases}
$$

Comparing the limits of these formulae at $t=x-a$ and $t=b-x$ we obtain that for any fixed $0<t<T, t \neq T / 2$, the function $x \mapsto u(x, t)$ is continuous only if

$$
\begin{equation*}
V^{1}(a)=u^{0}(a) \quad \text { and } \quad V^{1}(b)=-u^{0}(b) . \tag{25}
\end{equation*}
$$

By hypothesis (6) $u^{1}$ has a unique primitive $V^{1}$ with this property. Choosing $V^{1}$ in this way, (22), (23) follow from (8), (9), and we deduce from (24) that

$$
2 u(x, t)= \begin{cases}0, & \text { if }(x, t) \in U  \tag{26}\\ u^{0}(x-t)+u^{0}(b)+\int_{x-t}^{b} u^{1}(s) \mathrm{ds}, & \text { if }(x, t) \in R \\ u^{0}(x+t)+u^{0}(x-t)+\int_{x-t}^{x+t} u^{1}(s) \mathrm{ds}, & \text { if }(x, t) \in D \\ u^{0}(x+t)+u^{0}(a)+\int_{a}^{x+t} u^{1}(s) \mathrm{ds}, & \text { if }(x, t) \in L\end{cases}
$$

using these formulae one can directly verify (5) and (10).
Remark 0.2. - The problem (1)-(3) is also exactly controllable if $T>b-a$. Indeed, it is sufficient to drive the system to rest in time $b-a$ and then to extend the control functions $v_{a}$, $v_{b}$ by zero for $b-a<t<T$. Since $v_{a}(b-a)=v_{b}(b-a)=0$ by (6), (8) and (9), the property (7) remains valid.

Remark 0.3. - The problem (1)-(3) is not exactly controllable if $T<b-a$. Indeed, assume that the system may be driven to rest in time $T$ from some given initial data $u^{0}, u^{1}$. Extending the corresponding controls $v_{a}, v_{b}$ by zero for $T<t<b-a$, we obtain that the system may be driven from the initial state $\left(u^{0}, u^{1}\right)$ to rest in time $b-a$ with controls $v_{a}, v_{b} \in L^{2}(0, b-a)$ satisfying

$$
\begin{equation*}
v_{a}=v_{b} \equiv 0 \quad \text { in some left neighbourhood of } b-a . \tag{27}
\end{equation*}
$$

On the other hand, the proof of theorem 0.1 (see (22) and (23)) shows that

$$
\begin{equation*}
2 v_{a}(t)+2 v_{b}(t)=u^{0}(a+t)+u^{0}(b-t)+\int_{b-t}^{a+t} u^{1}(s) \mathrm{ds}, \quad 0<t<b-a . \tag{28}
\end{equation*}
$$

But (27) and (28) are not satisfied simultaneously for all initial data $\left(u^{0}, u^{1}\right) \in H^{1}(I) \times L^{2}(I)$ satisfying (6).

The following result improves theorem 0.1 by driving the system to general final states :

Theorem 0.4. - Let $T=b-a,\left(u^{0}, u^{1}\right),\left(u_{T}^{0}, u_{T}^{1}\right) \in H^{1}(I) \times L^{2}(I)$ and assume that

$$
u^{0}(a)+u^{0}(b)+\int_{a}^{b} u^{1}(s) \mathrm{ds}=0 \quad \text { and } \quad u_{T}^{0}(a)+u_{T}^{0}(b)+\int_{a}^{b} u_{T}^{1}(s) \mathrm{ds}=0
$$

Then there exist unique functions $v_{a}, v_{b} \in H^{1}(0, T)$ such that the solution of (1)-(3) satisfies

$$
\begin{equation*}
u(x, T)=u_{T}^{0}(x) \quad \text { and } \quad u_{t}(x, T)=u_{T}^{1}(x), \quad x \in I . \tag{29}
\end{equation*}
$$

Proof. - First we solve the problem

$$
\begin{gathered}
\left(z_{t t}-z_{x x}\right)(x, t)=0, \quad(x, t) \in I \times(0, T), \\
z(a, t)=z(b, t)=0, \quad t \in(0, T), \\
z(x, T)=u_{T}^{0}(x) \quad \text { and } \quad z_{t}(x, T)=u_{T}^{1}(x), \quad x \in I,
\end{gathered}
$$

and then we choose (using theorem 0.1) $v_{a}$ and $v_{b}$ such that the solution of the problem

$$
\begin{gathered}
\left(y_{t t}-y_{x x}\right)(x, t)=0, \quad(x, t) \in I \times(0, T), \\
y(a, t)=v_{a}(t), \quad y(b, t)=v_{b}(t), \quad t \in(0, T), \\
y(x, 0)=u^{0}(x)-z(x, 0) \quad \text { and } \quad y_{t}(x, 0)=u^{1}(x)-z_{t}(x, 0), \quad x \in I
\end{gathered}
$$

satisfies

$$
y(x, T)=y_{t}(x, T)=0, \quad x \in I .
$$

Then $u:=y+z$ is the solution of (1)-(3) and it satisfies (29).
Let us return to theorem 0.1. From the point of view of applications it would be useful to find controls $v_{a}, v_{b}$ defined by some "feedback law" $v_{a}=F_{a}(u), v_{b}=F_{b}(u)$ with explicitly given functions $F_{a}, F_{b}$ : this would realize an "automatic" control of the system. For the problem (1)-(3) such feedbacks can be found easily : let us consider the problem

$$
\begin{gather*}
\left(u_{t t}-u_{x x}\right)(x, t)=0, \quad(x, t) \in I \times \mathbb{R}_{+}  \tag{30}\\
\left(u_{x}-u_{t}\right)(a, t)=\left(u_{x}+u_{t}\right)(b, t)=0, \quad t \in \mathbb{R}_{+}  \tag{31}\\
u(x, 0)=u^{0}(x) \quad \text { and } \quad u_{t}(x, 0)=u^{1}(x), \quad x \in I \tag{32}
\end{gather*}
$$

By definition, a solution of this problem is a function

$$
u \in C\left([0,+\infty] ; H^{1}(I)\right) \cap C^{1}\left([0,+\infty] ; L^{2}(I)\right)
$$

satisfying (30) in the distributional sense, (32) almost everywhere and (31) in the following sense : in some neighbourhood of $\{a\} \times \mathbb{R}_{+} u(x, t)$ is a function of $x+t$ and in some neighbourhood of $\{b\} \times \mathbb{R}_{+}$it is a function of $x-t$.

Theorem 0.5. - Let $\left(u^{0}, u^{1}\right) \in H^{1}(I) \times L^{2}(I)$ be such that

$$
\begin{equation*}
u^{0}(a)+u^{0}(b)+\int_{a}^{b} u^{1}(s) d s=0 \tag{33}
\end{equation*}
$$

Then the solution of (30)-(32) satisfies

$$
\begin{equation*}
u(x, t)=0, \quad \forall(x, t) \in I \times[b-a,+\infty) \tag{34}
\end{equation*}
$$

Proof. - Adapting the proof of theorem 0.1 we obtain that the solution of (30)-(32) has again the form (11) with suitable functions $f:(a,+\infty) \rightarrow \mathbb{R}$ et $g:(-\infty, b) \rightarrow \mathbb{R}$ satisfying (14), (15) and

$$
\begin{align*}
& g^{\prime}(a-t)=0, \quad t \in \mathbb{R}_{+},  \tag{35}\\
& f^{\prime}(b+t)=0, \quad t \in \mathbb{R}_{+} . \tag{36}
\end{align*}
$$

It follows that $f, g$ satisfy (16),(17) with a suitable primitive $U^{1}$ of $u^{1}$, and that

$$
\begin{array}{ll}
g(s)=g(a), & s \leq a \\
f(s)=f(b), & s \geq b \tag{38}
\end{array}
$$

Introducing the subsets $U, R, D, L$ as before, but replacing $T$ by $+\infty$, we deduce from (11), (16), (17), (37) and (38) the following formulae :

$$
2 u(x, t)= \begin{cases}u^{0}(a)+u^{0}(b)+\int_{a}^{b} u^{1}(s) \mathrm{ds}, & \text { if }(x, t) \in U  \tag{39}\\ u^{0}(x-t)+u^{0}(b)+\int_{x-t}^{b} u^{1}(s) \mathrm{ds}, & \text { if }(x, t) \in R \\ u^{0}(x+t)+u^{0}(x-t)+\int_{x-t}^{x+t} u^{1}(s) \mathrm{ds}, & \text { if }(x, t) \in D \\ u^{0}(x+t)+u^{0}(a)+\int_{a}^{x+t} u^{1}(s) \mathrm{ds}, & \text { if }(x, t) \in L\end{cases}
$$

These show in particular the uniqueness of the solution of (30)-(32).
One can verify by a direct computation that the function defined by (39) is indeed a solution of (30)-(32). Finally, (34) follows from (33) and (39).

## 1. Linear evolutionary problems

The results of this chapter are standard; see e.g. Lions and Magenes [1] for proof.

### 1.1. The diagram $\mathrm{V} \subset \mathbf{H}=\mathbf{H}^{\prime} \subset \mathrm{V}^{\prime}$

Let $V$ be an infinite-dimensional, separable, real or complex Hilbert space and introduce the duality mapping $A: V \rightarrow V^{\prime}$ defined by

$$
\begin{equation*}
\langle A u, v\rangle_{V^{\prime}, V}:=(u, v)_{V}, u, v \in V . \tag{1}
\end{equation*}
$$

By the Riesz-Fréchet representation theorem $A$ is an isometric isomorphism of $V$ onto $V^{\prime}$.

Let $H$ be another Hilbert space with a dense and compact imbedding $V \subset H$. (The compactness means that every bounded subset of $V$ is precompact in $H$.) Then in addition $H$ is of infinite dimension, separable, and the imbedding $H^{\prime} \subset V^{\prime}$ is also dense and compact. Identifying $H$ with $H^{\prime}$, we obtain the diagram

$$
\begin{equation*}
V \subset H=H^{\prime} \subset V^{\prime} \tag{2}
\end{equation*}
$$

We deduce from (1) that

$$
\begin{equation*}
(A u, v)_{H}=(u, v)_{V}, \forall u \in V \quad \text { such that } \quad A u \in H, \forall v \in V . \tag{3}
\end{equation*}
$$

Denoting by $i$ the compact imbedding of $V$ into $V^{\prime}$, the linear map $T:=A^{-1} \circ i: V \rightarrow V$ is also compact. Moreover, it is selfadjoint. Indeed, given $u, v \in V$ arbitrarily, we deduce from (3) that

$$
(T u, v)_{V}=\left(A^{-1} u, v\right)_{V}=(u, v)_{H}
$$

and

$$
(u, T v)_{V}=\left(u, A^{-1} v\right)_{V}=(u, v)_{H}
$$

whence

$$
(T u, v)_{V}=(u, T v)_{V}
$$

Applying the spectral theorem to $T=A^{-1} \circ i$ we conclude that there exists a sequence $\lambda_{1}, \lambda_{2}, \ldots$ of distinct real numbers (the eigenvalues of $A$ ) and a sequence $Z_{1}, Z_{2}, \ldots$ of subspaces of $V$ such that

$$
\begin{gather*}
\left|\lambda_{k}\right| \rightarrow+\infty, \\
A z=\lambda_{k} z, \forall z \in Z_{k}, \forall k \geq 1,  \tag{4}\\
\operatorname{dim} Z_{k}<+\infty, \forall k \geq 1,  \tag{5}\\
Z_{k} \perp Z_{l} \quad \text { in } \quad V \quad \text { if } \quad k \neq l \tag{6}
\end{gather*}
$$

and
the vector space $Z$ generated by $\cup Z_{k}$ is dense in $V$.
It follows from (3) and (4) that the eigenvalues are positive. We may thus assume that

$$
\begin{equation*}
0<\lambda_{1}<\lambda_{2}<\cdots \quad \text { and } \quad \lambda_{k} \rightarrow+\infty \tag{8}
\end{equation*}
$$

It follows from (6) and (7) that every $v \in V$ has a unique orthogonal expansion

$$
\begin{equation*}
v=\sum v_{k}, v_{k} \in Z_{k}, \forall k \geq 1 \tag{9}
\end{equation*}
$$

converging in $V$; furthermore, we deduce from (6), (9), (3) and (4) that

$$
\begin{equation*}
\|v\|_{V}^{2}=\sum\left\|v_{k}\right\|_{V}^{2}=\sum \lambda_{k}\left\|v_{k}\right\|_{H}^{2} . \tag{10}
\end{equation*}
$$

Now observe that

$$
\begin{equation*}
Z_{k} \perp Z_{l} \quad \text { in } \quad H \quad \text { and also in } \quad V^{\prime} \quad \text { if } k \neq l . \tag{11}
\end{equation*}
$$

Indeed, the orthogonality in $H$ follows from (3), (4), (6) and (8), while the orthogonality in $V^{\prime}$ follows from (4) and (6), using the isometric property of $A$ : for $u \in Z_{k}$ and $v \in Z_{l}, k \neq l$, we have

$$
(u, v)_{V^{\prime}}=\left(A^{-1} u, A^{-1} v\right)_{V}=\lambda_{k}^{-1} \lambda_{l}^{-1}(u, v)_{V}=0
$$

Next we deduce from (9) and from the density of the imbeddings in (2) that $Z$ is also dense in $H$ and in $V^{\prime}$. Therefore every $v \in H$ (resp. every $v \in V^{\prime}$ )
has a unique orthogonal expansion of the form (9), converging in $H$ (resp. in $V^{\prime}$ ), and we have

$$
\begin{equation*}
\|v\|_{H}^{2}=\sum\left\|v_{k}\right\|_{H}^{2} \tag{12}
\end{equation*}
$$

(resp.

$$
\begin{equation*}
\left.\|v\|_{V^{\prime}}^{2}=\sum\left\|v_{k}\right\|_{V^{\prime}}^{2}=\sum \lambda_{k}^{-1}\left\|v_{k}\right\|_{H}^{2}\right) \tag{13}
\end{equation*}
$$

The last equality of (13) follows from (10) :

$$
\left\|v_{k}\right\|_{V^{\prime}}^{2}=\left\|A^{-1} v_{k}\right\|_{V}^{2}=\lambda_{k}^{-2}\left\|v_{k}\right\|_{V}^{2}=\lambda_{k}^{-1}\left\|v_{k}\right\|_{H}^{2}
$$

Let $\alpha \in \mathbb{R}$ and define a euclidean norm $\|\cdot\|_{\alpha}$ on $Z$ by putting

$$
\|v\|_{\alpha}:=\left(\sum \lambda_{k}^{2 \alpha}\left\|v_{k}\right\|_{H}^{2}\right)^{1 / 2}
$$

(here we use the expansion (9) of $v$ ). Completing $Z$ with respect to this norm we obtain a Hilbert space which will be denoted by $D_{\alpha}$. It is easy to verify the orthogonality relations

$$
Z_{k} \perp Z_{l} \quad \text { dans } \quad D_{\alpha} \quad \text { si } \quad k \neq l .
$$

It is clear that for $\alpha>\beta$ the norm $\|\cdot\|_{\alpha}$ is stronger than $\|\cdot\|_{\beta}$. Thus we may assume that

$$
D_{\alpha} \subset D_{\beta} \quad \text { if } \quad \alpha>\beta
$$

with a dense and continuous imbedding. (One can readily verify that these imbeddings are in fact compact.) Set

$$
D_{-\infty}:=\cup_{\alpha} D_{\alpha}
$$

For each fixed real number $\alpha$ let us introduce a linear mapping $A^{\alpha}: D_{-\infty} \rightarrow D_{-\infty}$ in the following way : first, for $v \in Z$ given by (9) we set

$$
\begin{equation*}
A^{\alpha} v:=\sum \lambda_{k}^{\alpha} v_{k} \tag{14}
\end{equation*}
$$

Then for any given $v \in D_{-\infty}$ we choose $\beta \in \mathbb{R}$ such that $v \in D_{\beta}$ and then we choose a sequence $\left(v_{j}\right)$ in $Z$ such that $\left\|v-v_{j}\right\|_{\beta} \rightarrow 0$. One can readily verify that $A^{\alpha} v_{j}$ is a Cauchy sequence in $D_{\beta-\alpha}$, hence it converges to a certain $w \in D_{\beta-\alpha}$. Furthermore, it is easy to show that the limit $w$ is independent of the particular choice of $\beta$ and of the sequence $v_{j}$. Define $A^{\alpha} v:=w$. It is clear that for $v \in Z$ this definition reduces to (14).

The following properties are easy to verify :
Given $\alpha, \beta \in \mathbb{R}$ arbitrarily, the restriction of $A^{\alpha}$ to $D_{\beta}$ is an isometric isomorphism of $D_{\beta}$ onto $D_{\beta-\alpha}$.

We have $A^{\alpha} A^{\beta}=A^{\alpha+\beta}$ for all $\alpha, \beta \in \mathbb{R}$.
We have

$$
\|v\|_{1 / 2}=\|v\|_{V}, \quad\|v\|_{0}=\|v\|_{H} \quad \text { and } \quad\|v\|_{-1 / 2}=\|v\|_{V^{\prime}}
$$

for every $v \in Z$.
As a consequence, using also the density of $Z$ in $V, H$ and in $V^{\prime}$, we may identify $D_{1 / 2}$ with $V, D_{0}$ with $H$ and $D_{-1 / 2}$ with $V^{\prime}$. Then $A^{1}$ is an extension of $A$ onto $D_{-\infty}$. In the sequel we shall also denote this extension by $A$ i.e. we shall write $A$ instead of $A^{1}$.

### 1.2. The equation $\mathbf{u}^{\prime \prime}+\mathbf{A u}=\mathbf{0}$

Let us first consider the homogeneous evolutionary problem

$$
\begin{equation*}
u^{\prime \prime}+A u=0 \quad \text { dans } \quad \mathbb{R}, \quad u(0)=u^{0} \text { et } u^{\prime}(0)=u^{1} \tag{15}
\end{equation*}
$$

with arbitrarily given initial data $u^{0}, u^{1} \in D_{-\infty}$. We shall use the orthogonal expansion of $u^{0}$ and $u^{1}$ :

$$
\begin{equation*}
u^{0}=\sum_{k=1}^{\infty} u_{k}^{0} \text { and } u^{1}=\sum_{k=1}^{\infty} u_{k}^{1}, \quad u_{k}^{0}, u_{k}^{1} \in Z_{k}, \quad k=1,2, \ldots \tag{16}
\end{equation*}
$$

Theorem 1.1. - Let $\alpha \in \mathbb{R}$ and $\left(u^{0}, u^{1}\right) \in D_{\alpha+1 / 2} \times D_{\alpha}$. Then the problem (15) has a unique solution such that

$$
\begin{equation*}
u \in C\left(\mathbb{R} ; D_{\alpha+1 / 2}\right) \cap C^{1}\left(\mathbb{R} ; D_{\alpha}\right) \cap C^{2}\left(\mathbb{R} ; D_{\alpha-1 / 2}\right) \tag{17}
\end{equation*}
$$

it is given by the series

$$
\begin{equation*}
u(t)=\sum_{k=1}^{\infty} u_{k}^{0} \cos \left(\sqrt{\lambda_{k}} t\right)+u_{k}^{1} \frac{\sin \left(\sqrt{\lambda_{k}} t\right)}{\sqrt{\lambda_{k}}}, t \in \mathbb{R} \tag{18}
\end{equation*}
$$

The energy $E_{\alpha}: \mathbb{R} \rightarrow \mathbb{R}_{+}$of the solution, defined by

$$
\begin{equation*}
E_{\alpha}(t):=\frac{1}{2}\|u(t)\|_{\alpha+1 / 2}^{2}+\frac{1}{2}\left\|u^{\prime}(t)\right\|_{\alpha}^{2} \tag{19}
\end{equation*}
$$

is in fact independent of $t \in \mathbb{R}$.
The linear mapping $\left(u^{0}, u^{1}\right) \mapsto u$ is continuous from $D_{\alpha+1 / 2} \times D_{\alpha}$ into

$$
C_{b}\left(\mathbb{R} ; D_{\alpha+1 / 2}\right) \cap C_{b}^{1}\left(\mathbb{R} ; D_{\alpha}\right) \cap C_{b}^{2}\left(\mathbb{R} ; D_{\alpha-1 / 2}\right)
$$

Since the energy is independent of $t \in \mathbb{R}$, we shall often write $E_{\alpha}(u)$ instead of $E_{\alpha}\left(u^{0}, u^{1}\right)$.

Let us recall that for $\alpha=0$ we have $D_{\alpha+1 / 2}=V, D_{\alpha}=H$ and $D_{\alpha-1 / 2}=V^{\prime}$. In this case we shall write $E$ instead of $E_{0}$.

Remark 1.2. - The formula (18) shows that if $u^{0} \perp Z_{k}$ and $u^{1} \perp Z_{k}$, then $u(t) \perp Z_{k}$ and $u^{\prime}(t) \perp Z_{k}$ for all $t \in \mathbb{R}$. Since the equation (15) is autonomous, we conclude more generally that if $u(T) \perp Z_{k}$ and $u^{\prime}(T) \perp Z_{k}$ for some $T \in \mathbb{R}$, then necessarily $u(t) \perp Z_{k}$ and $u^{\prime}(t) \perp Z_{k}$ for all $t \in \mathbb{R}$.

Remark 1.3.- One can readily verify that if $\left(u^{0}, u^{1}\right) \in D_{\alpha+1 / 2} \times D_{\alpha}$ for some $\alpha \in \mathbb{R}$ and if $u$ is the corresponding solution of (15), then for any fixed $\beta \in \mathbb{R}, A^{\beta} u$ is the solution of (15) with $\left(u^{0}, u^{1}\right)$ replaced by

$$
\left(A^{\beta} u^{0}, A^{\beta} u^{1}\right) \in \times D_{\alpha-\beta+1 / 2} \times D_{\alpha-\beta}
$$

### 1.3. The wave equation

Let $\Omega$ be a bounded domain (that is, a non-empty open connected set) of class $C^{2}$ in $\mathbb{R}^{n}$; we denote by $\nu$ the outward unit normal vector to its boundary $\Gamma$. Let $\left\{\Gamma_{0}, \Gamma_{1}\right\}$ be a partition of $\Gamma$ (the cases $\Gamma_{0}=\emptyset$ or $\Gamma_{1}=\emptyset$ are not excluded) and let $q: \Omega \rightarrow \mathbb{R}, a: \Gamma_{1} \rightarrow \mathbb{R}$ be two given nonnegative functions.

We consider the problem

$$
\begin{gather*}
u^{\prime \prime}-\Delta u+q u=0 \quad \text { in } \Omega \times \mathbb{R},  \tag{20}\\
u=0 \quad \text { on } \quad \Gamma_{0} \times \mathbb{R},  \tag{21}\\
\partial_{\nu} u+a u=0 \quad \text { on } \quad \Gamma_{1} \times \mathbb{R},  \tag{22}\\
u(0)=u^{0} \quad \text { and } \quad u^{\prime}(0)=u^{1} . \tag{23}
\end{gather*}
$$

In order to avoid some difficulties (studied in detail by Grisvard [1]), we shall assume throughout this book that

$$
\begin{equation*}
\overline{\Gamma_{0}} \cap \overline{\Gamma_{1}}=\emptyset \tag{24}
\end{equation*}
$$

and that

$$
\begin{equation*}
q \in L^{\infty}(\Omega), \quad a \in C^{1}\left(\Gamma_{1}\right) \tag{25}
\end{equation*}
$$

We introduce two real or complex Hilbert spaces $H$ and $V$ by the formulae

$$
\begin{align*}
H & :=\left\{v \in L^{2}(\Omega): \int_{\Omega} v \mathrm{dx}=0\right. \\
V & :=\left\{v \in H^{1}(\Omega): \int_{\Omega} v \mathrm{dx}=0\right\} \tag{26}
\end{align*}
$$

if $\Gamma_{0}=\emptyset, q \equiv 0$ and $a \equiv 0$, and

$$
\begin{equation*}
H:=L^{2}(\Omega) \quad \text { and } \quad V:=H^{1}(\Omega) \tag{27}
\end{equation*}
$$

otherwise ; in both cases we define their norms by

$$
\begin{equation*}
\|v\|_{H}:=\left(\int_{\Omega}|v|^{2} \mathrm{dx}\right)^{1 / 2} \tag{28}
\end{equation*}
$$

and

$$
\begin{equation*}
\|v\|_{V}:=\left(\int_{\Omega}|\nabla v|^{2}+q|v|^{2} \mathrm{dx}+\int_{\Gamma_{1}} a|v|^{2} \mathrm{~d} \Gamma\right)^{1 / 2} \tag{29}
\end{equation*}
$$

One can readily verify that the seminorm $\|\cdot\|_{V}$ is indeed a norm and that it is equivalent to the norm induced by $H^{1}(\Omega)$. Applying Rellich's theorem it follows that the imbedding $V \subset H$ is dense and compact.

We introduce the corresponding linear map $A$ and we define the solution of (20)-(23) as the (unique) solution of the problem

$$
\begin{equation*}
u^{\prime \prime}+A u=0 \quad \text { in } \quad \mathbb{R}, \quad u(0)=u^{0}, \quad u^{\prime}(0)=u^{1} \tag{30}
\end{equation*}
$$

in the sense of theorem 1.1. We justify this definition by showing that every sufficiently smooth classical solution of (20)-(23) (for example of class $C^{2}$ on $\bar{\Omega} \times \mathbb{R})$ is also a solution of (30).

Fixing $t \in \mathbb{R}$ and $v \in V$ arbitrarily, we have

$$
\begin{gathered}
0=\int_{\Omega}\left(u^{\prime \prime}-\Delta u+q u\right) \bar{v} \mathrm{dx} \\
=\int_{\Omega} u^{\prime \prime} \bar{v}+\nabla u \nabla \bar{v}+q u \bar{v} \mathrm{dx}-\int_{\Gamma} \partial_{\nu} u \bar{v} \mathrm{~d} \Gamma \\
=\int_{\Omega}\left(u^{\prime \prime}+q u\right) \bar{v}+\nabla u \nabla \bar{v} \mathrm{dx}+\int_{\Gamma_{1}} a u \bar{v} \mathrm{~d} \Gamma \\
=\left\langle u^{\prime \prime}, v\right\rangle_{V^{\prime}, V}+(u, v)_{V}=\left\langle u^{\prime \prime}+A u, v\right\rangle_{V^{\prime}, V}
\end{gathered}
$$

on $\mathbb{R}$; hence $u^{\prime \prime}(t)+A u(t)=0$ (as an element of $V^{\prime}$ ) for every $t \in \mathbb{R}$.
Let us note that $A u=-\Delta u+q u$ for every $u \in V \cap H^{2}(\Omega)$. Indeed, we have for every $v \in V$ the following equality :

$$
\begin{gathered}
\int_{\Omega}(-\Delta u+q u) \bar{v} \mathrm{dx}=\int_{\Omega} \nabla u \nabla \bar{v}+q u \bar{v} \mathrm{dx}-\int_{\Gamma} \partial_{\nu} u \bar{v} \mathrm{~d} \Gamma \\
\quad=\int_{\Omega} \nabla u \nabla \bar{v}+q u \bar{v} \mathrm{dx}+\int_{\Gamma_{1}} a u \bar{v} \mathrm{~d} \Gamma=\langle A u, v\rangle_{V^{\prime}, V}
\end{gathered}
$$

We remark that for this choice of $H, V$ and $A$ the energy of the solutions of (30) is given by the formula

$$
\begin{equation*}
E(t):=\frac{1}{2} \int_{\Omega}\left|u^{\prime}(t)\right|^{2}+|\nabla u(t)|^{2}+q|u(t)|^{2} \mathrm{dx}+\frac{1}{2} \int_{\Gamma_{1}} a|u(t)|^{2} \mathrm{~d} \Gamma \tag{31}
\end{equation*}
$$

Remark 1.4. - We recall from the elliptic regularity theory that under conditions (24) and (25) for any given $g \in L^{2}(\Omega)$ the solution $v \in V$ of the problem

$$
\begin{aligned}
&-\Delta v+q v=g \text { in } \quad \Omega, \\
& v=0 \quad \text { on } \Gamma_{0} \\
& \partial_{\nu} v+a v=0 \text { on } \\
& \Gamma_{1}
\end{aligned}
$$

belongs to $H^{2}(\Omega)$ and that we have the estimate

$$
\|v\|_{H^{2}(\Omega)} \leq c\|g\|_{L^{2}(\Omega)}
$$

with a constant $c$, independent of the choice of $g$. It follows that the eigenfunctions of $A$ belong to $H^{2}(\Omega)$.

If $u^{0}, u^{1} \in Z$, then it follows easily from formula (18) (which now reduces to a finite sum) that the solution of (20)-(23) satisfies

$$
\begin{equation*}
u \in C^{\infty}\left(\mathbb{R} ; H^{2}(\Omega)\right) \tag{32}
\end{equation*}
$$

### 1.4. A Petrovsky system

Let $\Omega$ be a bounded domain in $\mathbb{R}^{n}$ having a boundary $\Gamma$ of class $C^{4}$ and consider the problem

$$
\begin{gather*}
u^{\prime \prime}+\Delta^{2} u=0 \quad \text { in } \quad \Omega \times \mathbb{R}  \tag{33}\\
u=\partial_{\nu} u=0 \quad \text { on } \quad \Gamma \times \mathbb{R}  \tag{34}\\
u(0)=u^{0} \quad \text { and } \quad u^{\prime}(0)=u^{1} \tag{35}
\end{gather*}
$$

We choose $H=L^{2}(\Omega)$ with its usual norm and $V=H_{0}^{2}(\Omega)$ endowed with the norm $\|v\|_{V}:=\|\Delta v\|_{L^{2}(\Omega)}$; using remark 1.4 with $\Gamma_{0}=\Gamma$ and $q=0$ we conclude that this is indeed a norm on $V$ (and even on $H^{2}(\Omega) \cap H_{0}^{1}(\Omega)$ ), and that this norm is equivalent to the norm induced by $H^{2}(\Omega)$.

It follows from Rellich's theorem that the imbedding $V \subset H$ is dense and compact. Introducing the corresponding operator $A$ we define the solution of (33)-(35) as the solution of the problem

$$
\begin{equation*}
u^{\prime \prime}+A u=0 \quad \text { in } \quad \mathbb{R}, \quad u(0)=u^{0}, \quad u^{\prime}(0)=u^{1} \tag{36}
\end{equation*}
$$

in the sense of theorem 1.1. To justify this definition we show that every sufficiently smooth (for example of class $C^{4}$ on $\bar{\Omega} \times \mathbb{R}$ ) classical solution of (33)-(35) is also a solution of (36). The only nontrivial property is the equality

$$
\begin{equation*}
u^{\prime \prime}+A u=0 \quad \text { in } \quad \mathbb{R} . \tag{37}
\end{equation*}
$$

Fixing $t \in \mathbb{R}$ and $v \in C_{c}^{\infty}(\Omega)$ arbitrarily, we deduce from (33)-(35) that

$$
\begin{aligned}
0 & =\int_{\Omega}\left(u^{\prime \prime}+\Delta^{2} u\right) \bar{v} \mathrm{dx}=\int_{\Omega} u^{\prime \prime} \bar{v}+\Delta u \Delta \bar{v} \mathrm{dx} \\
& =\left\langle u^{\prime \prime}, v\right\rangle_{V^{\prime}, V}+(u, v)_{V}=\left\langle u^{\prime \prime}+A u, v\right\rangle_{V^{\prime}, V}
\end{aligned}
$$

Using the density of $C_{c}^{\infty}(\Omega)$ in $V$ hence (37) follows.
One can show similarly (as in the preceding section) that

$$
A v=\Delta^{2} v, \quad \forall v \in H^{4}(\Omega) \cap V
$$

For this choice of $H, V$ and $A$ the energy of the solutions of (36) is given by the formula

$$
\begin{equation*}
E(t):=\frac{1}{2} \int_{\Omega}\left|u^{\prime}(t)\right|^{2}+|\Delta u(t)|^{2} \mathrm{dx} \tag{38}
\end{equation*}
$$

Remark 1.5. - We recall from the elliptic regularity theory that for $g \in L^{2}(\Omega)$ the solution $v \in V$ of the problem

$$
\begin{gathered}
\Delta^{2} v=g \quad \text { dans } \quad \Omega, \\
v=\partial_{\nu} v=0 \quad \text { sur } \quad \Gamma
\end{gathered}
$$

belongs to $H^{4}(\Omega)$; moreover, we have the estimate

$$
\|v\|_{H^{4}(\Omega)} \leq c\|g\|_{L^{2}(\Omega)}
$$

with a constant $c$, independent of the choice of $g$. In particular, the eigenfunctions of $A$ belong to $H^{4}(\Omega)$.

As in remark 1.4, hence we conclude that for $u^{0}, u^{1} \in Z$ the solution of (33)-(35) satisfies

$$
\begin{equation*}
u \in C^{\infty}\left(\mathbb{R} ; H^{4}(\Omega)\right) \tag{39}
\end{equation*}
$$

### 1.5. Another Petrovsky system

Let $\Omega$ be again a bounded domain in $\mathbb{R}^{n}$ having a boundary $\Gamma$ of class $C^{4}$, and consider the problem

$$
\begin{gather*}
u^{\prime \prime}+\Delta^{2} u=0 \quad \text { in } \quad \Omega \times \mathbb{R}  \tag{40}\\
u=\Delta u=0 \quad \text { on } \quad \Gamma \times \mathbb{R}  \tag{41}\\
u(0)=u^{0} \quad \text { and } \quad u^{\prime}(0)=u^{1} \tag{42}
\end{gather*}
$$

Set $H=L^{2}(\Omega), V=H^{2}(\Omega) \cap H_{0}^{1}(\Omega)$ and $\|v\|_{V}:=\|\Delta v\|_{L^{2}(\Omega)}$. It follows from remark 1.4 that the latter norm is equivalent to the norm induced by $H^{2}(\Omega)$ on $V$.

It follows from Rellich's theorem that the imbedding $V \subset H$ is dense and compact. Introducing the corresponding operator $A$ we define the solution of (40)-(42) as the solution of the problem

$$
\begin{equation*}
u^{\prime \prime}+A u=0 \quad \text { in } \quad \mathbb{R}, \quad u(0)=u^{0}, \quad u^{\prime}(0)=u^{1} \tag{43}
\end{equation*}
$$

in the sense of theorem 1.1. We justify this definition by showing that every sufficiently smooth (say of class $C^{4}$ ) classical solution of (40)-(42) is also a solution of (43). The only nontrivial property is the equality

$$
\begin{equation*}
u^{\prime \prime}+A u=0 \quad \text { in } \quad \mathbb{R} \tag{44}
\end{equation*}
$$

Fixing $t \in \mathbb{R}$ and $v \in V$ arbitrarily, we deduce from (40)-(42) (using also the relations $\Delta u=v=0$ on $\Gamma$ ) that

$$
\begin{aligned}
0 & =\int_{\Omega}\left(u^{\prime \prime}+\Delta^{2} u\right) \bar{v} \mathrm{dx}=\int_{\Omega} u^{\prime \prime} \bar{v}+\Delta u \Delta \bar{v} \mathrm{dx} \\
& =\left\langle u^{\prime \prime}, v\right\rangle_{V^{\prime}, V}+(u, v)_{V}=\left\langle u^{\prime \prime}+A u, v\right\rangle_{V^{\prime}, V}
\end{aligned}
$$

proving (44).
One can show similarly that

$$
A v=\Delta^{2} v, \quad \forall v \in H^{4}(\Omega) \cap V
$$

In this case the energy of the solutions of (44) is given by the formula

$$
\begin{equation*}
E(t):=\frac{1}{2} \int_{\Omega}\left|u^{\prime}(t)\right|^{2}+|\Delta u(t)|^{2} \mathrm{dx} \tag{45}
\end{equation*}
$$

Remark 1.6. - It follows from the elliptic regularity theory that for any given $g \in L^{2}(\Omega)$ the solution $v \in V$ of the problem

$$
\begin{gathered}
\Delta^{2} v=g \quad \text { in } \quad \Omega, \\
v=\Delta v=0 \quad \text { on } \quad \Gamma
\end{gathered}
$$

belongs to $H^{4}(\Omega)$ and that the estimate

$$
\|v\|_{H^{4}(\Omega)} \leq c\|g\|_{L^{2}(\Omega)}
$$

holds true with a constant $c$, independent of $g$. In particular, the eigenfunctions of $A$ belong to $H^{4}(\Omega)$.

Using formula (18) hence we conclude that if $u^{0}, u^{1} \in Z$, then the solution of (40)-(42) satisfies

$$
\begin{equation*}
u \in C^{\infty}\left(\mathbb{R} ; H^{4}(\Omega)\right) \tag{46}
\end{equation*}
$$

We end this section with a technical result; it shows in particular that the operator $A^{1 / 2}$ coincides with the operator $A$ corresponding to the choice $H=L^{2}(\Omega)$ and $V=H_{0}^{1}(\Omega)$ (the latter being the particular case $\Gamma_{0}=\Gamma$ of the situation studied in section 1.3).

Lemma 1.7. - We have

$$
\begin{gathered}
D_{1 / 4}=H_{0}^{1}(\Omega), \quad\|v\|_{1 / 4}=\|\nabla v\|_{L^{2}(\Omega)} \\
D_{3 / 4}=\left\{v \in H^{3}(\Omega): v=\Delta v=0 \quad \text { on } \quad \Gamma\right\}, \quad\|v\|_{3 / 4}=\|\nabla \Delta v\|_{L^{2}(\Omega)}
\end{gathered}
$$

and $D_{-1 / 4}=H^{-1}(\Omega)$ with the corresponding dual norm. Moreover, we have

$$
\begin{equation*}
A^{1 / 2} v=-\Delta v, \quad \forall v \in Z \tag{47}
\end{equation*}
$$

Proof. - Let us denote by $Z_{1}^{\prime}, Z_{2}^{\prime}, \ldots$ the sequence of the eigenspaces of $-\Delta$ in $H_{0}^{1}(\Omega)$, by $(0<) \lambda_{1}^{\prime}<\lambda_{2}^{\prime}<\cdots$ the sequence of the corresponding eigenvalues and by $Z^{\prime}$ the linear hull of $\cup Z_{k}^{\prime}$. (This is the particular case $\Gamma_{0}=\Gamma$ of the situation studied in §1.3.) If $v \in Z_{k}^{\prime}$ for some $k$, then

$$
\begin{aligned}
& (v, w)_{V}=(\Delta v, \Delta w)_{H}=-\int_{\Omega} \lambda_{k}^{\prime} v(\Delta w) \mathrm{dx} \\
= & \lambda_{k}^{\prime} \int_{\Omega}(-\Delta v) w \mathrm{dx}=\left(\lambda_{k}^{\prime}\right)^{2}(v, w)_{H}, \quad \forall w \in V
\end{aligned}
$$

whence

$$
\begin{equation*}
A v=\left(\lambda_{k}^{\prime}\right)^{2} v, \quad \forall v \in Z_{k}^{\prime}, \quad k=1,2, \ldots \tag{48}
\end{equation*}
$$

Thus the eigenfunctions of $-\left.\Delta\right|_{H_{0}^{1}(\Omega)}$ are also eigenfunctions of $A$ and $\lambda_{k}=\left(\lambda_{k}^{\prime}\right)^{2}$. Since $Z^{\prime}$ is dense in $L^{2}(\Omega), A$ cannot have other eigenfunctions. Therefore $Z_{k}^{\prime}=Z_{k}, \forall k$ and $Z^{\prime}=Z$. Then (47) follows from (48).

For any $v=\sum v_{k} \in Z$ we have

$$
\|v\|_{1 / 4}^{2}=\left\|\sum v_{k}\right\|_{1 / 4}^{2}=\sum \lambda_{k}^{1 / 2}\left\|v_{k}\right\|_{H}^{2}=\sum \lambda_{k}^{\prime}\left\|v_{k}\right\|_{H}^{2}=\|v\|_{H_{0}^{1}(\Omega)}^{2}
$$

hence $D_{1 / 4}=H_{0}^{1}(\Omega)$ because $Z=Z^{\prime}$ is dense in $H_{0}^{1}(\Omega)$.
Since $A^{1 / 2}$ is an isometric isomorphism of $D_{1 / 4}$ onto $D_{-1 / 4}$ and $-\Delta$ is an isometric isomorphism of $H_{0}^{1}(\Omega)$ onto $H^{-1}(\Omega)$, using (47) hence we deduce that $H^{-1}(\Omega)=D_{-1 / 4}$.

Using (47) we obtain also that $D_{3 / 4}$ is the set of solutions of the problem

$$
-\Delta v=g \quad \text { in } \quad \Omega, \quad v=0 \quad \text { on } \quad \Gamma
$$

where $g$ runs over $D_{1 / 4}=H_{0}^{1}(\Omega)$. It follows that

$$
\left\{v \in H^{3}(\Omega): v=\Delta v=0 \quad \text { sur } \quad \Gamma\right\} \subset D_{3 / 4}
$$

the inverse inclusion follows from remark 1.4.
The norm equalities in the formulation of the lemma follow from the above computations.

## 2. Hidden regularity. Weak solutions

The results of $\S \S 2.2-2.4$ are due essentially to Lions [2] ; see also Lasiecka and Triggiani [1]. We consider the real case only : the complex case then follows easily by considering separately the real and imaginary parts of the solutions of the corresponding equations.

### 2.1. A special vector field

Let $\Omega \subset \mathbb{R}^{n}$ be a bounded open domain of class $C^{k}, k \geq 1$ and let $\left\{\Gamma_{0}, \Gamma_{1}\right\}$ be a partition of its boundary $\Gamma$ such that

$$
\overline{\Gamma_{0}} \cap \overline{\Gamma_{1}}=\emptyset .
$$

Thus, we have in particular $\overline{\Gamma_{0}}=\Gamma_{0}$ and $\overline{\Gamma_{1}}=\Gamma_{1}$.
We recall for the reader's convenience the following standard construction :
Lemma 2.1. - There exists a vector field $h: \bar{\Omega} \rightarrow \mathbb{R}^{n}$ of class $C^{k-1}$ such that

$$
h=\nu \quad \text { on } \quad \Gamma_{0} \quad \text { and } \quad h=0 \quad \text { on } \quad \Gamma_{1} .
$$

Proof. - Since $\Omega$ is of class $C^{k}$, for every fixed $x^{0} \in \Gamma_{0}$ there is an open neighbourhood $V$ of $x^{0}$ in $\mathbb{R}^{n}$ and a function $\varphi: V \rightarrow \mathbb{R}$ of class $C^{k}$ such that

$$
\nabla \varphi(x) \neq 0, \quad \forall x \in V
$$

and

$$
\varphi(x)=0 \Leftrightarrow x \in V \cap \Gamma .
$$

Replacing $\varphi$ by $-\varphi$ if needed, we may assume that

$$
\nu\left(x^{0}\right) \cdot \nabla \varphi\left(x^{0}\right)>0 .
$$

Choosing $V$ sufficiently small we may assume also that

$$
\bar{V} \cap \Gamma_{1}=\emptyset
$$

and that
$V \cap \Gamma \quad$ is connected.

Then the function $\psi: V \rightarrow \mathbb{R}^{n}$ defined by $\psi:=\nabla \varphi /|\nabla \varphi|$ is of class $C^{k-1}$ and $\psi=\nu$ on $V \cap \Gamma$.

Since $\Omega$ is bounded, $\Gamma_{0}$ is compact ; therefore it can be covered with a finite number of neighbourhoods $V_{1}, \ldots, V_{m}$ of this type. Denoting by $\psi_{1}, \ldots, \psi_{m}$ the corresponding functions we have

$$
\begin{gathered}
\Gamma_{0} \subset V_{1} \cup \ldots \cup V_{m}, \\
\left(V_{1} \cup \ldots \cup V_{m}\right) \cap \Gamma_{1}=\emptyset
\end{gathered}
$$

and

$$
\psi_{i}=\nu \quad \text { on } \quad V_{i} \cap \Gamma_{0}, \quad i=1, \ldots, m
$$

Then we fix an open set $V_{0}$ in $\mathbb{R}^{n}$ such that

$$
\begin{gathered}
\bar{\Omega} \subset V_{0} \cap \ldots \cap V_{m}, \\
V_{0} \cap \Gamma_{0}=\emptyset
\end{gathered}
$$

and we define $\psi_{0}: V_{0} \rightarrow \mathbb{R}^{n}$ by $\psi_{0}(x)=0, \forall x \in V_{0}$.
Let $\theta_{0}, \ldots, \theta_{m}$ be a partition of unity of class $C^{k}$, corresponding to the covering $V_{0}, \ldots, V_{m}$ of $\bar{\Omega}$ :

$$
\theta_{i} \in C_{c}^{k}\left(V_{i}\right) \quad \text { and } \quad 0 \leq \theta_{i} \leq 1, \quad i=0, \ldots, m
$$

and

$$
\theta_{0}+\cdots+\theta_{m}=1 \quad \text { on } \quad \bar{\Omega}
$$

One can readily verify that the vector field $h$ defined by

$$
h:=\left.\left(\sum_{i=0}^{m} \theta_{i} \psi_{i}\right)\right|_{\bar{\Omega}}
$$

has the desired properties.

### 2.2. The wave equation. Multiplier method

We consider the problem

$$
\begin{gather*}
u^{\prime \prime}-\Delta u+q u=0 \quad \text { in } \quad \Omega \times \mathbb{R}  \tag{1}\\
u=0 \quad \text { on } \quad \Gamma_{0} \times \mathbb{R}  \tag{2}\\
\partial_{\nu} u+a u=0 \quad \text { on } \quad \Gamma_{1} \times \mathbb{R}  \tag{3}\\
u(0)=u^{0} \quad \text { and } \quad u^{\prime}(0)=u^{1} \tag{4}
\end{gather*}
$$

introduced in $\S 1.3$. We recall (see remark 1.4) that for arbitrarily given $u^{0}, u^{1} \in Z$ the solution of (1)-(4) satisfies

$$
\begin{equation*}
u \in C^{\infty}\left(\mathbb{R} ; H^{2}(\Omega)\right) \tag{5}
\end{equation*}
$$

In particular the normal derivative of the solution is well-defined.
The following result will permit us to define the normal derivative of less regular solutions, too.

Theorem 2.2. - Let $T>0$. There is a constant $c=c(T)>0$ such that for every $u^{0}, u^{1} \in Z$ the solution of (1)-(4) satisfies the inequality

$$
\begin{equation*}
\int_{-T}^{T} \int_{\Gamma_{0}}\left|\partial_{\nu} u\right|^{2} \mathrm{~d} \Gamma \mathrm{dt} \leq c\left(\left\|u^{0}\right\|_{V}^{2}+\left\|u^{1}\right\|_{H}^{2}\right) \tag{6}
\end{equation*}
$$

Consequently, there is a unique continuous linear map

$$
L: V \times H \rightarrow L_{l o c}^{2}\left(\mathbb{R} ; L^{2}\left(\Gamma_{0}\right)\right)
$$

such that

$$
L\left(u^{0}, u^{1}\right)=\partial_{\nu} u, \forall\left(u^{0}, u^{1}\right) \in Z \times Z
$$

For the proof we need the following identity :
Lemma 2.3. - Let $u \in H_{l o c}^{2}\left(\mathbb{R} ; H^{2}(\Omega)\right)$ be a function satisfying (1) and let $h: \bar{\Omega} \rightarrow \mathbb{R}^{n}$ be a vector field of class $C^{1}$. Then for any fixed $-\infty<S<T<\infty$ the following identity holds true :

$$
\begin{gather*}
\int_{S}^{T} \int_{\Gamma} 2\left(\partial_{\nu} u\right) h \cdot \nabla u+(h \cdot \nu)\left(u^{\prime}\right)^{2}-(h \cdot \nu)|\nabla u|^{2} \mathrm{~d} \Gamma \mathrm{dt} \\
=\left[\int_{\Omega} 2 u^{\prime} h \cdot \nabla u \mathrm{dx}\right]_{S}^{T}  \tag{7}\\
+\int_{S}^{T} \int_{\Omega}(\operatorname{div} h)\left(\left(u^{\prime}\right)^{2}-|\nabla u|^{2}\right)+2 q u h \cdot \nabla u+\sum_{i, j=1}^{n} 2\left(\partial_{i} h_{j}\right)\left(\partial_{i} u\right)\left(\partial_{j} u\right) \mathrm{dx} \mathrm{dt}
\end{gather*}
$$

(The dot denotes the usual scalar product in $\mathbb{R}^{n}$.)
Proof. - We multiply (1) by $2 h \cdot \nabla u=2 \sum_{j=1}^{n} h_{j} \partial_{j} u$ and we integrate by parts as follows :

$$
\begin{aligned}
& \int_{S}^{T} \int_{\Omega}-2 q u h \cdot \nabla u \mathrm{dx} \mathrm{dt}=\int_{S}^{T} \int_{\Omega} 2\left(u^{\prime \prime}-\Delta u\right) h \cdot \nabla u \mathrm{dx} \mathrm{dt} \\
= & {\left[\int_{\Omega} 2 u^{\prime} h \cdot \nabla u \mathrm{dx}\right]_{S}^{T}+\int_{S}^{T} \int_{\Omega}-2 u^{\prime} h \cdot \nabla u^{\prime}-2(\Delta u) h \cdot \nabla u \mathrm{dx} \mathrm{dt} }
\end{aligned}
$$

$$
\begin{gathered}
=\left[\int_{\Omega} 2 u^{\prime} h \cdot \nabla u \mathrm{dx}\right]_{S}^{T}-\int_{S}^{T} \int_{\Gamma} 2 \partial_{\nu} u h \cdot \nabla u \mathrm{~d} \Gamma \mathrm{dt} \\
\quad+\int_{S}^{T} \int_{\Omega}-h \cdot \nabla\left(u^{\prime}\right)^{2}+2 \nabla u \cdot \nabla(h \cdot \nabla u) \mathrm{dx} \mathrm{dt} \\
=\left[\int_{\Omega} 2 u^{\prime} h \cdot \nabla u \mathrm{dx}\right]_{S}^{T}-\int_{S}^{T} \int_{\Gamma} 2 \partial_{\nu} u h \cdot \nabla u+(h \cdot \nu)\left(u^{\prime}\right)^{2} \mathrm{~d} \Gamma \mathrm{dt} \\
\\
\quad+\int_{S}^{T} \int_{\Omega}(\operatorname{div} h)\left(u^{\prime}\right)^{2}+2 \sum_{i, j=1}^{n} \partial_{i} u \partial_{i}\left(h_{j} \partial_{j} u\right) \mathrm{dx} \mathrm{dt} \\
=\left[\int_{\Omega} 2 u^{\prime} h \cdot \nabla u \mathrm{dx}\right]_{S}^{T}-\int_{S}^{T} \int_{\Gamma} 2 \partial_{\nu} u h \cdot \nabla u+(h \cdot \nu)\left(u^{\prime}\right)^{2} \mathrm{~d} \Gamma \mathrm{dt} \\
+\int_{S}^{T} \int_{\Omega}(\operatorname{div} h)\left(u^{\prime}\right)^{2}+\sum_{i, j=1}^{n} 2\left(\partial_{i} h_{j}\right)\left(\partial_{i} u\right)\left(\partial_{j} u\right)+h_{j} \partial_{j}\left(\left(\partial_{i} u\right)^{2}\right) \mathrm{dx} \mathrm{dt} \\
=\left[\int_{\Omega} 2 u^{\prime} h \cdot \nabla u \mathrm{dx}\right]_{S}^{T}-\int_{S}^{T} \int_{\Gamma} 2 \partial_{\nu} u h \cdot \nabla u+(h \cdot \nu)\left(\left(u^{\prime}\right)^{2}-|\nabla u|^{2}\right) \mathrm{d} \Gamma \mathrm{dt} \\
+ \\
+\int_{S}^{T} \int_{\Omega}(\operatorname{div} h)\left(\left(u^{\prime}\right)^{2}-|\nabla u|^{2}\right)+\sum_{i, j=1}^{n} 2\left(\partial_{i} h_{j}\right)\left(\partial_{i} u\right)\left(\partial_{j} u\right) \mathrm{dx} \mathrm{dt} . \quad \square
\end{gathered}
$$

Proof of theorem 2.2. - Applying theorem 1.1 we have

$$
\begin{equation*}
\|u(t)\|_{V}^{2}+\left\|u^{\prime}(t)\right\|_{H}^{2}=\left\|u^{0}\right\|_{V}^{2}+\left\|u^{1}\right\|_{H}^{2} \tag{8}
\end{equation*}
$$

for every $t \in \mathbb{R}$. (The constant $c_{1}$ does not depend on the choice of $u^{0}, u^{1}$ and f.)

Applying the identity (7) with $S=-T$ and with the vector field $h$ constructed in lemma 2.1 the left-hand side of (7) becomes

$$
\int_{-T}^{T} \int_{\Gamma_{0}}\left|\partial_{\nu} u\right|^{2} \mathrm{~d} \Gamma \mathrm{dt}
$$

(Note that $u^{\prime}=0$ and $\nabla u=\left(\partial_{\nu} u\right) \nu$ on $\Gamma_{0} \times \mathbb{R}$ because of (2).)
Since $h$ is of class $C^{1}$, there exists a constant $c_{1}$ such that

$$
|h(x)| \leq c_{1} \quad \text { and } \quad \sum_{i, j=1}^{n}\left|\partial_{i} h_{j}(x)\right| \leq c_{1}, \forall x \in \bar{\Omega}
$$

Using (8) and these inequalities to majorize the right-hand side of (7) we obtain easily the estimate (6) with a suitable constant $c$.

The last part of the theorem follows from the inequality (6), using the density of $Z \times Z$ in $V \times H$.

Remark 2.4. - Teorem 2.2 justifies the notation $\partial_{\nu} u$ or $\partial u / \partial \nu$ instead of $L\left(u^{0}, u^{1}\right)$ for any $\left(u^{0}, u^{1}\right) \in V \times H$. Then we have the following trace theorem :

$$
\begin{equation*}
\left(u^{0}, u^{1}\right) \in V \times H \Rightarrow \partial_{\nu} u \in L_{l o c}^{2}\left(\mathbb{R} ; L^{2}\left(\Gamma_{0}\right)\right) . \tag{9}
\end{equation*}
$$

This result does not follow from the usual trace theorems of the Sobolev spaces : it is a "hidden regularity" result. The corresponding inequality (6) is often called a direct inequality.

Now we consider the following non-homogeneous problem :

$$
\begin{gather*}
y^{\prime \prime}-\Delta y+q y=0 \quad \text { in } \quad \Omega \times \mathbb{R}  \tag{10}\\
y=v \quad \text { on } \quad \Gamma_{0} \times \mathbb{R}  \tag{11}\\
\partial_{\nu} y+a y=0 \quad \text { on } \quad \Gamma_{1} \times \mathbb{R}  \tag{12}\\
y(0)=y^{0}, \quad y^{\prime}(0)=y^{1} \quad \text { in } \quad \Omega . \tag{13}
\end{gather*}
$$

In order to find a reasonably general definition of the solution we begin with a formal computation. Let us first assume that $y$ is a sufficiently smooth function (say of class $C^{2}$ in $\bar{\Omega} \times \mathbb{R}$ ) satisfying (10)-(13). Fix $\left(u^{0}, u^{1}\right) \in V \times H$ arbitrarily and multiply the corresponding solution of (1)-(4) by $y$. Integrating by parts we obtain for every fixed $S \in \mathbb{R}$ the equality

$$
\begin{gathered}
0=\int_{0}^{S} \int_{\Omega}\left(u^{\prime \prime}-\Delta u+q u\right) y \mathrm{dx} \mathrm{dt}=\left[\int_{\Omega} u^{\prime} y-u y^{\prime} \mathrm{dx}\right]_{0}^{S} \\
\quad-\int_{0}^{S} \int_{\Gamma}\left(\partial_{\nu} u\right) y-u\left(\partial_{\nu} y\right) \mathrm{d} \Gamma \mathrm{dt} \\
\\
+\int_{0}^{S} \int_{\Omega} u\left(y^{\prime \prime}-\Delta y+q y\right) \mathrm{dx} \mathrm{dt} \\
=\int_{\Omega} u^{\prime}(S) y(S)-u(S) y^{\prime}(S)+u^{0} y^{1}-u^{1} y^{0} \mathrm{dx}-\int_{0}^{S} \int_{\Gamma_{0}}\left(\partial_{\nu} u\right) v \mathrm{~d} \Gamma \mathrm{dt}
\end{gathered}
$$

Putting

$$
L_{S}\left(u^{0}, u^{1}\right):=\int_{0}^{S} \int_{\Gamma_{0}}\left(\partial_{\nu} u\right) v \mathrm{~d} \Gamma \mathrm{dt}+\left\langle\left(-y^{1}, y^{0}\right),\left(u^{0}, u^{1}\right)\right\rangle_{V^{\prime} \times H, V \times H}
$$

we may rewrite this identity as

$$
\begin{equation*}
L_{S}\left(u^{0}, u^{1}\right)=\left\langle\left(-y^{\prime}(S), y(S)\right),\left(u(S), u^{\prime}(S)\right)\right\rangle_{V^{\prime} \times H, V \times H} . \tag{14}
\end{equation*}
$$

This leads to the following definition:
Definition. - We say that $\left(y, y^{\prime}\right)$ is a solution of (10)-(13) if $\left(y, y^{\prime}\right) \in C\left(\mathbb{R} ; H \times V^{\prime}\right)$ and if (14) is satisfied for every $S \in \mathbb{R}$ and for every $\left(u^{0}, u^{1}\right) \in V \times H$.

This definition is justified by
Theorem 2.5. - Given $\left(y^{0}, y^{1}\right) \in H \times V^{\prime}$ and $v \in L_{l o c}^{2}\left(\mathbb{R} ; L^{2}\left(\Gamma_{0}\right)\right)$ arbitrarily, the problem (10)-(13) has a unique solution.

Furthermore, the linear map $\left(y^{0}, y^{1}, v\right) \mapsto\left(y, y^{\prime}\right)$ is continuous with respect to these topologies.

Proof. - It follows from theorem 2.2 that for every fixed $S \in \mathbb{R}$ the linear form $L_{S}$ is bounded in $V \times H$. Furthermore, it follows from theorem 1.1 that the linear map

$$
\left(u(S), u^{\prime}(S)\right) \mapsto\left(u^{0}, u^{1}\right)
$$

is an isomorphism of $V \times H$ onto itself. Hence the linear form

$$
\left(u(S), u^{\prime}(S)\right) \mapsto L_{S}\left(u^{0}, u^{1}\right)
$$

is also bounded on $V \times H$. Therefore there is a unique pair $\left(y(S), y^{\prime}(S)\right) \in$ $H \times V^{\prime}$ satisfying (14).

Next we show that the function

$$
\mathbb{R} \ni S \mapsto\left\|\left(y(S), y^{\prime}(S)\right)\right\|_{H \times V^{\prime}}
$$

is bounded in every bounded interval. More precisely, for every bounded interval $I$ and for all $S \in I$ we have the estimate

$$
\begin{equation*}
\left\|\left(y(S), y^{\prime}(S)\right)\right\|_{H \times V^{\prime}} \leq c(I)\left(\|v\|_{L^{2}\left(I ; L^{2}\left(\Gamma_{0}\right)\right)}+\left\|\left(y^{0}, y^{1}\right)\right\|_{H \times V^{\prime}}\right) \tag{15}
\end{equation*}
$$

where $c(I)$ is a constant independent of $S, v, y^{0}$ and $y^{1}$. Indeed, choose $\left(u^{0}, u^{1}\right) \in V \times H$ arbitrarily and write for brevity

$$
Y=\left(-y^{\prime}, y\right), \quad Y^{0}=\left(-y^{1}, y^{0}\right), \quad U=\left(u, u^{\prime}\right), \quad U^{0}=\left(u^{0}, u^{1}\right)
$$

Then, using theorems 1.1 and 2.2 we have the following estimate (we denote by $\|\cdot\|_{I}$ the norm of $L^{2}\left(I ; L^{2}\left(\Gamma_{0}\right)\right)$ for brevity) :

$$
\begin{aligned}
&\left|\langle Y(S), U(S)\rangle_{V^{\prime} \times H, V \times H}\right|=\left|\int_{0}^{S} \int_{\Gamma_{0}}\left(\partial_{\nu} u\right) v \mathrm{~d} \Gamma \mathrm{dt}+\left\langle Y^{0}, U^{0}\right\rangle_{V^{\prime} \times H, V \times H}\right| \\
& \leq\left\|\partial_{\nu} u\right\|_{I}\|v\|_{I}+\left\|Y^{0}\right\|_{V^{\prime} \times H}\left\|U^{0}\right\|_{V \times H} \leq c(I)\left(\|v\|_{I}+\left\|Y^{0}\right\|_{V^{\prime} \times H}\right)\left\|U^{0}\right\|_{V \times H} .
\end{aligned}
$$

Hence (15) follows.
Next we recall (cf. e.g. Lions and Magenes [1]) that if $\left(y^{0}, y^{1}\right) \in V \times H$ and $v \in C_{c}^{\infty}\left(\mathbb{R} ; H^{3 / 2}\left(\Gamma_{0}\right)\right)$ are such that $v(0)=0$, then (4)-(6) has a regular solution

$$
y \in C(\mathbf{R} ; \mathbf{V}) \cap \mathbf{C}^{\mathbf{1}}(\mathbf{R} ; \mathbf{H})
$$

we have in particular $\left(y, y^{\prime}\right) \in C\left(\mathbb{R} ; H \times V^{\prime}\right)$. Since the set of these data $\left(y^{0}, y^{1}, v\right)$ is dense in $H \times V^{\prime} \times L_{l o c}^{2}\left(\mathbb{R} ; L^{2}\left(\Gamma_{0}\right)\right)$, using (15) hence we conclude that $\left(y, y^{\prime}\right) \in C\left(\mathbb{R} ; H \times V^{\prime}\right)$ in the general case, too.

Finally, the continuous dependence of the solution on $\left(y^{0}, y^{1}, v\right)$ also follows from (15).

### 2.3. The first Petrovsky system

Consider the problem introduced in § 1.4 :

$$
\begin{gather*}
u^{\prime \prime}+\Delta^{2} u=0 \quad \text { in } \quad \Omega \times \mathbb{R}  \tag{16}\\
u=\partial_{\nu} u=0 \quad \text { on } \quad \Gamma \times \mathbb{R}  \tag{17}\\
u(0)=u^{0} \quad \text { and } \quad u^{\prime}(0)=u^{1} . \tag{18}
\end{gather*}
$$

We recall (see remark 1.5) that for every $u^{0}, u^{1} \in Z$ the solution of (16)-(18) satisfies

$$
u \in C^{\infty}\left(\mathbb{R} ; H^{4}(\Omega)\right)
$$

In particular, $\Delta u$ is well-defined on $\Gamma \times \mathbb{R}$.
The following result shows that this trace may be defined for weaker solutions as well.

Theorem 2.6. - Given $T>0$ arbitrarily, there exists a constant $c=c(T)>0$ such that for every $\left(u^{0}, u^{1}\right) \in Z \times Z$ the solution of (16)(18) satisfies the inequality

$$
\begin{equation*}
\int_{-T}^{T} \int_{\Gamma}|\Delta u|^{2} \mathrm{~d} \Gamma \mathrm{dt} \leq c\left(\left\|u^{0}\right\|_{H_{0}^{2}(\Omega)}^{2}+\left\|u^{1}\right\|_{L^{2}(\Omega)}^{2}\right) \tag{19}
\end{equation*}
$$

Consequently, there is a unique continuous linear map

$$
L: H_{0}^{2}(\Omega) \times L^{2}(\Omega) \rightarrow L_{l o c}^{2}\left(\mathbb{R} ; L^{2}(\Gamma)\right)
$$

such that

$$
L\left(u^{0}, u^{1}\right)=\Delta u, \forall\left(u^{0}, u^{1}\right) \in Z \times Z .
$$

We begin by establishing an identity analogous to that of lemma 2.3.
Lemma 2.7. - Let $u \in H_{l o c}^{2}\left(\mathbb{R} ; H^{4}(\Omega)\right)$ be a function satisfying (16), (17), and let $h: \bar{\Omega} \rightarrow \mathbb{R}^{n}$ be a vector field of class $C^{2}$. Then for any fixed $-\infty<S<T<\infty$ the following identity holds true :

$$
\begin{gather*}
\int_{S}^{T} \int_{\Gamma}(h \cdot \nu)|\Delta u|^{2} \mathrm{~d} \Gamma \mathrm{dt}=\left[\int_{\Omega} 2 u^{\prime} h \cdot \nabla u \mathrm{dx}\right]_{S}^{T} \\
\quad+\int_{S}^{T} \int_{\Omega}(\operatorname{div} h)\left(\left(u^{\prime}\right)^{2}-|\Delta u|^{2}\right)  \tag{20}\\
+4 \sum_{i, j=1}^{n}\left(\partial_{i} h_{j}\right)(\Delta u)\left(\partial_{i} \partial_{j} u\right)+2 \sum_{i=1}^{n}\left(\Delta h_{i}\right)(\Delta u) \partial_{i} u \mathrm{dx} \mathrm{dt} .
\end{gather*}
$$

Proof. - We multiply (16) by $2 h \cdot \nabla u=2 \sum_{j=1}^{n} h_{j} \partial_{j} u$ and we integrate by parts :

$$
\begin{gathered}
0=\int_{S}^{T} \int_{\Omega} 2\left(u^{\prime \prime}+\Delta^{2} u\right) h \cdot \nabla u \mathrm{dx} \mathrm{dt} \\
=\left[\int_{\Omega} 2 u^{\prime} h \cdot \nabla u \mathrm{dx}\right]_{S}^{T}+\int_{S}^{T} \int_{\Omega}-2 u^{\prime} h \cdot \nabla u^{\prime}+2\left(\Delta^{2} u\right) h \cdot \nabla u \mathrm{dx} \mathrm{dt} \\
=\left[\int_{\Omega} 2 u^{\prime} h \cdot \nabla u \mathrm{dx}\right]_{S}^{T}+2 \int_{S}^{T} \int_{\Gamma}\left(\partial_{\nu} \Delta u\right) h \cdot \nabla u-(\Delta u) \partial_{\nu}(h \cdot \nabla u) \mathrm{d} \Gamma \mathrm{dt} \\
+\int_{S}^{T} \int_{\Omega}-h \cdot \nabla\left(u^{\prime}\right)^{2}+2(\Delta u) \Delta(h \cdot \nabla u) \mathrm{dx} \mathrm{dt} .
\end{gathered}
$$

Remark that

$$
\begin{gathered}
\Delta(h \cdot \nabla u)=\sum_{i, j=1}^{n} \partial_{i}^{2}\left(h_{j} \partial_{j} u\right) \\
=2 \sum_{i, j=1}^{n}\left(\partial_{i} h_{j}\right)\left(\partial_{i} \partial_{j} u\right)+\sum_{j=1}^{n}\left(\Delta h_{j}\right)\left(\partial_{j} u\right)+h_{j} \partial_{j}(\Delta u)
\end{gathered}
$$

whence

$$
\begin{gathered}
2(\Delta u) \Delta(h \cdot \nabla u)=h \cdot \nabla(\Delta u)^{2} \\
+ \\
2 \sum_{j=1}^{n}\left(\Delta h_{j}\right)\left(\partial_{j} u\right)(\Delta u)+4 \sum_{i, j=1}^{n}\left(\partial_{i} h_{j}\right)\left(\partial_{i} \partial_{j} u\right)(\Delta u)
\end{gathered}
$$

and the above identity becomes

$$
\begin{gathered}
0=\left[\int_{\Omega} 2 u^{\prime} h \cdot \nabla u \mathrm{dx}\right]_{S}^{T}+2 \int_{S}^{T} \int_{\Gamma}\left(\partial_{\nu} \Delta u\right) h \cdot \nabla u-(\Delta u) \partial_{\nu}(h \cdot \nabla u) \mathrm{d} \Gamma \mathrm{dt} \\
+\int_{S}^{T} \int_{\Omega} h \cdot \nabla\left((\Delta u)^{2}-\left(u^{\prime}\right)^{2}\right)+2 \sum_{j=1}^{n}\left(\Delta h_{j}\right)\left(\partial_{j} u\right)(\Delta u) \\
+4 \sum_{i, j=1}^{n}\left(\partial_{i} h_{j}\right)\left(\partial_{i} \partial_{j} u\right)(\Delta u) \mathrm{dx} \mathrm{dt} \\
=\left[\int_{\Omega} 2 u^{\prime} h \cdot \nabla u \mathrm{dx}\right]_{S}^{T} \\
+\int_{S}^{T} \int_{\Gamma} 2\left(\partial_{\nu} \Delta u\right) h \cdot \nabla u-2(\Delta u) \partial_{\nu}(h \cdot \nabla u)+(h \cdot \nu)\left((\Delta u)^{2}-\left(u^{\prime}\right)^{2}\right) \mathrm{d} \Gamma \mathrm{dt} \\
\quad+\int_{S}^{T} \int_{\Omega}(\operatorname{divh})\left(\left(u^{\prime}\right)^{2}-(\Delta u)^{2}\right)+2 \sum_{j=1}^{n}\left(\Delta h_{j}\right)\left(\partial_{j} u\right)(\Delta u) \\
+4 \sum_{i, j=1}^{n}\left(\partial_{i} h_{j}\right)\left(\partial_{i} \partial_{j} u\right)(\Delta u) \mathrm{dx} \mathrm{dt} .
\end{gathered}
$$

Comparing with (20) it remains to show that

$$
\begin{align*}
\int_{S}^{T} \int_{\Gamma} 2\left(\partial_{\nu} \Delta u\right) h \cdot \nabla u & -2(\Delta u) \partial_{\nu}(h \cdot \nabla u)+(h \cdot \nu)\left((\Delta u)^{2}-\left(u^{\prime}\right)^{2}\right) \mathrm{d} \Gamma \mathrm{dt} \\
& =-\int_{S}^{T} \int_{\Gamma}(h \cdot \nu)(\Delta u)^{2} \mathrm{~d} \Gamma \mathrm{dt} \tag{21}
\end{align*}
$$

For this we need the boundary conditions (17).
First, (17) implies that $u^{\prime}=|\nabla u|=0$ on $\Gamma \times \mathbb{R}$ and therefore (21) reduces to

$$
\begin{equation*}
\int_{S}^{T} \int_{\Gamma}(h \cdot \nu)(\Delta u)^{2} \mathrm{~d} \Gamma \mathrm{dt}=\int_{S}^{T} \int_{\Gamma}(\Delta u) \partial_{\nu}(h \cdot \nabla u) \mathrm{d} \Gamma \mathrm{dt} . \tag{22}
\end{equation*}
$$

To compute $\partial_{\nu}(h \cdot \nabla u)$ at a fixed point $x \in \Gamma$, let us choose the coordinates such that $\nu(x)=(0, \ldots, 0,1)$. Since $\partial_{1} u=\cdots=\partial_{n} u=0$ on $\Gamma$, we have

$$
\nabla\left(\partial_{j} u\right)=\left(\partial_{\nu} \partial_{j} u\right) \nu \operatorname{sur} \Gamma, j=1, \ldots, n
$$

In particular, we have

$$
\partial_{i} \partial_{j} u(x)=0, i=1, \ldots, n-1, j=1, \ldots, n .
$$

Since $\partial_{i} \partial_{j} u=\partial_{j} \partial_{i} u$, hence we conclude that

$$
\begin{equation*}
\partial_{i} \partial_{j} u(x)=0 \text { always, unless } i=j=n . \tag{23}
\end{equation*}
$$

Using (23) and the property $\nu(x)=(0, \ldots, 0,1)$ we obtain that

$$
\begin{gathered}
\partial_{\nu}(h \cdot \nabla u)(x)=\left(\sum_{j=1}^{n}\left(\partial_{\nu} h_{j}\right)\left(\partial_{j} u\right)+h_{j} \partial_{n} \partial_{j} u\right)(x) \\
=\sum_{j=1}^{n} h_{j} \partial_{n} \partial_{j} u(x)=h_{n} \partial_{n}^{2} u(x)=h_{n} \Delta u(x)=(h \cdot \nu) \Delta u(x) .
\end{gathered}
$$

Since $x \in \Gamma$ was chosen arbitrarily, we conclude that $\partial_{\nu}(h \cdot \nabla u)=(h \cdot \nu) \Delta u$ on $\Gamma$, proving (22).

Proof of theorem 2.6. - As in $\S 2.2$, it is sufficient to prove the inequalities (19). Let us apply the identity (20) with $S:=-T$ and with the vector field of lemma 2.1, corresponding to the case $\Gamma_{0}=\Gamma$. Then the left-hand side of (23) coincides with that of (22), while the right-hand side of (23) is easily majorized by

$$
c_{1}\left(\|u\|_{C\left([-T, T] ; H^{2}(\Omega)\right)}^{2}+\left\|u^{\prime}\right\|_{C\left([-T, T] ; L^{2}(\Omega)\right)}^{2}\right)
$$

where $c_{1}$ is a constant depending on $\|h\|_{C^{2}(\bar{\Omega})}$ and $T$ only. Applying theorem 1.1 this last expression is majorized by the right-hand side of (19) for a suitable constant $c$.

Remark 2.8. - In the sequel we shall write $\Delta u$ on $\Gamma$ instead of $L\left(u^{0}, u^{1}\right)$ for every $\left(u^{0}, u^{1}\right) \in H_{0}^{2}(\Omega) \times L^{2}(\Omega)$. Then we have the following trace theorem :

$$
\left(u^{0}, u^{1}\right) \in H_{0}^{2}(\Omega) \times L^{2}(\Omega) \Rightarrow \Delta u \in L_{l o c}^{2}\left(\mathbb{R} ; L^{2}(\Gamma)\right)
$$

Now we consider the non-homogeneous boundary value problem

$$
\begin{gather*}
y^{\prime \prime}+\Delta^{2} y=0 \quad \text { in } \quad \Omega \times \mathbb{R}  \tag{24}\\
y=0 \quad \text { and } \quad \partial_{\nu} y=v \quad \text { on } \Gamma \times \mathbb{R}  \tag{25}\\
y(0)=y^{0}, \quad y^{\prime}(0)=y^{1} . \tag{26}
\end{gather*}
$$

As usual, we begin with a formal computation. Let $y \in C^{4}(\bar{\Omega} \times[0, T])$ be a function satisfying (24)-(26). Fix $\left(u^{0}, u^{1}, f\right) \in V \times H=H_{0}^{2}(\Omega) \times L^{2}(\Omega)$ arbitrarily and multiply (16) by $y$. Using the boundary conditions (17) and
(25) we obtain for every fixed $S \in \mathbb{R}$ that

$$
\begin{gathered}
0=\int_{0}^{S} \int_{\Omega}\left(u^{\prime \prime}+\Delta^{2} u\right) y \mathrm{dx} \mathrm{dt} \\
=\left[\int_{\Omega} u^{\prime} y-u y^{\prime} \mathrm{dx}\right]_{0}^{T}+\int_{0}^{S} \int_{\Omega} u\left(y^{\prime \prime}+\Delta^{2} y\right) \mathrm{dx} \mathrm{dt} \\
+\int_{0}^{S} \int_{\Gamma}\left(\partial_{\nu} \Delta u\right) y-(\Delta u)\left(\partial_{\nu} y\right)+\left(\partial_{\nu} u\right)(\Delta y)-u\left(\partial_{\nu} \Delta y\right) \mathrm{d} \Gamma \mathrm{dt} \\
=\int_{\Omega} u^{\prime}(S) y(S)-u(S) y^{\prime}(S)+u^{0} y^{1}-u^{1} y^{0} \mathrm{dx}-\int_{0}^{S} \int_{\Gamma}(\Delta u) v \mathrm{~d} \Gamma \mathrm{dt}
\end{gathered}
$$

Introducing the linear form

$$
L_{S}\left(u^{0}, u^{1}\right):=\int_{0}^{S} \int_{\Gamma}(\Delta u) v \mathrm{~d} \Gamma \mathrm{dt}+\int_{\Omega} u^{1} y^{0}-u^{0} y^{1} \mathrm{dx}
$$

we may rewrite this identity in the following form :

$$
\begin{equation*}
L_{S}\left(u^{0}, u^{1}\right)=\left\langle\left(-y^{\prime}(S), y(S)\right),\left(u(S), u^{\prime}(S)\right)\right\rangle_{V^{\prime} \times H, V \times H} \tag{27}
\end{equation*}
$$

(Note that $V=H_{0}^{2}(\Omega), H=L^{2}(\Omega)$ and $V^{\prime}=H^{-2}(\Omega)$.)
Definition. - We say that $\left(y, y^{\prime}\right)$ is a solution of (24)-(26) if $\left(y, y^{\prime}\right) \in C\left(\mathbb{R} ; H \times V^{\prime}\right)$ and if $(27)$ is satisfied for all $S \in \mathbb{R}$ and for all $\left(u^{0}, u^{1}\right) \in V \times H$.

To justify this definition we prove the
Theorem 2.9. - Given $\left(y^{0}, y^{1}\right) \in L^{2}(\Omega) \times H^{-2}(\Omega)$ and $v \in L_{\text {loc }}^{2}\left(\mathbb{R} ; L^{2}(\Gamma)\right)$ arbitrarily, the problem (24)-(26) has a unique solution.

Furthermore, the linear map $\left(y^{0}, y^{1}, v\right) \mapsto\left(y, y^{\prime}\right)$ is continuous with respect to these topologies.

Proof. - It follows from theorem 2.6 that for every fixed $S \in \mathbb{R}$ the linear form $L_{S}$ is bounded in $V \times H$. Furthermore, it follows from theorem 1.1 that the linear map

$$
\left(u(S), u^{\prime}(S)\right) \mapsto\left(u^{0}, u^{1}\right)
$$

is an isomorphism of $V \times H$ onto itself. Hence the linear form

$$
\left(u(S), u^{\prime}(S)\right) \mapsto L_{S}\left(u^{0}, u^{1}\right)
$$

is also bounded on $V \times H$. Therefore there is a unique pair $\left(y(S), y^{\prime}(S)\right) \in$ $H \times V^{\prime}$ satisfying (27).

The rest of the theorem may be proved exactly as the corresponding part of theorem 2.5.

### 2.4. The second Petrovsky system

Here we consider the problem introduced in § 1.5 :

$$
\begin{gather*}
u^{\prime \prime}+\Delta^{2} u=0 \quad \text { in } \quad \Omega \times \mathbb{R}  \tag{28}\\
u=\Delta u=0 \quad \text { on } \quad \Gamma \times \mathbb{R}  \tag{29}\\
u(0)=u^{0} \quad \text { and } \quad u^{\prime}(0)=u^{1} . \tag{30}
\end{gather*}
$$

If $u^{0} \in H^{2}(\Omega) \cap H_{0}^{1}(\Omega)(=V)$ and $u^{1} \in L^{2}(\Omega)(=H)$, then by theorem 1.1 this problem has a unique solution satisfying

$$
u \in C\left(\mathbb{R} ; H^{2}(\Omega) \cap H_{0}^{1}(\Omega)\right) \cap C^{1}\left(\mathbb{R} ; L^{2}(\Omega)\right)
$$

In particular, $\partial_{\nu} u$ is well-defined. The following result permits us to define the normal derivative of weaker solutions, too.

Theorem 2.10. - Given $T>0$ arbitrarily, there exists a constant $c=c(T)>0$ such that for every $\left(u^{0}, u^{1}\right) \in Z \times Z$ the solution of (28)-(30) satisfies the inequality

$$
\begin{equation*}
\int_{-T}^{T} \int_{\Gamma}\left|\partial_{\nu} u\right|^{2} \mathrm{~d} \Gamma \mathrm{dt} \leq c\left(\left\|u^{0}\right\|_{H_{0}^{1}(\Omega)}^{2}+\left\|u^{1}\right\|_{H^{-1}(\Omega)}^{2}\right) \tag{31}
\end{equation*}
$$

Consequently, there is a unique continuous linear map

$$
L: H_{0}^{1}(\Omega) \times H^{-1}(\Omega) \rightarrow L_{l o c}^{2}\left(\mathbb{R} ; L^{2}(\Gamma)\right)
$$

such that

$$
L\left(u^{0}, u^{1}\right)=\partial_{\nu} u, \forall\left(u^{0}, u^{1}\right) \in Z \times Z .
$$

Remark 2.11. - Let us recall from lemma 1.7 that $H_{0}^{1}(\Omega)=D_{1 / 4}$ and $H^{-1}(\Omega)=D_{-1 / 4}$. Thus for any given $\left(u^{0}, u^{1}\right) \in H_{0}^{1}(\Omega) \times H^{-1}(\Omega)(28)-(30)$ has a unique solution such that

$$
u \in C\left(\mathbb{R} ; H_{0}^{1}(\Omega)\right) \cap C^{1}\left(\mathbb{R} ; H^{-1}(\Omega)\right)
$$

(See remark 1.3.) Therefore, writing $\partial_{\nu} u$ instead of $L\left(u^{0}, u^{1}\right)$ we will have the following hidden regularity result :

$$
\left(u^{0}, u^{1}\right) \in H_{0}^{1}(\Omega) \times H^{-1}(\Omega) \Rightarrow \partial_{\nu} u \in L_{l o c}^{2}\left(\mathbb{R} ; L^{2}(\Gamma)\right) .
$$

Lemma 2.12. - Let $u \in H_{l o c}^{2}\left(\mathbb{R} ; H^{4}(\Omega)\right)$ be a function satisfying (28), (29) and let $h: \bar{\Omega} \rightarrow \mathbb{R}^{n}$ be a vector field of class $C^{2}$. Then for any given $-\infty<S<T<\infty$ the following identity holds :

$$
\begin{gather*}
\int_{S}^{T} \int_{\Gamma}(h \cdot \nu)\left(\left|\partial_{\nu} u^{\prime}\right|^{2}+\left|\partial_{\nu} \Delta u\right|^{2}\right) \mathrm{d} \Gamma \mathrm{dt}=-\left[\int_{\Omega} 2 u^{\prime} h \cdot \nabla \Delta u \mathrm{dx}\right]_{S}^{T} \\
+\int_{S}^{T} \int_{\Omega}(\operatorname{div} h)\left(\left|\nabla u^{\prime}\right|^{2}-|\nabla \Delta u|^{2}\right)  \tag{32}\\
+2 \sum_{i, j=1}^{n}\left(\partial_{i} h_{j}\right)\left(\partial_{i} u^{\prime}\right)\left(\partial_{j} u^{\prime}\right)+\left(\partial_{i} h_{j}\right)\left(\partial_{i} \Delta u\right)\left(\partial_{j} \Delta u\right)+2 \sum_{i=1}^{n}\left(\Delta h_{i}\right) u^{\prime} \partial_{i} u^{\prime} \mathrm{dx} \mathrm{dt} .
\end{gather*}
$$

Proof. - We multiply (28) by $2 h \cdot \nabla \Delta u$ and we integrate by parts. We obtain

$$
\begin{gather*}
0=\int_{S}^{T} \int_{\Omega} 2\left(u^{\prime \prime}+\Delta^{2} u\right) h \cdot \nabla \Delta u \mathrm{dx} \mathrm{dt} \\
=\left[\int_{\Omega} 2 u^{\prime} h \cdot \nabla \Delta u \mathrm{dx}\right]_{S}^{T}+\int_{S}^{T} \int_{\Omega}-2 u^{\prime} h \cdot \nabla \Delta u^{\prime}+2\left(\Delta^{2} u\right) h \cdot \nabla \Delta u \mathrm{dx} \mathrm{dt} \\
=\left[\int_{\Omega} 2 u^{\prime} h \cdot \nabla \Delta u \mathrm{dx}\right]_{S}^{T}+2 \int_{S}^{T} \int_{\Gamma}\left(\partial_{\nu} \Delta u\right) h \cdot \nabla \Delta u \mathrm{~d} \Gamma \mathrm{dt} \\
-\int_{S}^{T} \int_{\Omega} 2 u^{\prime} h \cdot \nabla \Delta u^{\prime}+2(\nabla \Delta u) \cdot \nabla(h \cdot \nabla \Delta u) \mathrm{dx} \mathrm{dt}  \tag{33}\\
=\left[\int_{\Omega} 2 u^{\prime} h \cdot \nabla \Delta u \mathrm{dx}\right]_{S}^{T}+2 \int_{S}^{T} \int_{\Gamma}(h \cdot \nu)\left|\partial_{\nu} \Delta u\right|^{2} \mathrm{~d} \Gamma \mathrm{dt} \\
-\int_{S}^{T} \int_{\Omega} 2 u^{\prime} h \cdot \nabla \Delta u^{\prime}+2(\nabla \Delta u) \cdot \nabla(h \cdot \nabla \Delta u) \mathrm{dx} \mathrm{dt}
\end{gather*}
$$

because

$$
\Delta u=0 \quad \text { on } \quad \Gamma \Rightarrow \nabla \Delta u=\left(\partial_{\nu} \Delta u\right) \nu \quad \text { on } \quad \Gamma .
$$

Furthermore, we have

$$
\int_{\Omega}-2 u^{\prime} h \cdot \nabla \Delta u^{\prime} \mathrm{dx}=\sum_{i, j=1}^{n} \int_{\Omega}-2 u^{\prime} h_{i} \partial_{i} \partial_{j}^{2} u^{\prime} \mathrm{dx}
$$

$$
\begin{gathered}
=-2 \sum_{i, j=1}^{n} \int_{\Gamma} u^{\prime} h_{i} \nu_{j} \partial_{i} \partial_{j} u^{\prime} \mathrm{d} \Gamma+2 \sum_{i, j=1}^{n} \int_{\Omega} \partial_{j} u^{\prime} h_{i} \partial_{i} \partial_{j} u^{\prime}+\left(\partial_{j} h_{i}\right) u^{\prime}\left(\partial_{i} \partial_{j} u^{\prime}\right) \mathrm{dx} \\
=-2 \sum_{\Gamma}^{n} \int^{n} u^{\prime} h_{i} \nu_{j} \partial_{i} \partial_{j} u^{\prime} \mathrm{d} \Gamma+\int_{\Omega} h \cdot \nabla\left|\nabla u^{\prime}\right|^{2}+2 \sum_{i, j=1}^{n}\left(\partial_{j} h_{i}\right) u^{\prime}\left(\partial_{i} \partial_{j} u^{\prime}\right) \mathrm{dx} \\
=-2 \sum_{i, j=1}^{n} \int_{\Gamma} u^{\prime} h_{i} \nu_{j} \partial_{i} \partial_{j} u^{\prime} \mathrm{d} \Gamma+\int_{\Gamma}(h \cdot \nu)\left|\nabla u^{\prime}\right|^{2} \mathrm{~d} \Gamma \\
\quad-\int_{\Omega}(\operatorname{div} h)\left|\nabla u^{\prime}\right|^{2} \mathrm{dx}+2 \sum_{i, j=1}^{n} \int_{\Gamma} \nu_{j}\left(\partial_{j} h_{i}\right) u^{\prime}\left(\partial_{i} u^{\prime}\right) \mathrm{d} \Gamma \\
\quad-2 \sum_{i=1}^{n} \int_{\Omega}\left(\Delta h_{i}\right) u^{\prime}\left(\partial_{i} u^{\prime}\right) \mathrm{dx}-2 \sum_{i, j=1}^{n} \int_{\Omega}\left(\partial_{j} h_{i}\right)\left(\partial_{j} u^{\prime}\right)\left(\partial_{i} u^{\prime}\right) \mathrm{dx} .
\end{gathered}
$$

Since the condition $u=0$ on $\Gamma$ implies that $u^{\prime}=0$ and $\nabla u^{\prime}=\left(\partial_{\nu} u^{\prime}\right) \nu$ on $\Gamma$, the first and third integrals on $\Gamma$ vanish, and the second one is equal to

$$
\int_{\Gamma}(h \cdot \nu)\left|\partial_{\nu} u^{\prime}\right|^{2} \mathrm{~d} \Gamma
$$

Therefore we have

$$
\begin{gathered}
\int_{\Omega}-2 u^{\prime} h \cdot \nabla \Delta u^{\prime} \mathrm{dx}=\int_{\Gamma}(h \cdot \nu)\left|\partial_{\nu} u^{\prime}\right|^{2} \mathrm{~d} \Gamma-\int_{\Omega}(\operatorname{div} h)\left|\nabla u^{\prime}\right|^{2} \mathrm{dx} \\
-2 \int_{\Omega} \sum_{i=1}^{n}\left(\Delta h_{i}\right) u^{\prime}\left(\partial_{i} u^{\prime}\right)+\sum_{i, j=1}^{n}\left(\partial_{j} h_{i}\right)\left(\partial_{j} u^{\prime}\right)\left(\partial_{i} u^{\prime}\right) \mathrm{dx}
\end{gathered}
$$

Substituting into (33) we obtain that

$$
\begin{gather*}
0=\left[\int_{\Omega} 2 u^{\prime} h \cdot \nabla \Delta u \mathrm{dx}\right]_{S}^{T}+\int_{S}^{T} \int_{\Gamma} 2(h \cdot \nu)\left|\partial_{\nu} \Delta u\right|^{2}+(h \cdot \nu)\left|\partial_{\nu} u^{\prime}\right|^{2} \mathrm{~d} \Gamma \mathrm{dt} \\
\quad-\int_{S}^{T} \int_{\Omega}(\operatorname{div} h)\left|\nabla u^{\prime}\right|^{2}+2 \sum_{i, j=1}^{n}\left(\partial_{j} h_{i}\right)\left(\partial_{j} u^{\prime}\right)\left(\partial_{i} u^{\prime}\right)  \tag{34}\\
\quad+2 \sum_{i=1}^{n}\left(\Delta h_{i}\right) u^{\prime}\left(\partial_{i} u^{\prime}\right)+2 \nabla \Delta u \cdot \nabla(h \cdot \nabla \Delta u) \mathrm{dx} \mathrm{dt}
\end{gather*}
$$

Let us consider the last term of (34). We have

$$
\begin{gathered}
-\int_{\Omega} 2 \nabla \Delta u \cdot \nabla(h \cdot \nabla \Delta u) \mathrm{dx}=\sum_{i, j=1}^{n} \int_{\Omega}-2\left(\partial_{j} \Delta u\right) \partial_{j}\left(h_{i} \partial_{i} \Delta u\right) \mathrm{dx} \\
=\sum_{i, j=1}^{n} \int_{\Omega}-2\left(\partial_{j} h_{i}\right)\left(\partial_{j} \Delta u\right)\left(\partial_{i} \Delta u\right)-h_{i} \partial_{i}\left|\partial_{j} \Delta u\right|^{2} \mathrm{dx}
\end{gathered}
$$

$$
\begin{gathered}
=\sum_{i, j=1}^{n} \int_{\Omega}-2\left(\partial_{j} h_{i}\right)\left(\partial_{j} \Delta u\right)\left(\partial_{i} \Delta u\right) \mathrm{dx}+\int_{\Omega}(\operatorname{divh})|\nabla \Delta u|^{2} \mathrm{dx} \\
-\int_{\Gamma}(h \cdot \nu)|\nabla \Delta u|^{2} \mathrm{~d} \Gamma
\end{gathered}
$$

Since $\Delta u=0$ on $\Gamma$, in the last integral we may write $\left|\partial_{\nu} \Delta u\right|$ instead of $|\nabla \Delta u|$. Using this equality we deduce from (34) that

$$
\begin{gathered}
0=\left[\int_{\Omega} 2 u^{\prime} h \cdot \nabla \Delta u \mathrm{dx}\right]_{S}^{T}+\int_{S}^{T} \int_{\Gamma}(h \cdot \nu)\left|\partial_{\nu} \Delta u\right|^{2}+(h \cdot \nu)\left|\partial_{\nu} u^{\prime}\right|^{2} \mathrm{~d} \Gamma \mathrm{dt} \\
\quad+\int_{S}^{T} \int_{\Omega}(\operatorname{div} h)\left(|\nabla \Delta u|^{2}-\left|\nabla u^{\prime}\right|^{2}\right) \mathrm{dx} \mathrm{dt} \\
-2 \int_{S}^{T} \int_{\Omega} \sum_{i, j=1}^{n}\left(\partial_{j} h_{i}\right)\left(\left(\partial_{j} u^{\prime}\right)\left(\partial_{i} u^{\prime}\right)+\left(\partial_{j} \Delta u\right)\left(\partial_{i} \Delta u\right)\right) \mathrm{dx} \mathrm{dt} \\
\quad-2 \int_{S}^{T} \int_{\Omega} \sum_{i=1}^{n} u^{\prime}\left(\Delta h_{i}\right)\left(\partial_{i} u^{\prime}\right) \mathrm{dx} \mathrm{dt}
\end{gathered}
$$

and (32) follows.
Theorem 2.13. - Given $T>0$ arbitrarily, there exists a constant $c=c(T)>0$ such that for every $\left(u^{0}, u^{1}\right) \in Z \times Z$ the solution of (28)-(30) satisfies the inequality

$$
\begin{equation*}
\int_{-T}^{T} \int_{\Gamma}\left|\partial_{\nu} u^{\prime}\right|^{2}+\left|\partial_{\nu} \Delta u\right|^{2} \mathrm{~d} \Gamma \mathrm{dt} \leq c\left(\left\|\nabla \Delta u^{0}\right\|_{L^{2}(\Omega)}^{2}+\left\|\nabla u^{1}\right\|_{L^{2}(\Omega)}^{2}\right) \tag{35}
\end{equation*}
$$

Proof. - Applying the identity (32) with $S=-T$ and with the vector field of lemma 2.1 corresponding to $\Gamma_{0}=\Gamma$, we obtain the estimate

$$
\begin{equation*}
\int_{-T}^{T} \int_{\Gamma}\left|\partial_{\nu} u^{\prime}\right|^{2}+\left|\partial_{\nu} \Delta u\right|^{2} \mathrm{~d} \Gamma \mathrm{dt} \leq c_{1}\left\|\left(u, u^{\prime}\right)\right\|_{C\left([-T, T] ; H^{3}(\Omega) \times H^{1}(\Omega)\right)}^{2} \tag{36}
\end{equation*}
$$

for some constant $c_{1}$. Using lemma 1.7 and theorem 1.1 hence (35) follows.
Proof of theorem 2.10. - It suffices to prove the inequalities (31). Given $\left(u^{0}, u^{1}\right) \in Z \times Z$ arbitrarily, let us apply the estimate (35) of theorem 2.13 with

$$
\left(A^{-1 / 2} u^{0}, A^{-1 / 2} u^{1}\right)
$$

instead of $\left(u^{0}, u^{1}\right)$. We obtain that (see also remark 1.3)

$$
\begin{align*}
& \int_{-T}^{T} \int_{\Gamma}\left|\partial_{\nu} A^{-1 / 2} u^{\prime}\right|^{2}+\left|\partial_{\nu} \Delta A^{-1 / 2} u\right|^{2} \mathrm{~d} \Gamma \mathrm{dt} \\
& \leq c\left(\left\|\nabla \Delta A^{-1 / 2} u^{0}\right\|_{L^{2}(\Omega)}^{2}+\left\|\nabla A^{-1 / 2} u^{1}\right\|_{L^{2}(\Omega)}^{2}\right. \tag{37}
\end{align*}
$$

Using lemma 1.7 we have

$$
\begin{gathered}
\Delta A^{-1 / 2} v=-v, \\
\|\nabla v\|_{L^{2}(\Omega)}=\|v\|_{1 / 4}=\|v\|_{H_{0}^{1}(\Omega)}
\end{gathered}
$$

and

$$
\left\|\nabla A^{-1 / 2} v\right\|_{L^{2}(\Omega)}=\left\|A^{-1 / 2} v\right\|_{1 / 4}=\|v\|_{-1 / 4}=\|v\|_{H^{-1}(\Omega)}
$$

for every $v \in Z$. Therefore we deduce from (37) that

$$
\int_{-T}^{T} \int_{\Gamma}\left|\partial_{\nu} A^{-1 / 2} u^{\prime}\right|^{2}+\left|\partial_{\nu} u\right|^{2} \mathrm{~d} \Gamma \mathrm{dt} \leq c\left(\left\|u^{0}\right\|_{H_{0}^{1}(\Omega)}^{2}+\left\|u^{1}\right\|_{H^{-1}(\Omega)}^{2}\right)
$$

Omitting the term $\left|\partial_{\nu} A^{-1 / 2} u^{\prime}\right|^{2}$ we find (31).
Now let us apply theorem 2.10 for the study of the non-homogeneous problem

$$
\begin{gather*}
y^{\prime \prime}+\Delta^{2} y=0 \quad \text { in } \quad \Omega \times \mathbb{R}  \tag{38}\\
y=0 \quad \text { and } \quad \Delta y=v \quad \text { on } \quad \Gamma \times \mathbb{R}  \tag{39}\\
y(0)=y^{0}, \quad y^{\prime}(0)=y^{1} . \tag{40}
\end{gather*}
$$

Let $\left(u^{0}, u^{1}\right) \in H_{0}^{1}(\Omega) \times H^{-1}(\Omega)$ and consider the solution of (28)-(30). If $y \in C^{4}(\bar{\Omega} \times[0, T])$ satisfies (38)-(40), then multiplying (28) by $y$ and integrating by parts we find for every fixed $S \in \mathbb{R}$ that

$$
\begin{gathered}
0=\int_{0}^{S} \int_{\Omega}\left(u^{\prime \prime}+\Delta^{2} u\right) y \mathrm{dx} \mathrm{dt} \\
=\left[\int_{\Omega} u^{\prime} y-u y^{\prime} \mathrm{dx}\right]_{0}^{S}+\int_{0}^{S} \int_{\Omega} u\left(y^{\prime \prime}+\Delta^{2} y\right) \mathrm{dx} \mathrm{dt} \\
+\int_{0}^{S} \int_{\Gamma}\left(\partial_{\nu} \Delta u\right) y-(\Delta u)\left(\partial_{\nu} y\right)+\left(\partial_{\nu} u\right)(\Delta y)-u\left(\partial_{\nu} \Delta y\right) \mathrm{d} \Gamma \mathrm{dt} \\
=\int_{\Omega} u^{\prime}(S) y(S)-u(T) y^{\prime}(S)+u^{0} y^{1}-u^{1} y^{0} \mathrm{dx}+\int_{0}^{S} \int_{\Gamma}\left(\partial_{\nu} u\right) v \mathrm{~d} \Gamma \mathrm{dt} .
\end{gathered}
$$

Setting

$$
L_{S}\left(u^{0}, u^{1}\right):=\int_{0}^{S} \int_{\Gamma}\left(\partial_{\nu} u\right) v \mathrm{~d} \Gamma \mathrm{dt}+\int_{\Omega} u^{0} y^{1}-u^{1} y^{0} \mathrm{dx}
$$

and writing for brevity $X:=H_{0}^{1}(\Omega) \times H^{-1}(\Omega)$ we may rewrite this identity in the following form :

$$
\begin{equation*}
L_{S}\left(u^{0}, u^{1}\right)=\left\langle\left(y^{\prime}(S),-y(S)\right),\left(u(S), u^{\prime}(S)\right)\right\rangle_{X^{\prime}, X}, \quad \forall\left(u^{0}, u^{1}\right) \in X \tag{41}
\end{equation*}
$$

Definition. - We say that $\left(y, y^{\prime}\right)$ is a solution of (38)-(40) if $\left(y, y^{\prime}\right) \in C\left(\mathbb{R} ; H_{0}^{1}(\Omega) \times H^{-1}(\Omega)\right)$ and if (41) is satisfied for all $S \in \mathbb{R}$.

Then we have
Theorem 2.14. - Let $\left(y^{0}, y^{1}\right) \in H_{0}^{1}(\Omega) \times H^{-1}(\Omega)$ and $v \in L_{l o c}^{2}\left(\mathbb{R} ; L^{2}(\Gamma)\right)$ be given arbitrarily. Then the problem (38)-(40) has a unique solution.

Furthermore, the linear map $\left(y^{0}, y^{1}, v\right) \mapsto\left(y, y^{\prime}\right)$ is continuous with respect to these topologies.

Proof. - It follows from theorem 2.10 that for every fixed $S \in \mathbb{R}$ the linear form $L_{S}$ is bounded in $X=H_{0}^{1}(\Omega) \times H^{-1}(\Omega)$. Furthermore, it follows from theorem 1.1 that the linear map

$$
\left(u(S), u^{\prime}(S)\right) \mapsto\left(u^{0}, u^{1}\right)
$$

is an isomorphism of $X$ onto itself. Hence the linear form

$$
\left(u(S), u^{\prime}(S)\right) \mapsto L_{S}\left(u^{0}, u^{1}\right)
$$

is also bounded on $X$. Therefore there is a unique pair $\left(y(S), y^{\prime}(S)\right) \in X$ satisfying (41).

The proof now can be completed as that of theorem 2.5 before; we omit the details.

## 3. Uniqueness theorems

The results of this chapter will serve as a basis for the Hilbert Uniqueness Method (HUM) of J.-L. Lions (see the following chapter). In control-theoretic terminology, we shall here establish observability theorems; see remark 3.5 below. These will be obtained (following Ho [1]) by using a special multiplier. The same multiplier has already been used by many authors for different reasons : see e.g. Rellich [1], Pohožaev [1], Lax, Morawetz and Phillips [1], Chen [1].

As in the preceding chapter, we consider the real case only : the complex case then follows easily by applying the results to the real and imaginary parts of the solutions.

We shall use the following notation : for any fixed $x^{0} \in \mathbb{R}^{n}$ we set

$$
\begin{gather*}
m(x):=x-x^{0} \quad\left(x \in \mathbb{R}^{n}\right),  \tag{1}\\
R=R\left(x^{0}\right):=\sup _{x \in \Omega}|m(x)|,  \tag{2}\\
\Gamma_{+}:=\{x \in \Gamma: m \cdot \nu>0\}, \tag{3}
\end{gather*}
$$

and we introduce on $\Gamma$ the (signed) surface measure

$$
\begin{equation*}
\mathrm{d} \Gamma_{m}:=(m \cdot \nu) \mathrm{d} \Gamma ; \tag{4}
\end{equation*}
$$

clearly, we have

$$
\left|\mathrm{d} \Gamma_{m}\right| \leq R \mathrm{~d} \Gamma
$$

Let us note the following obvious relations :

$$
\begin{equation*}
\partial_{i} m_{j}=\delta_{i j} \text { and } \operatorname{div} m \equiv n . \tag{5}
\end{equation*}
$$

Here and in the rest of this book we shall always mean by an interval a bounded interval of strictly positive length.

### 3.1. The wave equation. Dirichlet condition

Consider the problem of sections 1.3 and 2.2 with $\Gamma_{0}=\Gamma$ and $\Gamma_{1}=\emptyset:$

$$
\begin{gather*}
u^{\prime \prime}-\Delta u+q u=0 \quad \text { in } \Omega \times \mathbb{R}  \tag{6}\\
u=0 \quad \text { on } \quad \Gamma \times \mathbb{R}  \tag{7}\\
u(0)=u^{0} \quad \text { and } \quad u^{\prime}(0)=u^{1} . \tag{8}
\end{gather*}
$$

Recall (see theorems 1.1 and 2.2) that for every $\left(u^{0}, u^{1}\right) \in H_{0}^{1}(\Omega) \times L^{2}(\Omega)$ this problem has a unique solution, whose energy

$$
\begin{equation*}
E:=\frac{1}{2} \int_{\Omega}\left|u^{\prime}\right|^{2}+|\nabla u|^{2}+q|u|^{2} \mathrm{dx} \tag{9}
\end{equation*}
$$

is conserved, and that for every interval $I$ the following (so-called direct) inequality holds true:

$$
\begin{equation*}
\int_{I} \int_{\Gamma}\left|\partial_{\nu} u\right|^{2} \mathrm{~d} \Gamma \mathrm{dt} \leq c E \tag{10}
\end{equation*}
$$

The aim of this section is to establish, under some further hypotheses, the inverse inequality

$$
\begin{equation*}
\int_{I} \int_{\Gamma}\left|\partial_{\nu} u\right|^{2} \mathrm{~d} \Gamma \mathrm{dt} \geq c^{\prime} E \tag{11}
\end{equation*}
$$

Fix $x^{0} \in \mathbb{R}^{n}$ arbitrarily and set $Q:=\sup _{\Omega} q$ and

$$
Q_{1}:= \begin{cases}2 R Q / \sqrt{\lambda_{1}}, & \text { if } n \geq 2 \\ 2 R Q / \sqrt{\lambda_{1}}+Q / \lambda_{1}, & \text { if } n=1\end{cases}
$$

(We recall from $\S 1.3$ that the function $q$ is supposed to be nonnegative, bounded and measurable.) Let us recall that $\lambda_{1}$ is the biggest constant such that

$$
\begin{equation*}
\int_{\Omega}|\nabla v|^{2}+q|v|^{2} \mathrm{dx} \geq \lambda_{1} \int_{\Omega}|v|^{2} \mathrm{dx}, \forall v \in H_{0}^{1}(\Omega) \tag{12}
\end{equation*}
$$

Theorem 3.1. - Assume that

$$
\begin{equation*}
Q_{1}<1 \tag{13}
\end{equation*}
$$

and let I be an interval of length

$$
\begin{equation*}
|I|>2 R /\left(1-Q_{1}\right) \tag{14}
\end{equation*}
$$

Then there is a constant $c^{\prime}>0$ such that

$$
\begin{equation*}
\int_{I} \int_{\Gamma}\left|\partial_{\nu} u\right|^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt} \geq c^{\prime} E, \forall\left(u^{0}, u^{1}\right) \in H_{0}^{1}(\Omega) \times L^{2}(\Omega) \tag{15}
\end{equation*}
$$

Clearly the estimate (15) implies (11) (with $c^{\prime} / R$ instead of $c^{\prime}$ ).
For $q \equiv 0$ this theorem was proved by Ho [1] under a stronger condition on the length of $I$; his condition was weakened by Lions [2], [3], using an
indirect compactness-uniqueness argument. The following constructive proof, based on lemma 3.2 below, was first given in Komornik [1].

Proof. - It is sufficient to prove the estimate (15) for $u^{0}, u^{1} \in Z$ : the general case then follows by a density argument, using theorem 2.2.

Let us write $I=[S, T]$. Applying lemma 2.3 with $h=m$ and using (5) we obtain the identity

$$
\begin{gather*}
\int_{S}^{T} \int_{\Gamma} 2 \partial_{\nu} u m \cdot \nabla u+(m \cdot \nu)\left(\left(u^{\prime}\right)^{2}-|\nabla u|^{2}\right) \mathrm{d} \Gamma \mathrm{dt} \\
=\left[\int_{\Omega} 2 u^{\prime} m \cdot \nabla u \mathrm{dx}\right]_{S}^{T}  \tag{16}\\
+\int_{S}^{T} \int_{\Omega} n\left(u^{\prime}\right)^{2}+(2-n)|\nabla u|^{2}+2 q u m \cdot \nabla u \mathrm{dx} \mathrm{dt} .
\end{gather*}
$$

Let us multiply the equation (6) by $u$; integrating by parts we obtain

$$
\begin{gathered}
0=\int_{S}^{T} \int_{\Omega} u\left(u^{\prime \prime}-\Delta u+q u\right) \mathrm{dx} \mathrm{dt}=\left[\int_{\Omega} u u^{\prime} \mathrm{dx}\right]_{S}^{T} \\
-\int_{S}^{T} \int_{\Gamma} u \partial_{\nu} u \mathrm{~d} \Gamma \mathrm{dt}+\int_{S}^{T} \int_{\Omega}-\left(u^{\prime}\right)^{2}+|\nabla u|^{2}+q u^{2} \mathrm{dx} \mathrm{dt}
\end{gathered}
$$

whence

$$
\begin{equation*}
\int_{S}^{T} \int_{\Gamma} u \partial_{\nu} u \mathrm{~d} \Gamma \mathrm{dt}=\left[\int_{\Omega} u u^{\prime} \mathrm{dx}\right]_{S}^{T}+\int_{S}^{T} \int_{\Omega}-\left(u^{\prime}\right)^{2}+|\nabla u|^{2}+q u^{2} \mathrm{dx} \mathrm{dt} . \tag{17}
\end{equation*}
$$

Putting

$$
M u:=2 m \cdot \nabla u+(n-1) u
$$

we deduce from (16) and (17) that

$$
\begin{gather*}
\int_{S}^{T} \int_{\Gamma}\left(\partial_{\nu} u\right) M u+(m \cdot \nu)\left(\left(u^{\prime}\right)^{2}-|\nabla u|^{2}\right) \mathrm{d} \Gamma \mathrm{dt} \\
=\left[\int_{\Omega} u^{\prime} M u \mathrm{dx}\right]_{S}^{T}  \tag{18}\\
+\int_{S}^{T} \int_{\Omega}\left(u^{\prime}\right)^{2}+|\nabla u|^{2}+(n-1) q u^{2}+2 q u m \cdot \nabla u \mathrm{dx} \mathrm{dt}
\end{gather*}
$$

Using (9) we may rewrite (18) in the following form :

$$
\begin{gather*}
\int_{S}^{T} \int_{\Gamma}\left(\partial_{\nu} u\right) M u+(m \cdot \nu)\left(\left(u^{\prime}\right)^{2}-|\nabla u|^{2}\right) \mathrm{d} \Gamma \mathrm{dt}  \tag{19}\\
=\left[\int_{\Omega} u^{\prime} M u \mathrm{dx}\right]_{S}^{T}+2|I| E+\int_{S}^{T} \int_{\Omega}(n-2) q u^{2}+2 q u m \cdot \nabla u \mathrm{dx} \mathrm{dt} .
\end{gather*}
$$

By the boundary conditions (7) we have $u=u^{\prime}=0$ and $\nabla u=\left(\partial_{\nu} u\right) \nu$ on $\Gamma$; this permits to reduce the left-hand side of (19) to $\int_{I} \int_{\Gamma}\left|\partial_{\nu} u\right|^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt}$.

On the other hand, using (12) the last integral in (19) is minorized by $-2|I| Q_{1} E$. Indeed, the case $q \equiv 0$ is trivial. If $q \not \equiv 0$ and $n \geq 2$, then we have

$$
\begin{gathered}
\int_{\Omega}(n-2) q u^{2}+2 q u m \cdot \nabla u \mathrm{dx} \geq-2 R Q\|u\|_{L^{2}(\Omega)}\|\nabla u\|_{L^{2}(\Omega)} \\
\geq-Q_{1} \int_{\Omega}|\nabla u|^{2}+q u^{2} \mathrm{dx} \geq-2 Q_{1} E
\end{gathered}
$$

If $q \not \equiv 0$ and $n=1$, then we have

$$
\int_{\Omega}(n-2) q u^{2}+2 q u m \cdot \nabla u \mathrm{dx} \geq \int_{\Omega}-Q u^{2}+2 q u m \cdot \nabla u \mathrm{dx}
$$

repeating the above computation and using the definition of $Q_{1}$ for $n=1$, this integral is again minorized by $-2 Q_{1} E$.

Thus we arrive at the following inequality :

$$
\begin{equation*}
\int_{I} \int_{\Gamma}\left|\partial_{\nu} u\right|^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt} \geq 2|I|\left(1-Q_{1}\right) E+\left[\int_{\Omega} u^{\prime} M u \mathrm{dx}\right]_{S}^{T} \tag{20}
\end{equation*}
$$

To minorize the last term of (20) we need the
Lemma 3.2. - The solution of (6)-(8) satisfies the estimate

$$
\begin{equation*}
\left|\int_{\Omega} u^{\prime} M u d x\right| \leq 2 R E, \forall t \in \mathbb{R} \tag{21}
\end{equation*}
$$

Proof. - First we show that

$$
\begin{equation*}
\|M u\|_{L^{2}(\Omega)} \leq\|2 m \cdot \nabla u\|_{L^{2}(\Omega)} \tag{22}
\end{equation*}
$$

Indeed, the application of Green's formula gives

$$
\begin{aligned}
& \|M u\|_{L^{2}(\Omega)}^{2}-\|2 m \cdot \nabla u\|_{L^{2}(\Omega)}^{2}=\|2 m \cdot \nabla u+(n-1) u\|_{L^{2}(\Omega)}^{2}-\|2 m \cdot \nabla u\|_{L^{2}(\Omega)}^{2} \\
& =\int_{\Omega}|2 m \cdot \nabla u+(n-1) u|^{2}-|2 m \cdot \nabla u|^{2} \mathrm{dx}=\int_{\Omega}(n-1)^{2} u^{2}+4(n-1) u m \cdot \nabla u \mathrm{dx} \\
& =\int_{\Omega}(n-1)^{2} u^{2}+2(n-1) m \cdot \nabla\left(u^{2}\right) \mathrm{dx} \\
& =2(n-1) \int_{\Gamma}(m \cdot \nu) u^{2} \mathrm{~d} \Gamma+\int_{\Omega}(n-1)^{2} u^{2}-2(n-1)(\operatorname{div} m) u^{2} \mathrm{dx} \\
& =2(n-1) \int_{\Gamma}(m \cdot \nu) u^{2} \mathrm{~d} \Gamma+\left(1-n^{2}\right) \int_{\Omega} u^{2} \mathrm{dx}=\left(1-n^{2}\right) \int_{\Omega} u^{2} \mathrm{dx} \leq 0 .
\end{aligned}
$$

(We used here (5) and (7).)
It follows from (22) and from the definition of the energy that

$$
\begin{aligned}
& \qquad\left|\int_{\Omega} u^{\prime} M u \mathrm{dx}\right| \leq\left\|u^{\prime}\right\|_{L^{2}(\Omega)}\|M u\|_{L^{2}(\Omega)} \\
& \leq\left\|u^{\prime}\right\|_{L^{2}(\Omega)}\|2 m \cdot \nabla u\|_{L^{2}(\Omega)} \leq R\left\|u^{\prime}\right\|_{L^{2}(\Omega)}^{2}+\frac{1}{R}\|m \cdot \nabla u\|_{L^{2}(\Omega)}^{2} \\
& =\int_{\Omega} R\left|u^{\prime}\right|^{2}+R^{-1}|m \cdot \nabla u|^{2} \mathrm{dx} \leq R \int_{\Omega}\left|u^{\prime}\right|^{2}+|\nabla u|^{2} \mathrm{dx} \leq 2 R E,
\end{aligned}
$$

proving (21).
Applying the lemma with $t=S$ and $t=T$ we deduce from (20) that

$$
\int_{I} \int_{\Gamma}\left|\partial_{\nu} u\right|^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt} \geq 2|I|\left(1-Q_{1}\right) E-4 R E
$$

and the theorem follows with

$$
c^{\prime}:=2|I|\left(1-Q_{1}\right)-4 R .
$$

Remark 3.3. - Let $k$ an arbitrary positive integer. Let us replace $\lambda_{1}$ by $\lambda_{k}$ in the definition of $Q_{1}$ and let us denote by $Q_{k}$ the corresponding quantity. Assume that the conditions (13), (14) are fulfilled with $Q_{k}$ instead of $Q_{1}$. (If $I$ is an arbitrary interval of length $|I|>2 R$, then (13) and (14) are satisfied for a sufficiently large $k$, but not necessarily for $k=1$.) Then the estimate (15) is valid (with some constant $c^{\prime}$ depending on $k$ ) for every couple ( $u^{0}, u^{1}$ ) in $H_{0}^{1}(\Omega) \times L^{2}(\Omega)$ satisfying the orthogonality conditions $u^{0}, u^{1} \perp Z_{j}$ for every $j<k$. To see this it suffices to repeat the proof of the theorem and to use for $v:=u(t)$ instead of (12) the inequality

$$
\int_{\Omega}|\nabla v|^{2}+q|v|^{2} \mathrm{dx} \geq \lambda_{k} \int_{\Omega}|v|^{2} \mathrm{dx}
$$

which is valid for every $v \in H_{0}^{1}(\Omega)$ such that $v \perp Z_{j}$ for every $j<k$. (This estimate follows easily from the identity (1.10).) Let us recall from remark 1.2 that $u(t)$ satisfies these orthogonality relations.

Remark 3.4. - Theorem 3.1 implies the following uniqueness result : If the solution of (6)-(8) satisfies the condition $\partial_{\nu} u=0$ on $\Gamma_{+} \times I$ with $I$ satisfying (13) and (14), then in fact $u^{0}=u^{1}=0$ and hence $u \equiv 0$ on $\Omega \times \mathbb{R}$.

Remark 3.5. - Theorem 3.1 shows that the "observation" of $\partial_{\nu} u=0$ on $\Gamma_{+} \times I$ permits one to distinguish the initial data provided $I$ is sufficiently
long. (Take the difference of two solutions corresponding to different initial data and apply the preceding remark.)

Remark 3.6. - The optimality of condition (14) in theorem 3.1 and analogous questions are studied by microlocal analysis in Bardos, Lebeau and Rauch [1], [2]; it is related to the finite propagation speed for the wave equation. In some cases the optimality of this condition can be shown easily. Consider for example the special case where $\Omega$ is an open ball of radious $R$, centered at the origin, and assume for simplicity that $q \equiv 0$ and that $I=(-R+\varepsilon, R-\varepsilon)$ for some $\varepsilon>0$. Fix smooth initial data $\left(u^{0}, u^{1}\right) \in H_{0}^{1}(\Omega) \times L^{2}(\Omega)$ and (following Cazenave [1]) set

$$
F(t):=\int_{\varepsilon+|t|<|x|<R}\left(\left(u^{\prime}\right)^{2}+|\nabla u|^{2}\right)(t, x) \mathrm{dx}, \quad t \in I .
$$

Then $F \geq 0$ and $F$ has a global maximum at $t=0$. Indeed, for almost every $t>0$ we have

$$
\begin{gathered}
F^{\prime}(t)=2 \int_{\varepsilon+t<|x|<R}\left(u u^{\prime}+\nabla u \cdot \nabla u^{\prime}\right)(t, x) \mathrm{dx}-\int_{\varepsilon+t=|x|}\left(\left(u^{\prime}\right)^{2}+|\nabla u|^{2}\right)(t, x) \mathrm{dx} \\
=\int_{\varepsilon+t=|x|}\left(2 u^{\prime} \partial_{\nu} u-\left(u^{\prime}\right)^{2}-|\nabla u|^{2}\right)(t, x) \mathrm{dx} \leq 0
\end{gathered}
$$

because $\left|\partial_{\nu} u\right| \leq|\nabla u|$. For $t<0$ we obtain similarly that

$$
F^{\prime}(t)=\int_{\varepsilon-t=|x|}\left(2 u^{\prime} \partial_{\nu} u+\left(u^{\prime}\right)^{2}+|\nabla u|^{2}\right)(t, x) \mathrm{dx} \geq 0
$$

Now choose non-zero initial data $\left(u^{0}, u^{1}\right)$ supported by the ball $\{|x|<\varepsilon\}$. Then $F(0)=0$ and therefore $F \equiv 0$ in $I$. It follows that $\partial_{\nu} u \equiv 0$ on $\Gamma \times I$ and therefore the estimate (15) of theorem 3.1 does not hold.

### 3.2. The first Petrovsky system

Consider the homogeneous problem

$$
\begin{gather*}
u^{\prime \prime}+\Delta^{2} u=0 \quad \text { in } \quad \Omega \times \mathbb{R}  \tag{23}\\
u=\partial_{\nu} u=0 \quad \text { on } \quad \Gamma \times \mathbb{R}  \tag{24}\\
u(0)=u^{0} \quad \text { and } \quad u^{\prime}(0)=u^{1} \tag{25}
\end{gather*}
$$

already studied in sections 1.4 and 2.3. Let us recall the theorem 2.6 : for every interval $I$ there exists a constant $c>0$ such that for every $\left(u^{0}, u^{1}\right) \in H_{0}^{2}(\Omega) \times L^{2}(\Omega)$ the solution of (23)-(25) satisfies

$$
\begin{equation*}
\int_{I} \int_{\Gamma}|\Delta u|^{2} \mathrm{~d} \Gamma \mathrm{dt} \leq c E \tag{26}
\end{equation*}
$$

where the energy $E$ of the solution is defined by

$$
\begin{equation*}
E=\frac{1}{2} \int_{\Omega}|\Delta u|^{2}+\left|u^{\prime}\right|^{2} \mathrm{dx} . \tag{27}
\end{equation*}
$$

We are going to prove the following inverse inequality :
Theorem 3.7. - Let $\mu_{1}$ denote the first eigenvalue of the problem

$$
\Delta^{2} v=-\mu \Delta v, \quad v \in H_{0}^{2}(\Omega)
$$

and let I be an interval satisfying

$$
\begin{equation*}
|I|>2 R \mu_{1}^{-1 / 2} \quad \text { if } \quad n=1 \quad \text { and } \quad|I|>R \mu_{1}^{-1 / 2} \quad \text { if } \quad n \geq 2 \tag{28}
\end{equation*}
$$

Then there is a constant $c^{\prime}>0$ such that

$$
\begin{equation*}
\int_{I} \int_{\Gamma}|\Delta u|^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt} \geq c^{\prime} E, \forall\left(u^{0}, u^{1}\right) \in H_{0}^{2}(\Omega) \times L^{2}(\Omega) \tag{29}
\end{equation*}
$$

Remark 3.8. - To be more precise, $\mu_{1}$ is the smallest eigenvalue of the operator $A$ associated to the spaces $H=H_{0}^{1}(\Omega),\|v\|_{H}=\|\nabla v\|_{L^{2}(\Omega)}$ and $V=H_{0}^{2}(\Omega),\|v\|_{V}=\|\Delta v\|_{L^{2}(\Omega)}$ in the sense of section 1.1. In particular, we have

$$
\begin{equation*}
\|\nabla v\|_{L^{2}(\Omega)} \leq \frac{1}{\sqrt{\mu_{1}}}\|\Delta v\|_{L^{2}(\Omega)}, \forall v \in H_{0}^{2}(\Omega) \tag{30}
\end{equation*}
$$

(see (1.10)).
Teorem 3.7 was proved by Lions [3], [4] under a stronger condition on the length of $I$. His assumption was later weakened in Komornik [1].

Proof of theorem 3.7. - It suffices to prove (29) for $u^{0}, u^{1} \in Z$.
Write $I=[S, T]$ and apply lemma 2.7 with $h=m$. Using (5) we obtain that

$$
\begin{align*}
& \int_{S}^{T} \int_{\Gamma}|\Delta u|^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt}=\left[\int_{\Omega} 2 u^{\prime} m \cdot \nabla u \mathrm{dx}\right]_{S}^{T}  \tag{31}\\
& \quad+\int_{S}^{T} \int_{\Omega} n\left(u^{\prime}\right)^{2}+(4-n)|\Delta u|^{2} \mathrm{dx} \mathrm{dt}
\end{align*}
$$

If $n=1$, then the last integral is minorized by $2|I| E$, while using (30) and (27) we obtain for every $t \in \mathbb{R}$ that

$$
\left|\int_{\Omega} 2 u^{\prime} m \cdot \nabla u \mathrm{dx}\right| \leq \frac{2 R}{\sqrt{\mu_{1}}}\left\|u^{\prime}\right\|_{L^{2}(\Omega)}\|\Delta v\|_{L^{2}(\Omega)} \leq \frac{2 R}{\sqrt{\mu_{1}}} E .
$$

Using the conservation of energy hence we conclude that

$$
\int_{S}^{T} \int_{\Gamma}|\Delta u|^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt} \geq\left(2|I|-\frac{4 R}{\sqrt{\mu_{1}}}\right) E
$$

and (29) follows with $c^{\prime}=2|I|-4 R / \sqrt{\mu_{1}}$.
Henceforth we assume that $n \geq 2$. We multiply the equation (23) by $u$ and we integrate by parts :

$$
\begin{gathered}
0=\int_{S}^{T} \int_{\Omega} u\left(u^{\prime \prime}+\Delta^{2} u\right) \mathrm{dx} \mathrm{dt}=\left[\int_{\Omega} u u^{\prime} \mathrm{dx}\right]_{S}^{T} \\
+\int_{S}^{T} \int_{\Gamma} u\left(\partial_{\nu} \Delta u\right)-\left(\partial_{\nu} u\right)(\Delta u) \mathrm{d} \Gamma \mathrm{dt}+\int_{S}^{T} \int_{\Omega}-\left(u^{\prime}\right)^{2}+|\Delta u|^{2} \mathrm{dx} \mathrm{dt} \\
=\left[\int_{\Omega} u u^{\prime} \mathrm{dx}\right]_{S}^{T}+\int_{S}^{T} \int_{\Omega}-\left(u^{\prime}\right)^{2}+|\Delta u|^{2} \mathrm{dx} \mathrm{dt} .
\end{gathered}
$$

(We used the boundary conditions (24).) Multiplying this identity by $n-2$ and adding the result to (31) we obtain, putting

$$
M u:=2 m \cdot \nabla u+(n-2) u
$$

for brevity, that

$$
\begin{equation*}
\int_{S}^{T} \int_{\Gamma}|\Delta u|^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt}=4|I| E+\left[\int_{\Omega} u^{\prime} M u \mathrm{dx}\right]_{S}^{T} \tag{32}
\end{equation*}
$$

Lemma 3.9. - If $n \geq 2$, then we have

$$
\begin{equation*}
\left|\int_{\Omega} u^{\prime} M u \mathrm{dx}\right| \leq \frac{2 R}{\sqrt{\mu_{1}}} E, \forall t \in \mathbb{R} \tag{33}
\end{equation*}
$$

Proof. - We have (compare with lemma 3.2)

$$
\begin{aligned}
& \|M u\|_{L^{2}(\Omega)}^{2}-\|2 m \cdot \nabla u\|_{L^{2}(\Omega)}^{2}=\|2 m \cdot \nabla u+(n-2) u\|_{L^{2}(\Omega)}^{2}-\|2 m \cdot \nabla u\|_{L^{2}(\Omega)}^{2} \\
& \begin{array}{c}
=\int_{\Omega}|2 m \cdot \nabla u+(n-2) u|^{2}-|2 m \cdot \nabla u|^{2} \mathrm{dx}=\int_{\Omega}(n-2)^{2} u^{2}+4(n-2) u m \cdot \nabla u \mathrm{dx} \\
=\int_{\Omega}(n-2)^{2} u^{2}+2(n-2) m \cdot \nabla\left(u^{2}\right) \mathrm{dx} \\
=2(n-2) \int_{\Gamma}(m \cdot \nu) u^{2} \mathrm{~d} \Gamma+\int_{\Omega}(n-2)^{2} u^{2}-2(n-2)(\operatorname{div} m) u^{2} \mathrm{dx} \\
=2(n-2) \int_{\Gamma}(m \cdot \nu) u^{2} \mathrm{~d} \Gamma+\left(4-n^{2}\right) \int_{\Omega} u^{2} \mathrm{dx}=\left(4-n^{2}\right) \int_{\Omega} u^{2} \mathrm{dx} \leq 0 .
\end{array}
\end{aligned}
$$

Using (27) and (30) hence we deduce that

$$
\begin{aligned}
& \left|\int_{\Omega} u^{\prime} M u \mathrm{dx}\right| \leq\left\|u^{\prime}\right\|_{L^{2}(\Omega)}\|M u\|_{L^{2}(\Omega)} \leq\left\|u^{\prime}\right\|_{L^{2}(\Omega)}\|2 m \cdot \nabla u\|_{L^{2}(\Omega)} \\
& \leq 2 R\left\|u^{\prime}\right\|_{L^{2}(\Omega)}\|\nabla u\|_{L^{2}(\Omega)} \leq \frac{2 R}{\sqrt{\mu_{1}}}\left\|u^{\prime}\right\|_{L^{2}(\Omega)}\|\Delta u\|_{L^{2}(\Omega)} \leq \frac{2 R}{\sqrt{\mu_{1}}} E .
\end{aligned}
$$

We conclude from (32) and (33) that

$$
\begin{equation*}
\int_{S}^{T} \int_{\Gamma}|\Delta u|^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt} \geq\left(4|I|-\frac{4 R}{\sqrt{\mu_{1}}}\right) E \tag{34}
\end{equation*}
$$

and (29) follows with $c^{\prime}=4|I|-4 R / \sqrt{\mu_{1}}$.
Remark 3.10. - Theorem 3.7 implies the following uniqueness result : If the solution of (23)-(25) satisfies the condition $\Delta u=0$ on $\Gamma \times I$ for some interval I satisfying (28), then in fact $u^{0}=u^{1}=0$ and hence $u \equiv 0$ on $\Omega \times \mathbb{R}$.

Remark 3.11. - Theorem 3.7 shows that the "observation" of $\Delta u=0$ on $\Gamma \times I$ permits one to distinguish the initial data provided $I$ is sufficiently long.

In connection with the inequality (30) let us note the following result :
Lemma 3.12. - Let $k$ be an arbitrary positive integer. If $v \in Z$ and $v \perp Z_{j}$ for every $j<k$, then

$$
\begin{equation*}
\|\nabla v\|_{L^{2}(\Omega)} \leq \lambda_{k}^{-1 / 4}\|\Delta v\|_{L^{2}(\Omega)} \tag{35}
\end{equation*}
$$

(Here the eigenvalues $\lambda_{j}$ and the eigenspaces $Z_{j}$ are those associated with $H=L^{2}(\Omega)$ and $\left.V=H_{0}^{2}(\Omega).\right)$

Proof. - Integrating by parts we obtain that

$$
\|\nabla v\|_{L^{2}(\Omega)}^{2}=-\int_{\Omega} v \Delta v \mathrm{dx} \leq\|v\|_{L^{2}(\Omega)}\|\Delta v\|_{L^{2}(\Omega)}
$$

The hypothesis on $v$ implies that (we apply (1.10) with $H$ and $V$ as above)

$$
\|v\|_{L^{2}(\Omega)} \leq \lambda_{k}^{-1 / 2}\|v\|_{H_{0}^{2}(\Omega)}\left(=\lambda_{k}^{-1 / 2}\|\Delta v\|_{L^{2}(\Omega)}\right)
$$

Hence

$$
\|\nabla v\|_{L^{2}(\Omega)}^{2} \leq \lambda_{k}^{-1 / 2}\|\Delta v\|_{L^{2}(\Omega)}^{2}
$$

and (35) follows.

Remark 3.13. - Let $k$ be a positive integer and assume that condition (28) is fulfilled with $\lambda_{k}^{-1 / 4}$ instead of $\mu_{1}^{-1 / 2}$. (Observe that every interval, of arbitrarily small length, satisfies this condition for a sufficiently large $k$.) Then the estimate (29) holds (with some constant $c^{\prime}$ depending on $k$ ) for every couple $\left(u^{0}, u^{1}\right)$ in $H_{0}^{2}(\Omega) \times L^{2}(\Omega)$ satisfying the orthogonality conditions $u^{0}, u^{1} \perp Z_{j}$ for every $j<k$. Indeed, it is sufficient to repeat the proof of theorem 3.7 and to apply for $v:=u(t)$ the inequality (35) instead of (30). (We recall from remark 1.2 that $u(t)$ satisfies these orthogonality relations.)

### 3.3. The second Petrovsky system

We consider here the problem

$$
\begin{gather*}
u^{\prime \prime}+\Delta^{2} u=0 \quad \text { in } \quad \Omega \times \mathbb{R}  \tag{36}\\
u=\Delta u=0 \quad \text { on } \quad \Gamma \times \mathbb{R}  \tag{37}\\
u(0)=u^{0} \quad \text { and } \quad u^{\prime}(0)=u^{1} \tag{38}
\end{gather*}
$$

studied earlier in sections 1.5 and 2.4. We recall from theorem 2.13 the following direct inequality :

$$
\begin{equation*}
\int_{I} \int_{\Gamma}\left|\partial_{\nu} u^{\prime}\right|^{2}+\left|\partial_{\nu} \Delta u\right|^{2} \mathrm{~d} \Gamma \mathrm{dt} \leq c E_{1 / 4} \tag{39}
\end{equation*}
$$

it is valid for every interval $I$ and for every solution of (36)-(38) corresponding to the initial data

$$
u^{0} \in D_{3 / 4}=\left\{v \in H^{3}(\Omega): v=\Delta v=0 \quad \text { on } \quad \Gamma\right\}
$$

and

$$
u^{1} \in D_{1 / 4}=H_{0}^{1}(\Omega)
$$

(The constant $c$ does not depend on $I$.)
Let us recall that the "increased" energy $E_{1 / 4}$ of the solution is defined by

$$
\begin{equation*}
E_{1 / 4}:=\frac{1}{2}\|u\|_{3 / 4}^{2}+\frac{1}{2}\left\|u^{\prime}\right\|_{1 / 4}^{2}\left(=\frac{1}{2} \int_{\Omega}|\nabla \Delta u|^{2}+\left|\nabla u^{\prime}\right|^{2} \mathrm{dx}\right) . \tag{40}
\end{equation*}
$$

In order to formulate an inverse result let us fix a point $x^{0} \in \mathbb{R}^{n}$ arbitrarily.
Theorem 3.14. - Let $\lambda_{1}$ denote the first eigenvalue of $-\Delta$ in $H_{0}^{1}(\Omega)$ and let $I$ be an interval of length

$$
\begin{equation*}
|I|>\frac{R}{\sqrt{\lambda_{1}}} \tag{41}
\end{equation*}
$$

Then there exists a constant $c^{\prime}>0$ such that

$$
\begin{equation*}
\int_{I} \int_{\Gamma}\left|\partial_{\nu} u^{\prime}\right|^{2}+\left|\partial_{\nu} \Delta u\right|^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt} \geq c^{\prime} E_{1 / 4}, \forall\left(u^{0}, u^{1}\right) \in D_{3 / 4} \times D_{1 / 4} \tag{42}
\end{equation*}
$$

Remark 3.15. - Here $\lambda_{1}$ is the smallest eigenvalue of the operator $A$ associated to $H:=L^{2}(\Omega)$ and $V:=H_{0}^{1}(\Omega)$; we have

$$
\begin{equation*}
\|v\|_{L^{2}(\Omega)} \leq \frac{1}{\sqrt{\lambda_{1}}}\|\nabla v\|_{L^{2}(\Omega)}, \quad \forall v \in H_{0}^{1}(\Omega) \tag{43}
\end{equation*}
$$

Theorem 3.14 was first proved, under a stronger condition on $|I|$, by Lions [3], [4]. His condition was later weakened in Komornik [1].
Proof of theorem 3.14. - It is sufficient to consider $u^{0}, u^{1} \in Z$. Let us apply lemma 2.12 with $h=m$ where $I=[S, T]$. We obtain

$$
\begin{gather*}
\int_{S}^{T} \int_{\Gamma}\left|\partial_{\nu} u^{\prime}\right|^{2}+\left|\partial_{\nu} \Delta u\right|^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt}=-\left[\int_{\Omega} 2 u^{\prime} m \cdot \nabla \Delta u \mathrm{dx}\right]_{S}^{T}  \tag{44}\\
\quad+\int_{S}^{T} \int_{\Omega}(2+n)\left|\nabla u^{\prime}\right|^{2}+(2-n)|\nabla \Delta u|^{2} \mathrm{dx} \mathrm{dt}
\end{gather*}
$$

On the other hand, multiplying the equation (36) by $\Delta u$ we obtain that

$$
\begin{gathered}
0=\int_{S}^{T} \int_{\Omega} \Delta u\left(u^{\prime \prime}+\Delta^{2} u\right) \mathrm{dx} \mathrm{dt} \\
=\left[\int_{\Omega} u^{\prime} \Delta u \mathrm{dx}\right]_{S}^{T}+\int_{S}^{T} \int_{\Gamma}(\Delta u)\left(\partial_{\nu} \Delta u\right) \mathrm{d} \Gamma \mathrm{dt} \\
\quad+\int_{S}^{T} \int_{\Omega}-u^{\prime} \Delta u^{\prime}-|\nabla \Delta u|^{2} \mathrm{dx} \mathrm{dt} \\
=\left[\int_{\Omega} u^{\prime} \Delta u \mathrm{dx}\right]_{S}^{T}+\int_{S}^{T} \int_{\Omega}\left|\nabla u^{\prime}\right|^{2}-|\nabla \Delta u|^{2} \mathrm{dx} \mathrm{dt} \\
=\left[\int_{\Omega}^{T} \int_{\Gamma}(\Delta u)\left(\partial_{\nu} \Delta u\right)-u^{\prime}\left(\partial_{\nu} u^{\prime}\right) \mathrm{d} \Gamma \mathrm{dt}\right. \\
\quad+]_{S}^{T}+\int_{S}^{T} \int_{\Omega}\left|\nabla u^{\prime}\right|^{2}-|\nabla \Delta u|^{2} \mathrm{dx} \mathrm{dt}
\end{gathered}
$$

because $u^{\prime}=\Delta u=0$ on $\Gamma \times \mathbb{R}$ by (37). Using the notation

$$
M u:=2 m \cdot \nabla \Delta u+n \Delta u,
$$

combining with (44) and using (40) we find that

$$
\begin{equation*}
\int_{S}^{T} \int_{\Gamma}\left|\partial_{\nu} u^{\prime}\right|^{2}+\left|\partial_{\nu} \Delta u\right|^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt}=4|I| E_{1 / 4}-\left[\int_{\Omega} u^{\prime} M u \mathrm{dx}\right]_{S}^{T} \tag{45}
\end{equation*}
$$

Lemma 3.16. - We have

$$
\begin{equation*}
\left|\int_{\Omega} u^{\prime} M u \mathrm{dx}\right| \leq \frac{2 R}{\sqrt{\lambda_{1}}} E_{1 / 4}, \quad \forall t \in \mathbb{R} \tag{46}
\end{equation*}
$$

Proof. - We proceed as in lemmas 3.2 and 3.9. First, we have

$$
\begin{equation*}
\|M u\|_{L^{2}(\Omega)} \leq\|2 m \cdot \nabla \Delta u\|_{L^{2}(\Omega)} \tag{47}
\end{equation*}
$$

because

$$
\begin{gathered}
\|M u\|_{L^{2}(\Omega)}^{2}-\|2 m \cdot \nabla \Delta u\|_{L^{2}(\Omega)}^{2}=\|2 m \cdot \nabla \Delta u+n \Delta u\|_{L^{2}(\Omega)}^{2}-\|2 m \cdot \nabla \Delta u\|_{L^{2}(\Omega)}^{2} \\
=\int_{\Omega}|2 m \cdot \nabla \Delta u+n \Delta u|^{2}-|2 m \cdot \nabla \Delta u|^{2} \mathrm{dx} \\
=\int_{\Omega} n^{2}(\Delta u)^{2}+4 n(\Delta u) m \cdot \nabla \Delta u \mathrm{dx}=\int_{\Omega} n^{2}(\Delta u)^{2}+2 n m \cdot \nabla(\Delta u)^{2} \mathrm{dx} \\
=2 n \int_{\Gamma}(m \cdot \nu)(\Delta u)^{2} \mathrm{~d} \Gamma+\int_{\Omega}\left(n^{2}-2 n \operatorname{div} m\right)(\Delta u)^{2} \mathrm{dx} \\
=-n^{2} \int_{\Omega}(\Delta u)^{2} \mathrm{dx} \leq 0
\end{gathered}
$$

Consequently, using (40) and (43) we obtain that

$$
\begin{gathered}
\left|\int_{\Omega} u^{\prime} M u \mathrm{dx}\right| \leq\left\|u^{\prime}\right\|_{L^{2}(\Omega)}\|M u\|_{L^{2}(\Omega)} \\
\leq\left\|u^{\prime}\right\|_{L^{2}(\Omega)}\|2 m \cdot \nabla \Delta u\|_{L^{2}(\Omega)} \leq 2 R\left\|u^{\prime}\right\|_{L^{2}(\Omega)}\|\nabla \Delta u\|_{L^{2}(\Omega)} \\
\leq\left(2 R / \sqrt{\lambda_{1}}\right)\left\|\nabla u^{\prime}\right\|_{L^{2}(\Omega)}\|\nabla \Delta u\|_{L^{2}(\Omega)} \leq\left(2 R / \sqrt{\lambda_{1}}\right) E_{1 / 4} .
\end{gathered}
$$

We deduce from (45) and (46) the inequality

$$
\int_{S}^{T} \int_{\Gamma}\left|\partial_{\nu} u^{\prime}\right|^{2}+\left|\partial_{\nu} \Delta u\right|^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt} \geq 4|I| E_{1 / 4}-\frac{4 R}{\sqrt{\lambda_{1}}} E_{1 / 4}
$$

and the theorem follows with $c^{\prime}=4|I|-4 R / \sqrt{\lambda_{1}}$.
Remark 3.17. - Let $k$ be a positive integer and assume that condition (41) is satisfied with $\lambda_{k}$ instead of $\lambda_{1}$. (Observe that every interval, arbitrarily small, satisfies this condition for a sufficiently large $k$.) Then the estimate
(42) holds (with a constant $c^{\prime}$ depending on $k$ ) for every couple ( $u^{0}, u^{1}$ ) satisfying the orthogonality conditions $u^{0}, u^{1} \perp Z_{j}$ for all $j<k$. Indeed, it is sufficient to repeat the proof of the theorem and to apply for $v:=u(t)$ the inequality

$$
\begin{equation*}
\|v\|_{L^{2}(\Omega)} \leq \frac{1}{\sqrt{\lambda_{k}}}\|\nabla v\|_{L^{2}(\Omega)}, \quad \forall v \in H_{0}^{1}(\Omega), v \perp Z_{j}, \forall j<k \tag{48}
\end{equation*}
$$

instead of (43). (Compare with remarks 3.3 and 3.13.)
Remark 3.18. - As before, theorem 3.15 yields some uniqueness and observability results.

### 3.4. The wave equation. Mixed boundary conditions

Here we consider the problem of section 1.3 with $\Gamma_{1} \neq \emptyset:$

$$
\begin{gather*}
u^{\prime \prime}-\Delta u+q u=0 \quad \text { in } \quad \Omega \times \mathbb{R}  \tag{49}\\
u=0 \quad \text { on } \quad \Gamma_{0} \times \mathbb{R}  \tag{50}\\
\partial_{\nu} u+a u=0 \quad \text { on } \quad \Gamma_{1} \times \mathbb{R}  \tag{51}\\
u(0)=u^{0} \quad \text { and } \quad u^{\prime}(0)=u^{1} \tag{52}
\end{gather*}
$$

In order to avoid some difficulites we shall assume that

$$
\begin{equation*}
m \cdot \nu \leq 0 \quad \text { on } \quad \Gamma_{0}, \quad m \cdot \nu \geq 0 \quad \text { on } \quad \Gamma_{1} \tag{53}
\end{equation*}
$$

and that $a$ has the form

$$
\begin{equation*}
a=(m \cdot \nu) b, \quad b \in C^{1}\left(\Gamma_{1}\right), \quad b \geq 0 \text { on } \Gamma_{1} . \tag{54}
\end{equation*}
$$

We recall that for every couple $\left(u^{0}, u^{1}\right) \in V \times H$ the problem (49)-(52) has a unique solution whose energy

$$
\begin{equation*}
E:=\frac{1}{2} \int_{\Omega}\left|u^{\prime}\right|^{2}+|\nabla u|^{2}+q|u|^{2} \mathrm{dx}+\frac{1}{2} \int_{\Gamma_{1}} a|u|^{2} \mathrm{~d} \Gamma \tag{55}
\end{equation*}
$$

is conserved.
Remark 3.19. - The conditions (53) and $\overline{\Gamma_{0}} \cap \overline{\Gamma_{1}}=\emptyset$ (see (1.25)) together impose a very restrictive geometrical property on $\Omega$. It is satisfied if $\Omega$ is star-shaped with respect to $x^{0}$ or if $\Omega$ has the form $\Omega_{1} \backslash \overline{\Omega_{0}}$ with two open sets $\Omega_{1}$ and $\Omega_{0}$, both star-shaped with respect to $x^{0}$ and such that $\overline{\Omega_{0}} \subset \Omega_{1}$. In
this case we may choose for $\Gamma_{0}, \Gamma_{1}$ the boundaries of $\Omega_{0}$ and $\Omega_{1}$, respectively. We shall return to this point later in remark 3.25.

Let us denote by $R_{1}$ the smallest positive constant such that

$$
\begin{gather*}
4 R^{2} \int_{\Omega}|\nabla v|^{2} \mathrm{dx}+(2 n-2) \int_{\Gamma_{1}}(m \cdot \nu)|v|^{2} \mathrm{~d} \Gamma  \tag{56}\\
\leq 4 R_{1}^{2}\left(\int_{\Omega}|\nabla v|^{2}+q|v|^{2} \mathrm{dx}+\int_{\Gamma_{1}} a|v|^{2} \mathrm{~d} \Gamma\right), \quad \forall v \in V .
\end{gather*}
$$

Remark 3.20. - If

$$
\begin{equation*}
b \geq \frac{n-1}{2 R^{2}} \quad \text { sur } \quad \Gamma_{1} \tag{57}
\end{equation*}
$$

then we may choose $R_{1}:=R$.
We define $Q:=\sup _{\Omega} q$ and

$$
Q_{k}:= \begin{cases}2 R Q / \sqrt{\lambda_{k}}, & \text { si } n \geq 2  \tag{58}\\ 2 R Q / \sqrt{\lambda_{k}}+Q / \lambda_{k}, & \text { si } n=1\end{cases}
$$

for $k=1,2, \ldots$, as in section 3.1. Let us recall that $\lambda_{k}$ is here the biggest constant such that

$$
\begin{equation*}
\int_{\Omega}|\nabla v|^{2}+q|v|^{2} \mathrm{dx}+\int_{\Gamma_{1}} a|v|^{2} \mathrm{~d} \Gamma \geq \lambda_{k} \int_{\Omega}|v|^{2} \mathrm{dx} \tag{59}
\end{equation*}
$$

for every $v \in V$ satisfying the orthogonality conditions $v \perp Z_{j}, \forall j<k$.
We are going to prove the
Proposition 3.21. - Assume that

$$
\begin{equation*}
Q_{1}<1 \tag{60}
\end{equation*}
$$

and let $I$ be an interval of length

$$
\begin{equation*}
|I|>2 R_{1} /\left(1-Q_{1}\right) \tag{61}
\end{equation*}
$$

Then there exists a constant $c^{\prime}>0$ such that

$$
\begin{equation*}
\int_{I} \int_{\Gamma_{1}}|u|^{2}+\left|u^{\prime}\right|^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt} \geq c^{\prime} E, \quad \forall\left(u^{0}, u^{1}\right) \in D_{1} \times D_{1 / 2} \tag{62}
\end{equation*}
$$

Let us recall that $D_{1 / 2}=V \subset H^{1}(\Omega)$ and that $D_{1} \subset H^{2}(\Omega) \cap V$ by remark 1.4.

Proof. - It is sufficient to consider the case $\left(u^{0}, u^{1}\right) \in Z \times Z$. Then the following computations are justified by remark 1.4.

We recall from the proof of theorem 3.1 the identity (18) :

$$
\begin{gathered}
\int_{S}^{T} \int_{\Gamma}\left(\partial_{\nu} u\right) M u+(m \cdot \nu)\left(\left(u^{\prime}\right)^{2}-|\nabla u|^{2}\right) \mathrm{d} \Gamma \mathrm{dt} \\
=\left[\int_{\Omega} u^{\prime} M u \mathrm{dx}\right]_{S}^{T}+\int_{S}^{T} \int_{\Omega}\left(u^{\prime}\right)^{2}+|\nabla u|^{2}+(n-1) q u^{2}+2 q u m \cdot \nabla u \mathrm{dx} \mathrm{dt}
\end{gathered}
$$

where

$$
M u=2 m \cdot \nabla u+(n-1) u ;
$$

using (55) hence we deduce that

$$
\begin{aligned}
& \int_{S}^{T} \int_{\Gamma}\left(\partial_{\nu} u\right) M u+(m \cdot \nu)\left(\left(u^{\prime}\right)^{2}-|\nabla u|^{2}\right) \mathrm{d} \Gamma \mathrm{dt}+\int_{I} \int_{\Gamma_{1}} a u^{2} \mathrm{~d} \Gamma \mathrm{dt} \\
& =\left[\int_{\Omega} u^{\prime} M u \mathrm{dx}\right]_{S}^{T}+2|I| E+\int_{S}^{T} \int_{\Omega}(n-2) q u^{2}+2 q u m \cdot \nabla u \mathrm{dx} \mathrm{dt}
\end{aligned}
$$

Using (50), (51) and (54) we obtain

$$
\begin{gather*}
\int_{I} \int_{\Gamma_{1}}\left(u^{\prime}\right)^{2}-|\nabla u|^{2}-2 b u(m \cdot \nabla u)+(2-n) b u^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt} \\
+\int_{I} \int_{\Gamma_{0}}\left(\partial_{\nu} u\right)^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt}  \tag{63}\\
=\left[\int_{\Omega} u^{\prime} M u \mathrm{dx}\right]_{S}^{T}+2|I| E+\int_{S}^{T} \int_{\Omega}(n-2) q u^{2}+2 q u m \cdot \nabla u \mathrm{dx} \mathrm{dt} .
\end{gather*}
$$

By (53) the second integral is $\leq 0$. Using the inequality

$$
-2 b u(m \cdot \nabla u) \leq 2 R|b u||\nabla u| \leq|\nabla u|^{2}+R^{2} b^{2} u^{2}
$$

we may majorize the first integral (and therefore the left-hand side of (63)) by

$$
\int_{I} \int_{\Gamma_{1}}\left(u^{\prime}\right)^{2}+\left(R^{2} b+2-n\right) b u^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt}
$$

Using (59) and (55) the last integral in (63) is minorized (as in section 3.1) by $-2|I| Q_{1} E$. Thus we deduce from (63) the estimate

$$
\begin{align*}
& \int_{I} \int_{\Gamma_{1}}\left(u^{\prime}\right)^{2}+\left(R^{2} b+2-n\right) b u^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt} \\
& \quad \geq 2|I|\left(1-Q_{1}\right) E+\left[\int_{\Omega} u^{\prime} M u \mathrm{dx}\right]_{S}^{T} \tag{64}
\end{align*}
$$

Lemma 3.22. - We have the estimate

$$
\begin{equation*}
\left|\int_{\Omega} u^{\prime} M u d x\right| \leq 2 R_{1} E, \quad \forall t \in \mathbb{R} . \tag{65}
\end{equation*}
$$

Proof. - We have

$$
\begin{gathered}
\int_{\Omega}|M u|^{2} \mathrm{dx} \\
=\int_{\Omega}|2 m \cdot \nabla u+(n-1) u|^{2} \mathrm{dx} \\
=\int_{\Omega}|2 m \cdot \nabla u|^{2}+(n-1)^{2} u^{2}+4(n-1) u m \cdot \nabla u \mathrm{dx} \\
=\int_{\Omega}|2 m \cdot \nabla u|^{2}+(n-1)^{2} u^{2}+(2 n-2) m \cdot \nabla\left(u^{2}\right) \mathrm{dx} \\
=\int_{\Omega}|2 m \cdot \nabla u|^{2}+(n-1)^{2} u^{2}-n(2 n-2) u^{2} \mathrm{dx}+(2 n-2) \int_{\Gamma}(m \cdot \nu) u^{2} \mathrm{~d} \Gamma \\
=\int_{\Omega}|2 m \cdot \nabla u|^{2}+\left(1-n^{2}\right) u^{2} \mathrm{dx}+(2 n-2) \int_{\Gamma_{1}}(m \cdot \nu) u^{2} \mathrm{~d} \Gamma \\
\leq 4 R^{2} \int_{\Omega}|\nabla u|^{2} \mathrm{dx}+(2 n-2) \int_{\Gamma_{1}}(m \cdot \nu) u^{2} \mathrm{~d} \Gamma \\
\leq 4 R_{1}^{2}\left(\int_{\Omega}|\nabla u|^{2}+q v^{2} \mathrm{dx}+\int_{\Gamma_{1}}(m \cdot \nu) b u^{2} \mathrm{~d} \Gamma\right)
\end{gathered}
$$

whence

$$
\begin{gathered}
\int_{\Omega}\left|u^{\prime} M u\right| \mathrm{dx} \\
\leq R_{1} \int_{\Omega}\left(u^{\prime}\right)^{2} \mathrm{dx}+\frac{4 R_{1}^{2}}{4 R_{1}}\left(\int_{\Omega}|\nabla u|^{2}+q u^{2} \mathrm{dx}+\int_{\Gamma_{1}}(m \cdot \nu) b u^{2} \mathrm{~d} \Gamma\right)=2 R_{1} E
\end{gathered}
$$

Applying lemma 3.22 with $t=S$ and $t=T$ we conclude from (64) that

$$
\begin{equation*}
\int_{I} \int_{\Gamma_{1}}\left(u^{\prime}\right)^{2}+\left(R^{2} b+2-n\right) b u^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt} \geq 2|I|\left(1-Q_{1}\right) E-4 R_{1} E . \tag{66}
\end{equation*}
$$

Let $c_{1} \geq 1$ be a majorant of $R^{2} b^{2}+(2-n) b$ on $\Gamma_{1}$, then (62) follows from (66) with

$$
c^{\prime}:=\left(2|I|\left(1-Q_{1}\right)-4 R_{1}\right) / c_{1}
$$

Remark 3.23. - If

$$
\begin{equation*}
R^{2} b^{2}+(2-n) b \leq 0 \quad \text { on } \quad \Gamma_{1}, \tag{67}
\end{equation*}
$$

then we may deduce from (66) the stronger estimate

$$
\begin{equation*}
\int_{I} \int_{\Gamma_{1}}\left|u^{\prime}\right|^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt} \geq c^{\prime} E, \quad \forall\left(u^{0}, u^{1}\right) \in D_{1} \times D_{1 / 2} \tag{68}
\end{equation*}
$$

with

$$
c^{\prime}:=2|I|\left(1-Q_{1}\right)-4 R_{1} .
$$

Remark 3.24. - Let us denote by $R_{k}(k=1,2 \ldots)$ the smallest positive constant such that

$$
\begin{aligned}
& 4 R^{2} \int_{\Omega}|\nabla v|^{2} \mathrm{dx}+(2 n-2) \int_{\Gamma_{1}}(m \cdot \nu)|v|^{2} \mathrm{~d} \Gamma \\
\leq & 4 R_{k}^{2}\left(\int_{\Omega}|\nabla v|^{2}+q|v|^{2} \mathrm{dx}+\int_{\Gamma_{1}} b(m \cdot \nu)|v|^{2} \mathrm{~d} \Gamma\right)
\end{aligned}
$$

for every $v \in V$ satisfying the orthogonality conditions $v \perp Z_{j}, \forall j<k$. Let $k$ be a positive integer and assume that the conditions (60) and (61) are satisfied with $Q_{k}, R_{k}$ instead of $Q_{1}, R_{1}$. (One can readily verify that for every interval $I$ of length $>2 R$ the conditions (60) and (61) are satified with a sufficiently large $k$.) Then (62) is valid (with a constant $c^{\prime}$ depending on $k$ ) for every couple ( $u^{0}, u^{1}$ ) satisfying the orthogonality conditions $u^{0}, u^{1} \perp Z_{j}$, $\forall j<k$. Indeed, it suffices to repeat the above proof and to observe that (by remark 1.2) we may replace $Q_{1}, R_{1}$ in the estimates by $Q_{k}, R_{k}$.

Remark 3.25. - Grisvard [1] proved by a delicate analysis of singularities that in dimension $n \leq 3$ the inequality (62) remains valid without the hypothesis $\overline{\Gamma_{0}} \cap \overline{\Gamma_{1}}=\emptyset$. The difficulty comes from the fact that the identity of lemma 2.3 and the identity (18) do not hold any more even if $\left(u^{0}, u^{1}\right) \in Z \times Z$. (The solutions are not sufficiently smooth to justify the integrations by parts.) We only have inequalities instead of these identities; fortunately, these are still sufficient for the proof of the desired estimates.

Theorem 3.26. - Assume that

$$
\begin{gather*}
Q_{1}<1  \tag{69}\\
R^{2} b^{2}+(2-n) b \leq 0 \quad \text { on } \quad \Gamma_{1} \tag{70}
\end{gather*}
$$

and let I be an interval of length

$$
\begin{equation*}
|I|>2 R_{1} /\left(1-Q_{1}\right) \tag{71}
\end{equation*}
$$

Then there is a constant $c^{\prime}>0$ such that

$$
\int_{I} \int_{\Gamma_{1}}|u|^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt} \geq c^{\prime} E_{-1 / 2}, \quad \forall\left(u^{0}, u^{1}\right) \in V \times H
$$

We recall that

$$
E_{-1 / 2}=\frac{1}{2}\left\|u^{0}\right\|_{H}^{2}+\frac{1}{2}\left\|u^{1}\right\|_{V^{\prime}}^{2} .
$$

The case $q \equiv 0$ and $b \equiv 0$ of this result is a weakened version of a theorem of Lions [4, p. 200]. We shall prove a stronger result later (see theorem 6.15).
Proof. - It suffices to consider $u^{0}, u^{1} \in Z$. Set $z^{0}:=-A^{-1} u^{1}, z^{1}:=u^{0}$, and apply the estimate (68) to the solution $z$ of (49)-(52) corresponding to the initial data $\left(z^{0}, z^{1}\right)$ instead of $\left(u^{0}, u^{1}\right)$. We obtain that

$$
\begin{gathered}
2 \int_{I} \int_{\Gamma_{1}}\left|z^{\prime}\right|^{2} \mathrm{~d} \Gamma \mathrm{dt} \geq c^{\prime}\left(\left\|z^{0}\right\|_{V}^{2}+\left\|z^{1}\right\|_{H}^{2}\right) \\
=c^{\prime}\left(\left\|A^{-1} u^{1}\right\|_{V}^{2}+\left\|u^{0}\right\|_{H}^{2}\right)=c^{\prime}\left(\left\|u^{0}\right\|_{H}^{2}+\left\|u^{1}\right\|_{V^{\prime}}^{2}\right)=c^{\prime} E_{-1 / 2}(u)
\end{gathered}
$$

it remains to verify that $z^{\prime} \equiv u$ or, by the uniqueness of the solutions, that $z^{\prime}$ satisfies (49)-(52). First, (49)-(51) are easily obtained by differentiating with respect to $t$ the analogous equations for $z$; the differentiation is permitted because $z \in C^{\infty}\left(\mathbb{R} ; H^{2}(\Omega)\right)$. Finally, (52) may be verified directly :

$$
z^{\prime}(0)=z^{1}=u^{0} \quad \text { and } \quad\left(z^{\prime}\right)^{\prime}(0)=(\Delta z-q z)(0)=-A z^{0}=u^{1}
$$

## 4. Exact controllability. Hilbert Uniqueness Method

We present here the Hilbert Uniqueness Method (HUM), introduced by Lions [3], [4], [5]. As in the preceding chapter, for any fixed $x^{0} \in \mathbb{R}^{n}$ we shall use the notation

$$
\begin{gather*}
m(x):=x-x^{0}, \quad x \in \mathbb{R}^{n},  \tag{1}\\
R=R\left(x^{0}\right):=\sup \left\{\left|x-x^{0}\right|: x \in \Omega\right\},  \tag{2}\\
\mathrm{d} \Gamma_{m}:=(m \cdot \nu) \mathrm{d} \Gamma,  \tag{3}\\
\Gamma_{+}:=\{x \in \Gamma: m(x) \cdot \nu(x)>0\}, \tag{4}
\end{gather*}
$$

and we set

$$
\begin{equation*}
\Gamma_{-}:=\{x \in \Gamma: m(x) \cdot \nu(x) \leq 0\} . \tag{5}
\end{equation*}
$$

As usual, we consider the real case only ; the complex case then follows easily.

### 4.1. The wave equation. Dirichlet control

Fix $T>0$ and consider the problem

$$
\begin{gather*}
y^{\prime \prime}-\Delta y+q y=0 \quad \text { in } \quad \Omega \times(0, T),  \tag{6}\\
y=v \quad \text { on } \quad \Gamma \times(0, T),  \tag{7}\\
y(0)=y^{0} \quad \text { and } \quad y^{\prime}(0)=y^{1} . \tag{8}
\end{gather*}
$$

It follows from theorem 2.5 that for any given $y^{0} \in L^{2}(\Omega), y^{1} \in H^{-1}(\Omega)$ and $v \in L^{2}\left(0, T ; L^{2}(\Gamma)\right)$ the problem (6)-(8) has a unique solution $\left(y, y^{\prime}\right) \in C\left([0, T] ; L^{2}(\Omega) \times H^{-1}(\Omega)\right)$. (Observe that here $\Gamma_{0}=\Gamma$ whence $\left.V=H^{-1}(\Omega)\right)$.

Definition. - The problem (6)-(8) is exactly controllable if for any given $\left(y^{0}, y^{1}\right),\left(y_{T}^{0}, y_{T}^{1}\right) \in L^{2}(\Omega) \times H^{-1}(\Omega)$ there exists $v \in L^{2}\left(0, T ; L^{2}(\Gamma)\right)$ such that the solution of (6)-(8) satisfies

$$
\begin{equation*}
y(T)=y_{T}^{0} \quad \text { and } \quad y^{\prime}(T)=y_{T}^{1} . \tag{9}
\end{equation*}
$$

As in section 3.1, fix $x^{0} \in \mathbb{R}^{n}$ arbitrarily and set $Q:=\sup _{\Omega} q$,

$$
Q_{1}:= \begin{cases}2 R Q / \sqrt{\lambda_{1}}, & \text { if } n \geq 2 \\ 2 R Q / \sqrt{\lambda_{1}}+Q / \lambda_{1}, & \text { if } n=1\end{cases}
$$

Theorem 4.1. - Assume that

$$
\begin{equation*}
Q_{1}<1 \tag{10}
\end{equation*}
$$

and let

$$
\begin{equation*}
T>2 R /\left(1-Q_{1}\right) \tag{11}
\end{equation*}
$$

Then for any given $\left(y^{0}, y^{1}\right),\left(y_{T}^{0}, y_{T}^{1}\right) \in L^{2}(\Omega) \times H^{-1}(\Omega)$ there exists $v \in L^{2}\left(0, T ; L^{2}(\Gamma)\right)$ such that

$$
\begin{equation*}
v=0 \quad \text { a.e. on } \quad \Gamma_{-} \times(0, T) \tag{12}
\end{equation*}
$$

and the solution of (6)-(8) satisfies

$$
\begin{equation*}
y(T)=y_{T}^{0} \quad \text { and } \quad y^{\prime}(T)=y_{T}^{1} \tag{13}
\end{equation*}
$$

This result shows that if $\Omega$ is contained in a ball of center $x^{0}$ and diameter $<T$, then the problem (6)-(8) is exactly controllable, even if we act on $\Gamma_{+}$ only.

Theorem 4.1 is due (for $q \equiv 0$ ) to Lions [4], [5].
Remark 4.2. - As for the choice of $T$ and $\Gamma_{+}$very precise conditions were given in Bardos, Lebeau and Rauch [1], using microlocal analysis. See also Cazenave [1], Graham and Russell [1], Joó [1] and Komornik [11] for estimates of $T$.

Let us consider the solution of the problem

$$
\begin{gathered}
y_{1}^{\prime \prime}-\Delta y_{1}+q y_{1}=0 \quad \text { in } \quad \Omega \times(0, T) \\
y_{1}=0 \quad \text { on } \quad \Gamma \times(0, T) \\
y_{1}(T)=y_{T}^{0} \quad \text { and } \quad y_{1}^{\prime}(T)=y_{T}^{1}
\end{gathered}
$$

and assume that there exists a unique function $v \in L^{2}\left(0, T ; L^{2}(\Gamma)\right)$ satisfying (12) and such that the solution of the problem

$$
\begin{gathered}
y_{2}^{\prime \prime}-\Delta y_{2}+q y_{2}=0 \quad \text { in } \quad \Omega \times(0, T) \\
y_{2}=v \quad \text { on } \quad \Gamma \times(0, T) \\
y_{2}(0)=y^{0}-y_{1}(0) \quad \text { and } \quad y_{2}^{\prime}(0)=y^{1}-y_{1}^{\prime}(0)
\end{gathered}
$$

satisfies $y_{2}(T)=y_{2}^{\prime}(T)=0$. Then $y:=y_{1}+y_{2}$ is a/the solution of (6)-(8) and it satisfies (13). In view of this remark it is sufficient to prove theorem 4.1 in the special case where $y_{T}^{0}=y_{T}^{1}=0$.

Henceforth we shall assume that $y_{T}^{0}=y_{T}^{1}=0$.
The first idea of HUM is to seek a control $v$ in the special form $v=\partial_{\nu} u$ where $u$ is the solution of the homogeneous problem

$$
\begin{gather*}
u^{\prime \prime}-\Delta u+q u=0 \quad \text { in } \quad \Omega \times(0, T)  \tag{14}\\
u=0 \quad \text { on } \quad \Gamma \times(0, T)  \tag{15}\\
u(0)=u^{0} \quad \text { and } \quad u^{\prime}(0)=u^{1} \tag{16}
\end{gather*}
$$

for a suitable choice of $\left(u^{0}, u^{1}\right) \in H_{0}^{1}(\Omega) \times L^{2}(\Omega)$. Let us recall (see theorems 1.1 and 2.2) that for any given $\left(u^{0}, u^{1}\right) \in H_{0}^{1}(\Omega) \times L^{2}(\Omega)$ the problem (14)(16) has a unique solution, that $\partial_{\nu} u \in L^{2}\left(0, T ; L^{2}(\Gamma)\right)$, and that the linear map $\left(u^{0}, u^{1}\right) \mapsto \partial_{\nu} u$ is continuous from $H_{0}^{1}(\Omega) \times L^{2}(\Omega)$ into $L^{2}\left(0, T ; L^{2}(\Gamma)\right)$.

Using theorem 2.5 hence we deduce that the second problem

$$
\begin{gather*}
y^{\prime \prime}-\Delta y+q y=0 \quad \text { in } \quad \Omega \times(0, T)  \tag{17}\\
y=\partial_{\nu} u \quad \text { on } \quad \Gamma_{+} \times(0, T)  \tag{18}\\
y=0 \quad \text { on } \quad \Gamma_{-} \times(0, T)  \tag{19}\\
y(T)=y^{\prime}(T)=0 \tag{20}
\end{gather*}
$$

has a unique solution satisfying $\left(y(0), y^{\prime}(0)\right) \in L^{2}(\Omega) \times H^{-1}(\Omega)$ and that the linear map $\left(u^{0}, u^{1}\right) \mapsto\left(y(0), y^{\prime}(0)\right)$ is continuous from $H_{0}^{1}(\Omega) \times L^{2}(\Omega)$ into $L^{2}(\Omega) \times H^{-1}(\Omega)$.

If $\left(u^{0}, u^{1}\right)$ is such that $\left(y(0), y^{\prime}(0)\right)=\left(y^{0}, y^{1}\right)$, then the control $v:=\partial_{\nu} u$ on $\Gamma_{+}$and $v=0$ on $\Gamma_{-}$drives the system (6)-(8) in rest. Thus theorem 4.1 will be proved if we show the surjectivity of the map

$$
H_{0}^{1}(\Omega) \times L^{2}(\Omega) \ni\left(u^{0}, u^{1}\right) \mapsto\left(y(0), y^{\prime}(0)\right) \in L^{2}(\Omega) \times H^{-1}(\Omega)
$$

For some technical reasons it is more convenient to study the surjectivity of the map

$$
\Lambda: H_{0}^{1}(\Omega) \times L^{2}(\Omega) \rightarrow H^{-1}(\Omega) \times L^{2}(\Omega)
$$

defined by

$$
\Lambda\left(u^{0}, u^{1}\right):=\left(y^{\prime}(0),-y(0)\right)
$$

Clearly, the two maps are surjective at the same time.
In fact, we shall prove a stronger result :
Lemma 4.3. - Assume (10) and (11). Then $\Lambda$ is an isomorphism of $H_{0}^{1}(\Omega) \times L^{2}(\Omega)$ onto $H^{-1}(\Omega) \times L^{2}(\Omega)$.

Proof. - Clearly $\Lambda$ is a bounded linear map. Applying the Lax-Milgram theorem (see Brézis [2]), it suffices to show the existence of a constant $c>0$ such that, putting for brevity

$$
F:=H_{0}^{1}(\Omega) \times L^{2}(\Omega)
$$

we have

$$
\begin{equation*}
\left\langle\Lambda\left(u^{0}, u^{1}\right),\left(u^{0}, u^{1}\right)\right\rangle_{F^{\prime}, F} \geq c\left\|\left(u^{0}, u^{1}\right)\right\|_{F}^{2} \tag{21}
\end{equation*}
$$

for every $\left(u^{0}, u^{1}\right) \in F$. Since $\Lambda: F \rightarrow F^{\prime}$ is continuous and $Z \times Z$ is dense in $F$, it is sufficient to prove this inequality for $u^{0}, u^{1} \in Z$.

Multiplying the equation (17) by $u$ and integrating by parts we obtain

$$
\begin{gathered}
0=\int_{0}^{T} \int_{\Omega} u\left(y^{\prime \prime}-\Delta y+q y\right) \mathrm{dx} \mathrm{dt}=\left[\int_{\Omega} u y^{\prime}-u^{\prime} y \mathrm{dx}\right]_{0}^{T} \\
+\int_{0}^{T} \int_{\Omega}\left(u^{\prime \prime}-\Delta u+q u\right) y \mathrm{dx} \mathrm{dt}+\int_{0}^{T} \int_{\Gamma}-u \partial_{\nu} y+y \partial_{\nu} u \mathrm{~d} \Gamma \mathrm{dt} \\
\quad=\int_{\Omega}-u^{0} y^{\prime}(0)+u^{1} y(0) \mathrm{dx}+\int_{0}^{T} \int_{\Gamma_{+}}\left|\partial_{\nu} u\right|^{2} \mathrm{~d} \Gamma \mathrm{dt}
\end{gathered}
$$

whence

$$
\left\langle\Lambda\left(u^{0}, u^{1}\right),\left(u^{0}, u^{1}\right)\right\rangle_{F^{\prime}, F}=\int_{0}^{T} \int_{\Gamma_{+}}\left|\partial_{\nu} u\right|^{2} \mathrm{~d} \Gamma \mathrm{dt}
$$

By (10) and (11) we may apply theorem 3.1. We obtain the estimate

$$
\left\langle\Lambda\left(u^{0}, u^{1}\right),\left(u^{0}, u^{1}\right)\right\rangle_{F^{\prime}, F} \geq c^{\prime} E
$$

with a positive constant $c^{\prime}=c^{\prime}(T)$. Using the definition (3.5) of the energy hence we deduce (21) with $c:=c^{\prime} / 2$.

Remark 4.4. - HUM is based on the idea that the observability of the homogeneous problem (14)-(16) is sufficient for the exact controllability of the non-homogeneous problem (6)-(8). It is useful to observe that the observability is also a necessary condition. Indeed, assume that the problem (14)-(16) is not observable and fix non-zero initial data $\left(u^{0}, u^{1}\right) \in H_{0}^{1}(\Omega) \times L^{2}(\Omega)$ such that $\partial_{\nu} u \equiv 0$ on $\Gamma \times(0, T)$. Choose $\left(y^{0}, y^{1}\right) \in L^{2}(\Omega) \times H^{-1}(\Omega)$ such that

$$
\int_{\Omega} u^{0} y^{1}-u^{1} y^{0} \mathrm{dx} \neq 0
$$

Then there is no control $v \in L^{2}\left(0, T ; L^{2}(\Gamma)\right)$ such that the solution of (6)(8) satisfies $y(T)=y^{\prime}(T)=0$ on $\Omega$. To show this, fix $v \in L^{2}\left(0, T ; L^{2}(\Gamma)\right)$ arbitrarily and multiply (6) by the solution of (14)-(16). Integrating by parts and using (14), (15) and the property $\partial_{\nu} u \equiv 0$ on $\Gamma \times(0, T)$, we obtain that

$$
\begin{gathered}
0=\int_{0}^{T} \int_{\Omega} u\left(y^{\prime \prime}-\Delta y+q y\right) \mathrm{dx} \mathrm{dt} \\
=\left[\int_{\Omega} u y^{\prime}-u^{\prime} y \mathrm{dx}\right]_{0}^{T}-\int_{0}^{T} \int_{\Gamma} u \partial_{\nu} y-\left(\partial_{\nu} u\right) y \mathrm{~d} \Gamma \mathrm{dt} \\
+\int_{0}^{T} \int_{\Omega}\left(u^{\prime \prime}-\Delta u+q u\right) y \mathrm{dx} \mathrm{dt} \\
=\left[\int_{\Omega} u y^{\prime}-u^{\prime} y \mathrm{dx}\right]_{0}^{T} .
\end{gathered}
$$

Hence

$$
\int_{\Omega} u(T) y^{\prime}(T)-u^{\prime}(T) y(T) \mathrm{dx}=\int_{\Omega} u^{0} y^{1}-u^{1} y^{0} \mathrm{dx} \neq 0
$$

and therefore we cannot have $y(T)=y^{\prime}(T)=0$ on $\Omega$.

### 4.2. The first Petrovsky system

Fix $T>0$ and consider the problem

$$
\begin{gather*}
y^{\prime \prime}+\Delta^{2} y=0 \quad \text { in } \quad \Omega \times(0, T)  \tag{22}\\
y=0 \quad \text { and } \quad \partial_{\nu} y=v \quad \text { on } \quad \Gamma \times(0, T)  \tag{23}\\
y(0)=y^{0} \quad \text { and } \quad y^{\prime}(0)=y^{1} . \tag{24}
\end{gather*}
$$

It follows from theorem 2.9 that for every $\left(y^{0}, y^{1}, v\right) \in L^{2}(\Omega) \times H^{-2}(\Omega) \times$ $L^{2}\left(0, T ; L^{2}(\Gamma)\right)$ the problem $(22)-(24)$ has a unique solution $\left(y, y^{\prime}\right) \in$ $C\left([0, T] ; L^{2}(\Omega) \times H^{-2}(\Omega)\right)$.

Definition. - The problem (22)-(24) is exactly controllable if for any given $\left(y^{0}, y^{1}\right)$ and $\left(y_{T}^{0}, y_{T}^{1}\right) \in L^{2}(\Omega) \times H^{-2}(\Omega)$ there exists $v \in L^{2}\left(0, T ; L^{2}(\Gamma)\right)$ such that the solution of (22)-(24) satisfies

$$
y(T)=y_{T}^{0} \quad \text { and } \quad y^{\prime}(T)=y_{T}^{1}
$$

Let us denote by $\mu_{1}$ the first eigenvalue of the problem

$$
\Delta^{2} v=-\mu \Delta v, \quad v \in H_{0}^{2}(\Omega)
$$

Fix $x^{0} \in \mathbb{R}^{n}$ arbitrarily.
Theorem 4.5. - If $T>2 R / \sqrt{\mu_{1}}$, then for any given $\left(y^{0}, y^{1}\right)$ and $\left(y_{T}^{0}, y_{T}^{1}\right) \in L^{2}(\Omega) \times H^{-2}(\Omega)$ there exists $v \in L^{2}\left(0, T ; L^{2}(\Gamma)\right)$ such that

$$
\begin{equation*}
v=0 \quad \text { a.e. on } \quad \Gamma_{-} \times(0, T) \tag{25}
\end{equation*}
$$

and that the solution of $(22)-(24)$ satisfies

$$
y(T)=y_{T}^{0} \quad \text { and } \quad y^{\prime}(T)=y_{T}^{1}
$$

In particular, if $\Omega$ is contained in some ball of diameter $<T \sqrt{\mu_{1}}$, then the problem (22)-(24) is exactly controllable.

Theorem 4.5 was first proved by Lions [4], [5] under a stronger assumption on $T$; this assumption was weakened in Komornik [1].

Using an indirect compactness-uniqueness argument, Zuazua [1] later proved that these results remain valid in fact for arbitrarily small $T>0$. We shall prove his results in a constructive way in chapter 6 (see theorem 6.8).

Proof of theorem 4.5. - Using the same argument as in the preceding section, we may assume that $y_{T}^{0}=y_{T}^{1}=0$.

Fix $\left(u^{0}, u^{1}\right) \in H_{0}^{2}(\Omega) \times L^{2}(\Omega)$ arbitrarily. Solve the problem

$$
\begin{gathered}
u^{\prime \prime}+\Delta^{2} u=0 \quad \text { in } \quad \Omega \times(0, T) \\
u=\partial_{\nu} u=0 \quad \text { on } \quad \Gamma \times(0, T) \\
u(0)=u^{0} \quad \text { and } \quad u^{\prime}(0)=u^{1}
\end{gathered}
$$

and then the problem

$$
\begin{gathered}
y^{\prime \prime}+\Delta^{2} y=0 \quad \text { in } \quad \Omega \times(0, T), \\
y=0 \quad \text { and } \quad \partial_{\nu} y=\Delta u \quad \text { on } \quad \Gamma_{+} \times(0, T), \\
y=\partial_{\nu} y=0 \quad \text { on } \quad \Gamma_{-} \times(0, T), \\
y(T)=y^{\prime}(T)=0
\end{gathered}
$$

It follows from theorems 1.1, 2.6 and 2.9 that the formula

$$
\Lambda\left(u^{0}, u^{1}\right):=\left(y^{\prime}(0),-y(0)\right)
$$

defines a linear and continuous map of $F:=H_{0}^{2}(\Omega) \times L^{2}(\Omega)$ into $F^{\prime}$. It is sufficient to prove that $\Lambda$ is surjective. We establish a stronger result :

Lemma 4.6. - If $T>2 R / \sqrt{\mu_{1}}$, Then $\Lambda$ is an isomorphism of $F$ onto $F^{\prime}$.

Proof. - Applying the Lax-Milgram theorem, using the continuity of $\Lambda$ and the density of $Z \times Z$ in $F$ it suffices to prove the estimate

$$
\begin{equation*}
\left\langle\Lambda\left(u^{0}, u^{1}\right),\left(u^{0}, u^{1}\right)\right\rangle_{F^{\prime}, F} \geq c\left\|\left(u^{0}, u^{1}\right)\right\|_{F}^{2} \tag{26}
\end{equation*}
$$

for every $\left(u^{0}, u^{1}\right) \in Z \times Z$, with a suitable positive constant $c$.
First we observe that the right-hand side of (26) is equal to $2 c E$. Furthermore, the left-hand side of (26) equals $\int_{0}^{T} \int_{\Gamma_{+}}|\Delta u|^{2} \mathrm{~d} \Gamma \mathrm{dt}$. Indeed, we have

$$
\begin{gathered}
0=\int_{0}^{T} \int_{\Omega} u\left(y^{\prime \prime}+\Delta^{2} y\right) \mathrm{dx} \mathrm{dt} \\
=\int_{0}^{T} \int_{\Omega}\left(u^{\prime \prime}+\Delta^{2} u\right) y \mathrm{dx} \mathrm{dt}+\left[\int_{\Omega} u y^{\prime}-u^{\prime} y \mathrm{dx}\right]_{0}^{T} \\
+\int_{0}^{T} \int_{\Gamma} u\left(\partial_{\nu} \Delta y\right)-\left(\partial_{\nu} u\right)(\Delta y)+(\Delta u)\left(\partial_{\nu} y\right)-\left(\partial_{\nu} \Delta u\right) y \mathrm{~d} \Gamma \mathrm{dt} \\
=\int_{\Omega}-u^{0} y^{\prime}(0)+u^{1} y(0) \mathrm{dx}+\int_{0}^{T} \int_{\Gamma_{+}}(\Delta u)^{2} \mathrm{~d} \Gamma \mathrm{dt}
\end{gathered}
$$

whence

$$
\left\langle\Lambda\left(u^{0}, u^{1}\right),\left(u^{0}, u^{1}\right)\right\rangle_{F^{\prime}, F}=\int_{0}^{T} \int_{\Gamma_{+}}(\Delta u)^{2} \mathrm{~d} \Gamma \mathrm{dt}
$$

It remains to prove the estimate

$$
\begin{equation*}
\int_{0}^{T} \int_{\Gamma_{+}}(\Delta u)^{2} \mathrm{~d} \Gamma \mathrm{dt} \geq 2 c E \tag{27}
\end{equation*}
$$

and this follows from theorem 3.7 and from the hypothesis $T>2 R / \sqrt{\mu_{1}}$.

### 4.3. The wave equation. Neumann or Robin control

Consider the problem

$$
\begin{gather*}
y^{\prime \prime}-\Delta y+q y=0 \quad \text { in } \Omega \times \mathbb{R},  \tag{28}\\
y=0 \quad \text { on } \quad \Gamma_{0} \times \mathbb{R},  \tag{29}\\
\partial_{\nu} y+a y=v \quad \text { on } \quad \Gamma_{1} \times \mathbb{R},  \tag{30}\\
y(0)=y^{0} \quad \text { and } \quad y^{\prime}(0)=y^{1} \quad \text { in } \quad \Omega \tag{31}
\end{gather*}
$$

with $\Gamma_{1} \neq \emptyset$.

We begin by defininig the solution of (28)-(31). Now we do not have suitable hidden regularity results leading to optimal existence and uniqueness results expressed in terms of the usual Sobolev spaces; therefore the transposition (or duality) method will provide less precise existence results for the weak solutions.

As usual, we begin with a formal computation. Fix $\left(u^{0}, u^{1}\right) \in V \times H$ arbitrarily, solve the problem (cf. section 1.3)

$$
\begin{gather*}
u^{\prime \prime}-\Delta u+q u=0 \quad \text { in } \Omega \times \mathbb{R},  \tag{32}\\
u=0 \quad \text { on } \quad \Gamma_{0} \times \mathbb{R},  \tag{33}\\
\partial_{\nu} u+a u=0 \quad \text { on } \quad \Gamma_{1} \times \mathbb{R},  \tag{34}\\
u(0)=u^{0} \quad \text { and } \quad u^{\prime}(0)=u^{1} \quad \text { in } \quad \Omega, \tag{35}
\end{gather*}
$$

and multiply (32) by the solution of (28)-(31). We obtain for every fixed $S \in \mathbb{R}$ the equality

$$
\begin{aligned}
& 0=\int_{0}^{S} \int_{\Omega}\left(u^{\prime \prime}-\Delta u+q u\right) y \mathrm{dx} \mathrm{dt}=\left[\int_{\Omega} u^{\prime} y-u y^{\prime} \mathrm{dx}\right]_{0}^{S} \\
- & \int_{0}^{S} \int_{\Gamma}\left(\partial_{\nu} u\right) y-u\left(\partial_{\nu} y\right) \mathrm{d} \Gamma \mathrm{dt}+\int_{0}^{S} \int_{\Omega} u\left(y^{\prime \prime}-\Delta y+q y\right) \mathrm{dx} \mathrm{dt} \\
= & \int_{\Omega} u^{\prime}(S) y(S)-u(S) y^{\prime}(S)+u^{0} y^{1}-u^{1} y^{0} \mathrm{dx}+\int_{0}^{T} \int_{\Gamma_{1}} u v \mathrm{~d} \Gamma \mathrm{dt} .
\end{aligned}
$$

Putting

$$
L_{S}\left(u^{0}, u^{1}\right):=\int_{0}^{S} \int_{\Gamma_{1}} u v \mathrm{~d} \Gamma \mathrm{dt}+\left\langle\left(y^{1},-y^{0}\right),\left(u^{0}, u^{1}\right)\right\rangle_{V^{\prime} \times H, V \times H}
$$

we may rewrite this identity as
$L_{S}\left(u^{0}, u^{1}\right)=\left\langle\left(y^{\prime}(S),-y(S)\right),\left(u(S), u^{\prime}(S)\right)\right\rangle_{V^{\prime} \times H, V \times H}, \quad \forall\left(u^{0}, u^{1}\right) \in V \rtimes 3(6)$.
This leads to the following definition :
Definition. - We say that $\left(y, y^{\prime}\right)$ is a solution of (28)-(31) if $\left(y, y^{\prime}\right) \in C\left(\mathbb{R} ; H \times V^{\prime}\right)$ and if (36) is satisfied for every $S \in \mathbb{R}$.

We have the
Theorem 4.7. - Given $\left(y^{0}, y^{1}\right) \in H \times V^{\prime}$ and $v \in L_{l o c}^{2}\left(\mathbb{R} ; L^{2}\left(\Gamma_{1}\right)\right)$ arbitrarily, the problem (28)-(31) has a unique solution $\left(y, y^{\prime}\right) \in C\left(\mathbb{R} ; H \times V^{\prime}\right)$.

Furthermore, the linear map $\left(y^{0}, y^{1}, v\right) \mapsto\left(y, y^{\prime}\right)$ is continuous with respect to these topologies.

Proof. - We may easily adapt the proof of theorem 2.5 with one modification : to show the continuity of the linear form $L_{S}$ now we apply the standard trace theorem $(V \subset) H^{1}(\Omega) \hookrightarrow L^{2}(\Gamma)$

Now assume, as in section 3.4, that there exists $x^{0} \in \mathbb{R}^{n}$ such that

$$
\begin{equation*}
m \cdot \nu \leq 0 \quad \text { on } \quad \Gamma_{0}, \quad m \cdot \nu \geq 0 \quad \text { on } \quad \Gamma_{1}, \tag{37}
\end{equation*}
$$

that $a$ has the form

$$
\begin{equation*}
a=(m \cdot \nu) b, \quad b \in C^{1}\left(\Gamma_{1}\right), \quad b \geq 0 \quad \text { on } \quad \Gamma_{1}, \tag{38}
\end{equation*}
$$

and that

$$
\begin{equation*}
R^{2} b^{2}+(2-n) b \leq 0 \quad \text { on } \quad \Gamma_{1} . \tag{39}
\end{equation*}
$$

Let us introduce the constants $R_{1}$ and $Q_{1}$ as in section 3.4 and assume that

$$
\begin{equation*}
Q_{1}<1 . \tag{40}
\end{equation*}
$$

Theorem 4.8. - Assume (37)-(40) and let

$$
\begin{equation*}
T>2 R_{1} /\left(1-Q_{1}\right) \tag{41}
\end{equation*}
$$

Then for any given $\left(y^{0}, y^{1}\right),\left(y_{T}^{0}, y_{T}^{1}\right) \in V \times H$ there exists a control $v \in L^{2}\left(0, T ; L^{2}\left(\Gamma_{1}\right)\right)$ such that (extending $v$ by zero outside $\left.(0, T)\right)$ the solution of (28)-(31) satisfies

$$
y(T)=y_{T}^{0} \quad \text { and } \quad y^{\prime}(T)=y_{T}^{1}
$$

For $\Gamma_{0} \neq \emptyset$ and $q \equiv 0$ we obtain a weakened version of a theorem of Lions [4; p. 203] : he proved this result (using indirect compactness-uniqueness arguments) with $R_{1}$ replaced by $R$ in (41). Theorem 4.8 will be improved later (see theorems 6.20 and 8.10).

Proof. - As usual, we may assume that $y_{T}^{0}=y_{T}^{1} \equiv 0$. We apply HUM.
Consider on $Z \times Z$ the seminorm defined by

$$
\left\|\left(u^{0}, u^{1}\right)\right\|_{F}:=\|u\|_{L^{2}\left(0, T ; L^{2}\left(\Gamma_{1}\right)\right)}
$$

where $u$ is the solution of (32)-(35). By theorem 3.26 it is a norm.
Completing $Z \times Z$ with respect to this norm we obtain a Hilbert space $F$ and we deduce from the trace theorem $V \hookrightarrow L^{2}(\Gamma)$ and from theorem 3.26 the algebraical and topological inclusions

$$
\begin{equation*}
V \times H \subset F \subset H \times V^{\prime} . \tag{42}
\end{equation*}
$$

Given $\left(u^{0}, u^{1}\right) \in F$ arbitrarily, we solve (32)-(35), then we solve

$$
\begin{gather*}
y^{\prime \prime}-\Delta y+q y=0 \quad \text { in } \quad \Omega \times \mathbb{R},  \tag{43}\\
y=0 \quad \text { on } \quad \Gamma_{0} \times \mathbb{R},  \tag{44}\\
\partial_{\nu} y+a y=-u \quad \text { on } \quad \Gamma_{1} \times \mathbb{R},  \tag{45}\\
y(T)=y^{\prime}(T)=0 \tag{46}
\end{gather*}
$$

and we set

$$
\Lambda\left(u^{0}, u^{1}\right)=\left(y^{\prime}(0),-y(0)\right) .
$$

By the definition of $F$ and by theorem $4.7 \Lambda$ is a bounded linear map of $F$ into $H \times V^{\prime}$. The proof will be completed if we show that $F^{\prime}$ is contained in the range of $\Lambda$. Indeed, then we will conclude from (42) that $H \times V \subset F^{\prime}$. In fact, the following stronger result holds true :
Lemma 4.9. - $\Lambda$ is an isomorphism of $F$ onto $F^{\prime}$.
Proof. - Compute

$$
\left\langle\Lambda\left(u^{0}, u^{1}\right),\left(v^{0}, v^{1}\right)\right\rangle_{V^{\prime} \times H, V \times H}
$$

for $\left(u^{0}, u^{1}\right),\left(v^{0}, v^{1}\right) \in Z \times Z$. Denoting by $u, v$ the corresponding solutions of (32)-(35) and considering the corresponding solution $y$ of (43)-(46), we have

$$
\begin{gathered}
0=\int_{0}^{T} \int_{\Omega}\left(y^{\prime \prime}-\Delta y+q y\right) v \mathrm{dx} \mathrm{dt}=\left[\int_{\Omega} y^{\prime} v-y v^{\prime} \mathrm{dx}\right]_{0}^{T} \\
-\int_{0}^{T} \int_{\Gamma}\left(\partial_{\nu} y\right) v-y \partial_{\nu} v \mathrm{~d} \Gamma \mathrm{dt}+\int_{0}^{T} \int_{\Omega} y\left(v^{\prime \prime}-\Delta v+q v\right) \mathrm{dx} \mathrm{dt} \\
=\int_{\Omega} y^{\prime}(T) v(T)-y(T) v^{\prime}(T)-y^{\prime}(0) v(0)+y(0) v\left(^{\prime} 0\right) \mathrm{dx}+\int_{0}^{T} \int_{\Gamma_{1}} u v \mathrm{~d} \Gamma \mathrm{dt} \\
=\int_{\Omega}-y^{\prime}(0) v(0)+y(0) v^{\prime}(0) \mathrm{dx}+\int_{0}^{T} \int_{\Gamma_{1}} u v \mathrm{~d} \Gamma \mathrm{dt}
\end{gathered}
$$

whence

$$
\left\langle\Lambda\left(u^{0}, u^{1}\right),\left(v^{0}, v^{1}\right)\right\rangle_{V^{\prime} \times H, V \times H}=(u, v)_{L^{2}\left(0, T ; L^{2}\left(\Gamma_{1}\right)\right)},
$$

i.e.

$$
\begin{equation*}
\left\langle\Lambda\left(u^{0}, u^{1}\right),\left(v^{0}, v^{1}\right)\right\rangle_{V^{\prime} \times H, V \times H}=\left(\left(u^{0}, u^{1}\right),\left(v^{0}, v^{1}\right)\right)_{F} . \tag{47}
\end{equation*}
$$

We conclude from (47) that $\Lambda\left(u^{0}, u^{1}\right) \in F^{\prime}$ and that

$$
\left\|\Lambda\left(u^{0}, u^{1}\right)\right\|_{F^{\prime}}=\left\|\left(u^{0}, u^{1}\right)\right\|_{F}
$$

Consequently, $\Lambda$ is a bounded linear map of $F$ into $F^{\prime}$. By (47) we may apply the Lax-Milgram theorem and the proof is completed.

Remark 4.10. - Let us recall from remark 3.20 that in certain cases we have $R_{1} \leq R$.

## 5. Norm inequalities

The aim of this chapter is to introduce a general method, which will permit us in the following chapter to improve and complete the uniqueness and exact controllability results obtained earlier, by weakening the hypotheses on the length of the intervals of uniqueness and on the sufficient time of exact controllability.

This method is closely related to an estimation method set out by Haraux [3]. In section 5.4 we shall also apply an idea of Lebeau [1].

Throughout this chapter all spaces are assumed to be complex.
In section 5.1 we outline the main ideas in the case of the first Petrovsky system. The precise results will be formulated and proved in sections 5.2-5.4.

### 5.1. Riesz sequences

Let $\left(\omega_{j}\right)_{j \geq 1}$ be a sequence of distinct real numbers and consider the functions of the form

$$
\begin{equation*}
u(t):=\sum_{j=1}^{\infty} z_{j} e^{i \omega_{j} t}, \quad t \in \mathbb{R} \tag{1}
\end{equation*}
$$

with complex scalar coefficients $z_{j}$. We say that $\left(e^{i \omega_{j} t}\right)_{j \geq 1}$ is a Riesz sequence on $(0, T)$ if there exist two positive constants $c_{1}(T)$ and $c_{2}(T)$ such that

$$
\begin{equation*}
c_{1}(T) \sum_{j=1}^{\infty}\left|z_{j}\right|^{2} \leq \int_{0}^{T}|u(t)|^{2} \mathrm{dt} \leq c_{2}(T) \sum_{j=1}^{\infty}\left|z_{j}\right|^{2} \tag{2}
\end{equation*}
$$

for every function $u$ of the form (1).
It is clear that if $\left(e^{i \omega_{j} t}\right)_{j>1}$ is a Riesz sequence on $(0, T)$, then for every $k>1$ the subsequence $\left(e^{i \omega_{j} t}\right)_{j \geq k}$ is also a Riesz sequence on $(0, T)$. Conversely, the following standard result holds true (see e.g. Ball and Slemrod [1] or Haraux [3]) :

Let $k>1$ and assume that the subsequence $\left(e^{i \omega_{j} t}\right)_{j \geq k}$ is a Riesz sequence on some interval $\left(0, T_{k}\right)$. Then the entire sequence $\left(e^{i \omega_{j} t}\right){ }_{j \geq 1}$ is a Riesz sequence on $(0, T)$ for all $T>T_{k}$, arbitrarily close to $T_{k}$.

Now consider the problem

$$
\begin{gather*}
u^{\prime \prime}+\Delta^{2} u=0 \quad \text { in } \quad \Omega \times \mathbb{R}  \tag{3}\\
u=\partial_{\nu} u=0 \quad \text { on } \quad \Gamma \times \mathbb{R}  \tag{4}\\
u(0)=u^{0} \quad \text { and } \quad u^{\prime}(0)=u^{1} \tag{5}
\end{gather*}
$$

the energy $E$ of the solution is defined by

$$
\begin{equation*}
E=\frac{1}{2} \int_{\Omega}|\Delta u|^{2}+\left|u^{\prime}\right|^{2} \mathrm{dx} \tag{6}
\end{equation*}
$$

Given $T>2 R \lambda_{1}^{-1 / 4}$ arbitrarily, by the results of sections 2.3 and 3.2 (see theorems 2.6, 3.6 and remark 3.12) there exist two constants $c_{1}(T), c_{2}(T)$ such that for every $u^{0}, u^{1} \in Z$ the solution of (3)-(5) satisfies the estimates

$$
\begin{equation*}
c_{1}(T) E \leq \int_{0}^{T} \int_{\Gamma_{+}}|\Delta u|^{2} \mathrm{~d} \Gamma \mathrm{dt} \leq c_{2}(T) E \tag{7}
\end{equation*}
$$

Let us recall from section 4.2 that the first inequality in (7) implies the exact controllability of the first Petrovsky system in time $T$.

In fact, the estimates (7) hold for arbitrarily small $T>0$. To convince ourselves, set

$$
\omega_{2 j-1}=\sqrt{\lambda_{j}}, \omega_{2 j}=-\sqrt{\lambda_{j}}, \widetilde{Z}_{2 j-1}=\widetilde{Z}_{2 j}=Z_{j}, j \geq 1
$$

and introduce in $Z$ the norm

$$
\|v\|:=\left(\int_{\Omega}|\Delta v|^{2} \mathrm{dx}\right)^{1 / 2}
$$

and the semi-norm

$$
|v|:=\left(\int_{\Gamma_{+}}|\Delta v|^{2} \mathrm{dx}\right)^{1 / 2}
$$

Then the solution of (3)-(5) may be written in the form (1) (see theorem 1.3) with vector coefficients $z_{j} \in \widetilde{Z}_{j}$, and one can readily verify that

$$
\begin{equation*}
E=\sum_{j=1}^{\infty}\left\|z_{j}\right\|^{2} \tag{8}
\end{equation*}
$$

Consequently, we may rewrite the estimates (7) in a form analogous to (2) :

$$
\begin{equation*}
c_{1}(T) \sum_{j=1}^{\infty}\left\|z_{j}\right\|^{2} \leq \int_{0}^{T}|u(t)|^{2} \mathrm{dt} \leq c_{2}(T) \sum_{j=1}^{\infty}\left\|z_{j}\right\|^{2} . \tag{9}
\end{equation*}
$$

By analogy with the scalar case, let us say that $\left(e^{i \omega_{j} t}, \widetilde{Z}_{j}\right)_{j \geq 1}$ is a vector Riesz sequence on $(0, T)$ if the estimates (9) are satisfied.

Now fix $T>0$ arbitrarily. Choose $k>1$ such that $T>2 R \lambda_{k}^{-1 / 4}$ and then choose $T_{k}$ such that $T>T_{k}>2 R \lambda_{k}^{-1 / 4}$. By theorems 2.6, 3.6 and remark 3.12 the subsequence $\left(e^{i \omega_{j} t}, \widetilde{Z}_{j}\right)_{j \geq 2 k-1}$ is a vector Riesz sequence on $\left(0, T_{k}\right)$ : observe that

$$
\begin{equation*}
z_{1}=\cdots=z_{2 k-2}=0 \Longleftrightarrow u^{0}, u^{1} \perp Z_{1}, \cdots, Z_{k-1} . \tag{10}
\end{equation*}
$$

Applying a generalization of the above mentioned scalar result, hence we will conclude that the entire sequence $\left(e^{i \omega_{j} t}\right)_{j \geq 1},\left(\widetilde{Z}_{j}\right)_{j \geq 1}$ is a vector Riesz sequence on $(0, T)$ (see theorem 5.2 below) ; in other words, the estimates (7) hold for arbitrarily small $T>0$.

### 5.2. Formulation of the results

Let $\mathcal{A}$ be a linear operator in an infinite-dimensional complex Hilbert space $\mathcal{H}$. Assume that $\mathcal{A}$ has an infinite sequence of purely imaginary eigenvalues $i \omega_{j}\left(\omega_{j} \in \mathbb{R}, j=1,2, \ldots\right)$ satisfying

$$
\begin{equation*}
\left|\omega_{j}\right| \rightarrow+\infty \quad \text { as } \quad j \rightarrow+\infty \tag{11}
\end{equation*}
$$

and a corresponding sequence of finite-dimensional, pairwise orthogonal eigenspaces $\mathcal{Z}_{j}$ whose linear hull $\mathcal{Z}$ is dense in $\mathcal{H}$. Then for any given $U^{0} \in \mathcal{Z}$ the initial value problem

$$
\begin{equation*}
U^{\prime}=\mathcal{A} U \quad \text { in } \quad \mathbb{R}, \quad U(0)=U^{0} \tag{12}
\end{equation*}
$$

has a unique solution $U \in C^{\infty}(\mathbb{R} ; \mathcal{H})$ and this solution has a unique expansion of the form

$$
\begin{equation*}
U(t)=\sum_{j} U_{j} e^{i \omega_{j} t}, \quad U_{j} \in \mathcal{Z}_{j} \tag{13}
\end{equation*}
$$

where finitely many coefficients $U_{j}$ are different from zero only. It follows that

$$
\begin{equation*}
\|U(t)\|_{\mathcal{H}}=\left\|U^{0}\right\|_{\mathcal{H}}, \quad \forall t \in \mathbb{R} \tag{14}
\end{equation*}
$$

(Our hypotheses mean that $\mathcal{A}$ has a skew-adjoint extension (i.e. $i \mathcal{A}$ is self-adjoint) in $\mathcal{H}$, having a compact resolvent.)

Example 5.1. - Consider the abstract problem

$$
\begin{equation*}
u^{\prime \prime}+A u=0 \quad \text { in } \quad \mathbb{R}, \quad u(0)=u^{0}, \quad u^{\prime}(0)=u^{1} \tag{15}
\end{equation*}
$$

introduced in section 1.2. We can rewrite it in the form (12) by setting

$$
\mathcal{A}(u, v)=(v,-A u), \quad U=\left(u, u^{\prime}\right), \quad \text { and } \quad U^{0}=\left(u^{0}, u^{1}\right) .
$$

Introducing the Hilbert space $\mathcal{H}=V \times H$ and the sequences $\left(\omega_{j}\right),\left(\mathcal{Z}_{j}\right)$ by

$$
\omega_{2 k-1}=\sqrt{\lambda_{k}}, \quad \mathcal{Z}_{2 k-1}=\left\{\left(v, i \omega_{2 k-1} v\right): v \in Z_{k}\right\}
$$

and

$$
\omega_{2 k}=-\sqrt{\lambda_{k}}, \quad \mathcal{Z}_{2 k}=\left\{\left(v, i \omega_{2 k} v\right): v \in Z_{k}\right\}
$$

for $k=1,2, \ldots$, the above mentioned conditions can be easily verified. Furthermore, the energy of the solutions of (15) is closely related to the norm of the "energy space" $\mathcal{H}$ :

$$
\left\|U^{0}\right\|_{\mathcal{H}}=\left\|\left(u^{0}, u^{1}\right)\right\|_{\mathcal{H}}=2 E .
$$

Let us observe that for any fixed $U^{0}=\left(u^{0}, u^{1}\right) \in \mathcal{Z}$ and for any positive integer $n$ the following properties are equivalent :

$$
\begin{gathered}
U^{0} \perp \mathcal{Z}_{j} \quad \text { in } \mathcal{H} \text { for } \quad j=1, \ldots, 2 n-2 \\
u^{0}, u^{1} \perp Z_{j} \quad \text { in } H(\text { or in } V) \text { for } \quad j=1, \ldots, n-1 .
\end{gathered}
$$

Indeed, writing

$$
U^{0}=\left(u^{0}, u^{1}\right)=\sum_{k}\left(v_{k}, i \omega_{2 k-1} v_{k}\right)+\left(w_{k}, i \omega_{2 k} w_{k}\right), \quad v_{k}, w_{k} \in Z_{k}
$$

for every fixed $m \geq 1$ we have the following equivalences:

$$
\begin{gathered}
u^{0}, u^{1} \perp Z_{m} \Longleftrightarrow v_{m}+w_{m}=v_{m}-w_{m}=0 \Longleftrightarrow v_{m}=w_{m}=0 \\
U^{0} \perp \mathcal{Z}_{2 m-1}, Z_{2 m} \Longleftrightarrow v_{m}=w_{m}=0 .
\end{gathered}
$$

Now let $p$ be a semi-norm in $\mathcal{Z}$. The following theorem will play a crucial role in the following chapter.

Theorem 5.2. - Assume that $p$ is a norm in each of the eigenspaces $\mathcal{Z}_{j}$ :

$$
\begin{equation*}
U \in \mathcal{Z}_{j} \quad \text { and } \quad p(U)=0 \Longrightarrow U=0 \quad(j=1,2, \ldots) . \tag{16}
\end{equation*}
$$

Assume that there exist an integer $k>1$, two intervals $I_{1}, I_{2}$ and two positive constants $c_{1}, c_{2}$ such that the solutions of (12) satisfy the inequalities

$$
\begin{equation*}
\int_{I_{1}} p(U(t))^{2} \mathrm{dt} \geq c_{1}\left\|U^{0}\right\|_{\mathcal{H}}^{2} \tag{17}
\end{equation*}
$$

and

$$
\begin{equation*}
\int_{I_{2}} p(U(t))^{2} \mathrm{dt} \leq c_{2}\left\|U^{0}\right\|_{\mathcal{H}}^{2} \tag{18}
\end{equation*}
$$

whenever

$$
\begin{equation*}
U^{0} \perp \mathcal{Z}_{1}, \ldots, \mathcal{Z}_{k-1} \tag{19}
\end{equation*}
$$

Then for every interval I of length $|I|>\left|I_{1}\right|$ there exist two positive constants $c_{3}$ and $c_{4}$ such that

$$
\begin{equation*}
c_{3}\left\|U^{0}\right\|_{\mathcal{H}}^{2} \leq \int_{I} p(U(t))^{2} \mathrm{dt} \leq c_{4}\left\|U^{0}\right\|_{\mathcal{H}}^{2} \tag{20}
\end{equation*}
$$

for all $U^{0} \in \mathcal{Z}$.
In fact, we shall prove the following more general result :
Theorem 5.3. - Assume that there exist an integer $k>1$, two intervals $I_{1}, I_{2}$, a semi-norm $q$ in $\mathcal{Z}$ and three positive constants $c_{0}, c_{1}, c_{2}$ such that

$$
\begin{gather*}
U \in \mathcal{Z}_{j} \quad \text { and } \quad q(U)=0 \Longrightarrow U=0 \quad(j=1,2, \ldots),  \tag{21}\\
q \leq c_{0} p \quad \text { in } \mathcal{Z} \tag{22}
\end{gather*}
$$

and that the solutions of (12) satisfy the inequalities

$$
\begin{equation*}
\int_{I_{1}} p(U(t))^{2} \mathrm{dt} \geq c_{1}\left\|U^{0}\right\|_{\mathcal{H}}^{2} \tag{23}
\end{equation*}
$$

and

$$
\begin{equation*}
\int_{I_{2}} q(U(t))^{2} \mathrm{dt} \leq c_{2}\left\|U^{0}\right\|_{\mathcal{H}}^{2} \tag{24}
\end{equation*}
$$

whenever

$$
\begin{equation*}
U^{0} \perp \mathcal{Z}_{1}, \ldots, \mathcal{Z}_{k-1} \tag{25}
\end{equation*}
$$

Then for every interval I of length $|I|>\left|I_{1}\right|$ there exists a positive constant $c_{3}$ such that

$$
\begin{equation*}
\int_{I} p(U(t))^{2} \mathrm{dt} \geq c_{3}\left\|U^{0}\right\|_{\mathcal{H}}^{2} \tag{26}
\end{equation*}
$$

and for every interval I there exists a positive constant $c_{4}$ such that

$$
\begin{equation*}
\int_{I} q(U(t))^{2} \mathrm{dt} \leq c_{4}\left\|U^{0}\right\|_{\mathcal{H}}^{2} \tag{27}
\end{equation*}
$$

for all $U^{0} \in \mathcal{Z}$.
For $q=p$ this theorem reduces to theorem 5.2.
Remark 5.4. - The following remark is often useful to verify the condition (21) or (16). Assume that for some given $j$ there is an interval $I^{\prime}$ and a positive constant $c^{\prime}$ such that

$$
\begin{equation*}
\int_{I^{\prime}} q(U(t))^{2} \mathrm{dt} \geq c^{\prime}\left\|U^{0}\right\|_{\mathcal{H}}^{2} \tag{28}
\end{equation*}
$$

for all $U^{0} \in \mathcal{Z}_{j}$. Then $q$ is a norm in $\mathcal{Z}_{j}$. Indeed, if $U^{0} \in \mathcal{Z}_{j}$ then the corresponding solution of (12) is

$$
U(t)=U^{0} e^{i \omega_{j} t}
$$

and therefore

$$
\left|I^{\prime}\right| q\left(U^{0}\right)^{2}=\int_{I^{\prime}} q(U(t))^{2} \mathrm{dt} \geq c^{\prime}\left\|U^{0}\right\|_{\mathcal{H}}^{2}
$$

Hence $U^{0} \neq 0$ implies $q\left(U^{0}\right) \neq 0$.
Sometimes we shall find $I^{\prime}$ and $c^{\prime}$ satisfying (28) for all $U^{0} \in \mathcal{H}$. However, the interval $I^{\prime}$ will be much longer than $I_{1}$ in theorems 5.2 and 5.3 , hence it will not replace (17) or (23).

We shall also prove a euclidean version of theorem 5.3. Let $p$ be a euclidean semi-norm in $\mathcal{Z}$ (i.e. defined by a positive semi-definite hermitian bilinear form $p(\cdot, \cdot)$ in $\mathcal{Z}$ such that $p(U)=p(U, U)^{1 / 2}$ for all $\left.U \in \mathcal{Z}\right)$ and assume that there exist two positive constants $\alpha$ and $c_{\alpha}$ such that

$$
\begin{equation*}
p(U) \leq c_{\alpha}\left|\omega_{j}\right|^{\alpha}\|U\|_{\mathcal{H}}, \quad \forall U \in \mathcal{Z}_{j}, \quad \forall \omega_{j} \neq 0 \tag{29}
\end{equation*}
$$

We also need the following assumption on the spectrum of $\mathcal{A}$ : there is a positive constant $d$ satisfying

$$
\begin{equation*}
\sum_{j=1}^{\infty}\left|\omega_{j}\right|^{-d}<+\infty \tag{30}
\end{equation*}
$$

Let us denote by $\mathcal{Z}_{+}$(resp. by $\mathcal{Z}_{-}$) the linear hull of the eigenspaces $\mathcal{Z}_{j}$ corresponding to the eigenvalues $i \omega_{j}$ with $\omega_{j} \geq 0$ (resp. $\omega_{j}<0$ ). Clearly we have

$$
\begin{equation*}
\mathcal{Z}_{+} \perp \mathcal{Z}_{-} \quad \text { and } \quad \mathcal{Z}=\mathcal{Z}_{+}+\mathcal{Z}_{-} \tag{31}
\end{equation*}
$$

Theorem 5.5. - Assume (30) and let p be a euclidean semi-norm in $\mathcal{Z}$ satisfying (29). Assume that there exist an integer $l>1$, two intervals $I_{1}^{\prime}, I_{2}$, a semi-norm $q$ in $\mathcal{Z}$ and three positive constants $c_{0}, c_{1}^{\prime}$ and $c_{2}^{\prime}$ satisfying (21), (22) and such that the solutions of (12) satisfy the inequalities

$$
\begin{equation*}
\int_{I_{1}^{\prime}} p(U(t))^{2} \mathrm{dt} \geq c_{1}^{\prime}\left\|U^{0}\right\|_{\mathcal{H}}^{2} \tag{32}
\end{equation*}
$$

and

$$
\begin{equation*}
\int_{I_{2}} q(U(t))^{2} \mathrm{dt} \leq c_{2}^{\prime}\left\|U^{0}\right\|_{\mathcal{H}}^{2} \tag{33}
\end{equation*}
$$

whenever

$$
\begin{equation*}
U^{0} \perp \mathcal{Z}_{1}, \ldots, \mathcal{Z}_{l-1} \tag{34}
\end{equation*}
$$

and

$$
\begin{equation*}
\text { either } \quad U^{0} \in \mathcal{Z}_{+} \quad \text { or } \quad U^{0} \in \mathcal{Z}_{-} . \tag{35}
\end{equation*}
$$

Then for every interval I of length $|I|>\left|I_{1}^{\prime}\right|$ there exist two positive constants $c_{3}$ and $c_{4}$ such that

$$
\begin{equation*}
\int_{I} p(U(t))^{2} \mathrm{dt} \geq c_{3}\left\|U^{0}\right\|_{\mathcal{H}}^{2} \tag{36}
\end{equation*}
$$

and

$$
\begin{equation*}
\int_{I} q(U(t))^{2} \mathrm{dt} \leq c_{4}\left\|U^{0}\right\|_{\mathcal{H}}^{2} \tag{37}
\end{equation*}
$$

for all $U^{0} \in \mathcal{Z}$.
The rest of this chapter is devoted to the proof of theorems 5.3 and 5.5.

### 5.3. Proof of theorem 5.3

First of all, we may assume by an obvious induction argument that $k=2$. Furthermore, we may assume that $\omega_{1}=0$. To see this let us consider the initial value problem

$$
\begin{equation*}
V^{\prime}=\left(\mathcal{A}-i \omega_{1} \mathcal{I}\right) V \quad \text { in } \quad \mathbb{R}, \quad V(0)=U^{0} \tag{38}
\end{equation*}
$$

The solutions of (12) and (38) are clearly connected by the relation

$$
V(t) \equiv U(t) e^{-i \omega_{1} t}
$$

Since

$$
p(V(t)) \equiv p(U(t)) \quad \text { and } \quad q(V(t)) \equiv q(U(t))
$$

the inequalities $(23),(24),(26),(27)$ of theorem 5.3 are the same for $U(t)$ and for $V(t)$. Furthermore, the eigenvalues of $\mathcal{A}-i \omega_{1} \mathcal{I}$ are those of $\mathcal{A}$ shifted by $-i \omega_{1}$; in particular the eigenvalue $i \omega_{1}$ of $\mathcal{A}$ corresponds to the zero eigenvalue of $\mathcal{A}-i \omega_{1} \mathcal{I}$.

Step 1. - First we establish a weakened version of inequality (26) : there is a positive constant $c^{\prime}$ such that

$$
\begin{equation*}
\int_{I} p(U(t))^{2} \mathrm{dt} \geq c^{\prime}\left\|U^{0}-U_{1}\right\|_{\mathcal{H}}^{2} \tag{39}
\end{equation*}
$$

for all $U^{0} \in \mathcal{Z}$ where $U_{1}$ denotes the orthogonal projection of $U^{0}$ onto $\mathcal{Z}_{1}$. (Since $\omega_{1}=0, U_{1}$ is also the orthogonal projection of $U(t)$ onto $\mathcal{Z}_{1}$ for every $t \in \mathbb{R}$. In other words, $U_{1}$ is the constant part of $U(t)$, cf. (13).)

Let us first note that inequality (23) remains true (with the same constant $c_{1}$ ) for every translate $I_{1}+\tau$ of $I_{1}$. Indeed, set

$$
\begin{equation*}
V(t):=U(t+\tau), \quad t \in \mathbb{R}, \tag{40}
\end{equation*}
$$

then $V$ is the solution of (12) with $U^{0}$ replaced by $U(\tau)$. Using (23) and (14) we have

$$
\int_{I_{1}+\tau} p(U(t))^{2} \mathrm{dt}=\int_{I_{1}} p(V(t))^{2} \mathrm{dt} \geq c_{1}\|V(0)\|_{\mathcal{H}}^{2}=c_{1}\|U(\tau)\|_{\mathcal{H}}^{2}=c_{1}\left\|U^{0}\right\|_{\mathcal{H}}^{2}
$$

Since $I$ is longer than $I_{1}$, we may therefore assume that $I$ contains the closure of $I_{1}$ in its interior, say

$$
\begin{equation*}
I_{1}=(a, b) \quad \text { and } \quad(a-\varepsilon, b+\varepsilon) \subset I, \quad \varepsilon>0 \tag{41}
\end{equation*}
$$

Now fix $U^{0} \in \mathcal{Z}$ arbitrarily, solve (12) and (following Haraux [3]) set

$$
\begin{equation*}
V(t):=U(t)-\frac{1}{2 \varepsilon} \int_{-\varepsilon}^{\varepsilon} U(t+s) \mathrm{ds}, t \in \mathbb{R} \tag{42}
\end{equation*}
$$

Observe that

$$
\begin{equation*}
\int_{I_{1}} p(V(t))^{2} \mathrm{dt} \leq 4 \int_{I} p(U(t))^{2} \mathrm{dt} \tag{43}
\end{equation*}
$$

Indeed, we have

$$
\begin{aligned}
& p(V(t))^{2} \leq 2 p(U(t))^{2}+2 p\left(\frac{1}{2 \varepsilon} \int_{-\varepsilon}^{\varepsilon} U(t+s) \mathrm{ds}\right)^{2} \\
& \leq 2 p(U(t))^{2}+\frac{1}{2 \varepsilon^{2}}\left(\int_{-\varepsilon}^{\varepsilon} p(U(t+s)) \mathrm{ds}\right)^{2} \\
& \quad \leq 2 p(U(t))^{2}+\frac{1}{\varepsilon} \int_{-\varepsilon}^{\varepsilon} p(U(t+s))^{2} \mathrm{ds}
\end{aligned}
$$

for every $t \in \mathbb{R}$, whence, using (41) and the Fubini theorem,

$$
\begin{gathered}
\int_{I_{1}} p(V(t))^{2} \mathrm{dt} \leq 2 \int_{I_{1}} p(U(t))^{2} \mathrm{dt}+\frac{1}{\varepsilon} \int_{I_{1}} \int_{-\varepsilon}^{\varepsilon} p(U(t+s))^{2} \mathrm{ds} \mathrm{dt} \\
\quad \leq 2 \int_{I_{1}} p(U(t))^{2} \mathrm{dt}+2 \int_{a-\varepsilon}^{b+\varepsilon} p\left(U\left(t^{\prime}\right)\right)^{2} \mathrm{dt}^{\prime} \leq 4 \int_{I} p(U(t))^{2} \mathrm{dt} .
\end{gathered}
$$

Using the expansion (13) of $U(t)$ one computes easily that

$$
V(t)=\sum_{j} V_{j} e^{i \omega_{j} t}, \quad V_{j} \in \mathcal{Z}_{j}
$$

with

$$
\begin{equation*}
V_{1}=0 \quad \text { and } \quad V_{j}=\left(1-\frac{\sin \omega_{j} \varepsilon}{\omega_{j} \varepsilon}\right) U_{j} \quad \text { if } \quad j>1 \tag{44}
\end{equation*}
$$

Hence $V$ is the solution of (12) with $U^{0}$ replaced by $V^{0}:=\sum_{j} V_{j}$, which satisfies the orthogonality condition (25). Applying (23) we obtain that

$$
\begin{equation*}
\int_{I_{1}} p(V(t))^{2} \mathrm{dt} \geq c_{1}\left\|V^{0}\right\|_{\mathcal{H}}^{2} \tag{45}
\end{equation*}
$$

Since $\omega_{j} \neq 0\left(=\omega_{1}\right)$ for $j \neq 1$ and since $\left|\omega_{j}\right| \rightarrow+\infty$ as $j \rightarrow+\infty($ cf. (11)), there exists a positive constant $\alpha$ such that

$$
\begin{equation*}
1-\frac{\sin \omega_{j} \varepsilon}{\omega_{j} \varepsilon} \geq \alpha, \quad \forall j \geq 2 \tag{46}
\end{equation*}
$$

Using the orthogonality of the eigenspaces $\mathcal{Z}_{j}$ we deduce from (44) and (46) that

$$
\begin{equation*}
\left\|V^{0}\right\|_{\mathcal{H}} \geq \alpha\left\|U^{0}-U_{1}\right\|_{\mathcal{H}} \tag{47}
\end{equation*}
$$

and (39) follows from (43), (45) and (47) with $c^{\prime}:=\alpha^{2} c_{1} / 4$.
Step 2. - Next we prove the (easy) estimate (27) for an arbitrary interval $I$. (Its length does not play any role here.) Let us first note that the estimate (24) remains true (with another constant) without the orthogonality assumption (25). Indeed, since the square root of the left hand side of (24) defines a semi-norm of $U^{0}$ and since this semi-norm is obviously majorized by the norm of $\mathcal{H}$ on the finite-dimensional vector space $\mathcal{Z}_{1}$, there exists a positive constant $c_{2}^{\prime}$ such that

$$
\int_{I_{2}} q\left(U_{1}\right)^{2} \mathrm{dt} \leq c_{2}^{\prime}\left\|U_{1}\right\|_{\mathcal{H}}^{2}, \quad \forall U_{1} \in \mathcal{Z}_{1}
$$

We may assume that $c_{2}^{\prime} \geq c_{2}$. Now, given $U^{0} \in \mathcal{Z}$ arbitrarily, the (constant) solution $U_{1}$ of (12) with the initial value $U_{1}$ (cf. (13)) satisfies the above estimate while the solution $U(t)-U_{1}$ of (12) with the initial value $U^{0}-U_{1}$ satisfies the estimate (24) by assumption. Hence, using the triangle inequality and then the orthogonality of the eigenspaces, we have

$$
\begin{gathered}
\int_{I_{2}} q(U(t))^{2} \mathrm{dt} \leq 2 \int_{I_{2}} q\left(U_{1}\right)^{2} \mathrm{dt}+2 \int_{I_{2}} q\left(U(t)-U_{1}\right)^{2} \mathrm{dt} \\
\leq 2 c_{2}^{\prime}\left\|U_{1}\right\|_{\mathcal{H}}^{2}+2 c_{2}\left\|U^{0}-U_{1}\right\|_{\mathcal{H}}^{2} \leq 2 c_{2}^{\prime}\left\|U^{0}\right\|_{\mathcal{H}}^{2} .
\end{gathered}
$$

Therefore (24) is satisfied for all $U^{0} \in \mathcal{Z}$ if we replace $c_{2}$ by $2 c_{2}^{\prime}$.
If $I$ is a translate of $I_{2}$, say $I=I_{2}+\tau$, then (27) is true with $c_{4}=c_{2}$. Indeed, set

$$
V(t):=U(t+\tau), \quad t \in \mathbb{R},
$$

then $V$ is the solution of (12) with $U^{0}$ replaced by $U(\tau)$. Using (24) and (14) we have

$$
\int_{I} q(U(t))^{2} \mathrm{dt}=\int_{I_{2}} q(V(t))^{2} \mathrm{dt} \leq c_{2}\|V(0)\|_{\mathcal{H}}^{2}=c_{2}\|U(\tau)\|_{\mathcal{H}}^{2}=c_{2}\left\|U^{0}\right\|_{\mathcal{H}}^{2} .
$$

In the general case let us cover $I$ with a finite number of translates of $I_{2}$, say

$$
I \subset \cup_{k=1}^{m}\left(I_{2}+\tau_{k}\right)
$$

Then we have

$$
\int_{I} q(U(t))^{2} \mathrm{dt} \leq \sum_{k=1}^{m} \int_{I_{2}+\tau_{k}} q(U(t))^{2} \mathrm{dt} \leq \sum_{k=1}^{m} c_{2}\left\|U^{0}\right\|_{\mathcal{H}}^{2}=m c_{2}\left\|U^{0}\right\|_{\mathcal{H}}^{2}
$$

i.e. (27) is satisfied with $c_{4}=m c_{2}$.

Step 3. - Now we establish the estimate (26). Since $\mathcal{Z}_{1}$ is finitedimensional, by assumption (21) there exists a positive constant $c$ satisfying

$$
\left\|U_{1}\right\|_{\mathcal{H}}^{2} \leq c \int_{I} q\left(U_{1}\right)^{2} \operatorname{dt}\left(=c|I| q\left(U_{1}\right)^{2}\right), \quad \forall U_{1} \in \mathcal{Z}_{1}
$$

Using this inequality we have for any given $U^{0} \in \mathcal{Z}$ the following estimate :

$$
\begin{aligned}
& \left\|U^{0}\right\|_{\mathcal{H}}^{2}=\left\|U_{1}\right\|_{\mathcal{H}}^{2}+\left\|U^{0}-U_{1}\right\|_{\mathcal{H}}^{2} \leq c \int_{I} q\left(U_{1}\right)^{2} \mathrm{dt}+\left\|U^{0}-U_{1}\right\|_{\mathcal{H}}^{2} \\
& \quad \leq 2 c \int_{I} q(U(t))^{2} \mathrm{dt}+2 c \int_{I} q\left(U(t)-U_{1}\right)^{2} \mathrm{dt}+\left\|U^{0}-U_{1}\right\|_{\mathcal{H}}^{2}
\end{aligned}
$$

Using (22), (27) and (39) to majorize the right-hand side of this estimate we obtain (26) :

$$
\begin{aligned}
\left\|U^{0}\right\|_{\mathcal{H}}^{2} & \leq 2 c c_{0}^{2} \int_{I} p(U(t))^{2} \mathrm{dt}+\left(1+2 c c_{2}\right)\left\|U^{0}-U_{1}\right\|_{\mathcal{H}}^{2} \\
& \leq\left(2 c c_{0}^{2}+\left(1+2 c c_{4}\right) / c^{\prime}\right) \int_{I} p(U(t))^{2} \mathrm{dt}
\end{aligned}
$$

### 5.4. Proof of theorem 5.5

Let us fix an interval $I_{1}$ containing the closure of $I_{1}^{\prime}$ in its interior and such that $\left|I_{1}^{\prime}\right|<\left|I_{1}\right|$. In view of theorem 5.3 it is sufficient to prove that conditions (23), (24) of theorem 5.3 are satisfied with a suitable integer $k$ and with suitable constants $c_{1}, c_{2}$. Let us fix an integer $k \geq l$, to be chosen later.

Condition (24) follows easily from (33). Indeed, given $U^{0} \in \mathcal{Z}$ satisfying (25), let us denote by $U_{+}^{0}$ (resp. by $U_{-}^{0}$ ) its orthogonal projection onto $\mathcal{Z}_{+}$ (resp. onto $\mathcal{Z}_{-}$). Then $U_{+}^{0}$ and $U_{-}^{0}$ satisfy (34) and (35). Observing that the corresponding solutions $U_{+}, U_{-}$of (12) satisfy $U=U_{+}+U_{-}$and applying inequality (33) for $U_{+}, U_{-}$we obtain easily the estimate (24) with $c_{2}=2 c_{2}^{\prime}$ :

$$
\begin{gathered}
\int_{I_{2}} q(U(t))^{2} \mathrm{dt} \leq 2 \int_{I_{2}} q\left(U_{+}(t)\right)^{2} \mathrm{dt}+2 \int_{I_{2}} q\left(U_{-}(t)\right)^{2} \mathrm{dt} \\
\leq 2 c_{2}^{\prime}\left(\left\|U_{+}^{0}\right\|_{\mathcal{H}}^{2}+\left\|U_{-}^{0}\right\|_{\mathcal{H}}^{2}\right)=2 c_{2}^{\prime}\left\|U^{0}\right\|_{\mathcal{H}}^{2} .
\end{gathered}
$$

Turning to the proof of (23) let us choose, following Lebeau [1], an even function $\varphi: \mathbb{R} \rightarrow \mathbb{R}$ of class $C^{\infty}$, satisfying the following conditions :

$$
\begin{equation*}
0 \leq \varphi \leq 1 \quad \text { in } \quad \mathbb{R}, \varphi=1 \quad \text { in } \quad I_{1}^{\prime}, \varphi=0 \quad \text { in } \quad \mathbb{R} \backslash I_{1} \tag{48}
\end{equation*}
$$

Then $\varphi \in \mathcal{S}(\mathbb{R})$ and its Fourier transform

$$
\hat{\varphi}(x)=\int_{-\infty}^{+\infty} \varphi(t) e^{i x t} \mathrm{dt}, x \in \mathbb{R}
$$

is also an even function belonging to $\mathcal{S}(\mathbb{R})$. Hence there is a constant $c_{5}>0$ such that

$$
\begin{equation*}
|\hat{\varphi}(x)|=|\hat{\varphi}(-x)| \leq c_{5}|x|^{-2 \alpha-d}, \forall x \in \mathbb{R} \backslash\{0\} \tag{49}
\end{equation*}
$$

Now given $U^{0} \in \mathcal{Z}$ satisfying (25), let us introduce $U_{+}, U_{-}, U_{+}(t)$ and $U_{-}(t)$ as above. Using (48), (32) and (29) we have

$$
\begin{gathered}
\int_{I_{1}} p(U(t))^{2} \mathrm{dt} \geq \int_{-\infty}^{+\infty} \varphi(t) p(U(t))^{2} \mathrm{dt} \\
=\int_{-\infty}^{+\infty} \varphi(t) p\left(U_{+}(t)\right)^{2} \mathrm{dt}+\int_{-\infty}^{+\infty} \varphi(t) p\left(U_{-}(t)\right)^{2} \mathrm{dt} \\
+\int_{-\infty}^{+\infty} \varphi(t) p\left(U_{+}(t), U_{-}(t)\right) \mathrm{dt}+\int_{-\infty}^{+\infty} \varphi(t) p\left(U_{-}(t), U_{+}(t)\right) \mathrm{dt} \\
\geq c_{1}^{\prime}\left\|U_{+}^{0}\right\|_{\mathcal{H}}^{2}+c_{1}^{\prime}\left\|U_{-}^{0}\right\|_{\mathcal{H}}^{2}+\sum_{\omega_{j} \geq 0} \sum_{\omega_{i}<0} \hat{\varphi}\left(\omega_{j}-\omega_{i}\right)\left(p\left(U_{j}, U_{i}\right)+p\left(U_{i}, U_{j}\right)\right) \\
\geq c_{1}^{\prime}\left\|U^{0}\right\|_{\mathcal{H}}^{2}-\sum_{\omega_{j} \geq 0} \sum_{\omega_{i}<0}\left|\hat{\varphi}\left(\omega_{j}-\omega_{i}\right)\right|\left(p\left(U_{j}\right)^{2}+p\left(U_{i}\right)^{2}\right) \\
\geq c_{1}^{\prime}\left\|U^{0}\right\|_{\mathcal{H}}^{2}-c_{\alpha}^{2} \sum_{\omega_{j} \geq 0} \sum_{\omega_{i}<0}\left|\hat{\varphi}\left(\omega_{j}-\omega_{i}\right)\right|\left(\left|\omega_{j}\right|^{2 \alpha}\left\|U_{j}\right\|_{\mathcal{H}}^{2}+\left|\omega_{i}\right|^{2 \alpha}\left\|U_{i}\right\|_{\mathcal{H}}^{2}\right)
\end{gathered}
$$

yielding

$$
\begin{gather*}
\int_{I_{1}} p(U(t))^{2} \mathrm{dt} \geq \sum_{\omega_{j} \geq 0}\left(c_{1}^{\prime}-c_{\alpha}^{2}\left|\omega_{j}\right|^{2 \alpha} \sum_{\omega_{i}<0}\left|\hat{\varphi}\left(\omega_{j}-\omega_{i}\right)\right|\right)\left\|U_{j}\right\|_{\mathcal{H}}^{2} \\
\quad+\sum_{\omega_{i}<0}\left(c_{1}^{\prime}-c_{\alpha}^{2}\left|\omega_{i}\right|^{2 \alpha} \sum_{\omega_{j} \geq 0}\left|\hat{\varphi}\left(\omega_{j}-\omega_{i}\right)\right|\right)\left\|U_{i}\right\|_{\mathcal{H}}^{2} \tag{50}
\end{gather*}
$$

Using (49) we have

$$
\begin{gathered}
\left|\omega_{j}\right|^{2 \alpha} \sum_{\omega_{i}<0, i \geq k}\left|\hat{\varphi}\left(\omega_{j}-\omega_{i}\right)\right| \\
\leq c_{5}\left|\omega_{j}\right|^{2 \alpha} \sum_{\omega_{i}<0, i \geq k}\left(\left|\omega_{j}\right|+\left|\omega_{i}\right|\right)^{-2 \alpha-d} \leq c_{5} \sum_{i \geq k}\left|\omega_{i}\right|^{-d}
\end{gathered}
$$

whenever $\omega_{j} \geq 0$ and analogously

$$
\left|\omega_{i}\right|^{2 \alpha} \sum_{\omega_{j}>0, i \geq k}\left|\hat{\varphi}\left(\omega_{j}-\omega_{i}\right)\right| \leq c_{5} \sum_{j \geq k}\left|\omega_{j}\right|^{-d}
$$

whenever $\omega_{i}<0$. Substituting these inequalities into (50) we find that

$$
\begin{equation*}
\int_{I_{1}} p(U(t))^{2} \mathrm{dt} \geq\left(c_{1}^{\prime}-c_{\alpha}^{2} c_{5} \sum_{j \geq k}\left|\omega_{j}\right|^{-d}\right)\left\|U^{0}\right\|_{\mathcal{H}}^{2}=: c_{1}\left\|U^{0}\right\|_{\mathcal{H}}^{2} \tag{51}
\end{equation*}
$$

By (30) the coefficient $c_{1}$ in (51) is positive if we choose a sufficiently large integer $k$. Then (23) follows from (51).

## 6. New uniqueness and exact controllability results

The aim of this chapter is to complete and improve the results of chapters 3 and 4 by applying the theorems of the preceding chapter. In particular, we give simple and constructive proofs of certain theorems of Lions [4] and ZuAzUA [1], originally proved by applying compactness arguments combined with Holmgren type unique continuation theorems. Our approach is general and may be applied in cases where Holmgren type theorems are not available, e.g. for equations with non-analytic coefficients. In other cases we only need elliptic unique continuation theorems but not deeper ones concerning evolutionary problems.

### 6.1. A unique continuation theorem

We begin by recalling the following standard result of Carleman [1] :
Theorem 6.1. - Let $G$ be an open domain in $\mathbb{R}^{n}, V \in L^{\infty}(G), u \in H^{2}(G)$, and assume that

$$
\begin{equation*}
-\Delta u+V u=0 \quad \text { in } \quad G . \tag{1}
\end{equation*}
$$

Assume that $u \equiv 0$ in some neighbourhood of a point $x^{\prime} \in G$. Then $u \equiv 0$ in $G$.

See e.g. Garofalo et Lin [1], [2] for a short proof based on the multiplier method.

Corollary 6.2. - Let $\Omega$ be a bounded open domain in $\mathbb{R}^{n}$ with a boundary $\Gamma$ of class $C^{2}$ and let $B$ be an (arbitrarily small) open ball such that

$$
\begin{equation*}
\Gamma \cap B \neq \emptyset . \tag{2}
\end{equation*}
$$

Let $V \in L^{\infty}(\Omega), u \in H^{2}(\Omega)$ and assume that

$$
-\Delta u+V u=0 \quad \text { in } \quad \Omega
$$

and

$$
\begin{equation*}
u=\partial_{\nu} u=0 \quad \text { on } \quad \Gamma \cap B . \tag{3}
\end{equation*}
$$

Then $u \equiv 0$ in $\Omega$.

Proof. - Set $G:=\Omega \cup B$ and define $V(x)=u(x)=0$ for $x \in B \backslash \Omega$. Clearly, we have $V \in L^{\infty}(G)$; it is sufficient to verify that $u \in H^{2}(G)$. Indeed, the we can conclude by applying theorem 6.1 with an arbitrary point $x^{\prime} \in B \backslash \bar{\Omega}$.

Let us denote by $g_{i}, g_{i j}$ the extensions by zero to $G$ of the functions $\partial_{i} u$, $\partial_{i} \partial_{j} u, i, j=1, \ldots, n$. Then $g_{i}, g_{i j} \in L^{2}(G)$ and it suffices the show that

$$
\int_{G} u\left(\partial_{j} \varphi\right) \mathrm{dx}=-\int_{G} g_{j} \varphi \mathrm{dx}, \quad \forall \varphi \in \mathcal{D}(G)
$$

and

$$
\int_{G}\left(\partial_{i} u\right)\left(\partial_{j} \varphi\right) \mathrm{dx}=-\int_{G} g_{i j} \varphi \mathrm{dx}, \quad \forall \varphi \in \mathcal{D}(G)
$$

These are obtained easily by integration by parts and by using the following properties : $\partial_{i} u \equiv 0$ and $\partial_{i} \partial_{j} u \equiv 0$ outside of $\Omega, \varphi \equiv 0$ on $\Gamma \backslash(\Gamma \cap B)(\subset \partial G)$, and $u=\partial_{i} u \equiv 0$ on $\Gamma \cap B$ (by (3)). We have

$$
\begin{aligned}
& \int_{G} u\left(\partial_{j} \varphi\right) \mathrm{dx}=\int_{\Omega} u\left(\partial_{j} \varphi\right) \mathrm{dx}=\int_{\Gamma} u \varphi \nu_{j} \mathrm{~d} \Gamma-\int_{\Omega}\left(\partial_{j} u\right) \varphi \mathrm{dx} \\
&=\int_{\Gamma \cap B} u \varphi \nu_{j} \mathrm{~d} \Gamma-\int_{\Omega}\left(\partial_{j} u\right) \varphi \mathrm{dx}=-\int_{\Omega}\left(\partial_{j} u\right) \varphi \mathrm{dx}=-\int_{G} g_{j} \varphi \mathrm{dx}
\end{aligned}
$$

and

$$
\begin{aligned}
& \int_{G}\left(\partial_{i} u\right)\left(\partial_{j} \varphi\right) \mathrm{dx}=\int_{\Omega}\left(\partial_{i} u\right)\left(\partial_{j} \varphi\right) \mathrm{dx}=\int_{\Gamma}\left(\partial_{i} u\right) \varphi \nu_{j} \mathrm{~d} \Gamma-\int_{\Omega}\left(\partial_{j} \partial_{i} u\right) \varphi \mathrm{dx} \\
= & \int_{\Gamma \cap B}\left(\partial_{i} u\right) \varphi \nu_{j} \mathrm{~d} \Gamma-\int_{\Omega}\left(\partial_{j} \partial_{i} u\right) \varphi \mathrm{dx}=-\int_{\Omega}\left(\partial_{j} \partial_{i} u\right) \varphi \mathrm{dx}=-\int_{G} g_{i j} \varphi \mathrm{dx} .
\end{aligned}
$$

Corollary 6.3. - Let $\Omega$ be a bounded open domain with a boundary $\Gamma$ of class $C^{2}$ and let $V \in L^{\infty}(\Omega)$. Fix $x^{0} \in \mathbb{R}^{n}$ arbitrarily and set

$$
\Gamma_{+}:=\{x \in \Gamma: m(x) \cdot \nu(x)>0\}
$$

as usual. Let $u \in H^{2}(\Omega)$ and assume that

$$
-\Delta u+V u=0 \quad \text { in } \quad \Omega
$$

and

$$
u=\partial_{\nu} u=0 \quad \text { on } \quad \Gamma_{+} .
$$

Then $u \equiv 0$ in $\Omega$.
Proof. - In view of the preceding corollary it is sufficient to find a ball $B$ satisfying $\Gamma \cap B \subset \Gamma_{+}$.

By the compactness of $\Gamma$ there exists a point $x^{1} \in \Gamma$ such that

$$
\left|x^{0}-x^{1}\right|=\max _{\Gamma}\left|x^{0}-x\right| .
$$

It is clear that $m\left(x^{1}\right) \cdot \nu\left(x^{1}\right)>0$. Since $\Omega$ is of class $C^{0}$ (in fact of class $C^{2}$ ), there exists an open ball $B$ of center $x^{1}$ such that $\Gamma \cap B$ is connected. If we choose the radious of $B$ sufficiently small, then we also have $m \cdot \nu>0$ on $\Gamma \cap B$ whence $\Gamma \cap B \subset \Gamma_{+}$.

### 6.2. The wave equation. Dirichlet condition

Let us return to the probem studied in section 3.1. First we shall apply theorem 5.2 in order to obtain a variant of the uniqueness theorem 3.1 concerning the problem

$$
\begin{gather*}
u^{\prime \prime}-\Delta u+q u=0 \quad \text { in } \quad \Omega \times \mathbb{R}  \tag{4}\\
u=0 \quad \text { on } \quad \Gamma \times \mathbb{R}  \tag{5}\\
u(0)=u^{0} \quad \text { and } \quad u^{\prime}(0)=u^{1} . \tag{6}
\end{gather*}
$$

Fix $x^{0} \in \mathbb{R}^{n}$ arbitrarily.
Theorem 6.4. - Let $I$ be an interval of length $|I|>2 R$. There exists a constant $c^{\prime}>0$ such that the solution of (4)-(6) satisfies

$$
\begin{equation*}
\int_{I} \int_{\Gamma_{+}}\left|\partial_{\nu} u\right|^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt} \geq c^{\prime} E, \forall\left(u^{0}, u^{1}\right) \in H_{0}^{1}(\Omega) \times L^{2}(\Omega) \tag{7}
\end{equation*}
$$

If $\Gamma_{+} \neq \Gamma$, then the estimate (7) is slightly weaker than the inequality

$$
\int_{I} \int_{\Gamma}\left|\partial_{\nu} u\right|^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt} \geq c^{\prime} E
$$

of theorem 3.1, because $\mathrm{d} \Gamma_{m} \leq 0$ on $\Gamma_{-}$; nevertheless, this weaker estimate is still sufficient to prove the exact controllability of the corresponding problem : see theorem 6.5 below. On the other hand, and this is the key point, the condition $|I|>2 R$ is weaker than the condition $\left(1-Q_{1}\right)|I|>2 R$ of theorem 3.1.

Proof. - We are going to apply theorem 5.2 with

$$
H=L^{2}(\Omega), \quad V=H_{0}^{1}(\Omega) \text { and } p\left(v^{0}, v^{1}\right):=\left(\int_{\Gamma_{+}}\left|\partial_{\nu} v^{0}\right|^{2} \mathrm{~d} \Gamma_{m}\right)^{1 / 2}
$$

cf. example 5.1. Condition (5.18) of theorem 5.2 is fulfilled by theorem 2.2, regardless of the choice of $k$ and for every interval $I_{2}$, even if $U^{0}$ does not satisfy any orthogonality condition of the type (5.19).

Now fix $l>1$ such that

$$
\begin{equation*}
Q_{l}<1 \quad \text { and } \quad|I|>2 R /\left(1-Q_{l}\right) \tag{8}
\end{equation*}
$$

and set $k=2 l-1$; then we have

$$
\left\{\omega_{1}, \ldots, \omega_{k-1}\right\}=\left\{ \pm \sqrt{\lambda_{1}}, \ldots, \pm \sqrt{\lambda_{l-1}}\right\}
$$

We deduce from theorem 3.1 and remark 3.3 that condition (5.17) is fulfilled for every interval $I_{1}$ of length

$$
\left|I_{1}\right|>2 R /\left(1-Q_{l}\right)
$$

By (8) we may choose $I_{1}$ such that $\left|I_{1}\right|<|I|$.
If $Q_{1}<1$, then we conclude from theorem 3.1 that condition (5.28) (for $q=p$ ) is fulfilled for every interval $I^{\prime}$ of length

$$
\left|I^{\prime}\right|>2 R /\left(1-Q_{1}\right)
$$

This implies (5.16). If $Q_{1} \geq 1$, then we verify (5.16) directly. Let $U=(v, \pm i \sqrt{\lambda} v)$ be an eigenvector of $\mathcal{A}$ satisfying $p(U)=0$. Then we deduce from the definition of $\mathcal{A}$ that

$$
\begin{gathered}
-\Delta v+q v=\lambda v \quad \text { in } \quad \Omega, \\
v=0 \quad \text { on } \quad \Gamma
\end{gathered}
$$

and

$$
\partial_{\nu} v=0 \quad \text { on } \quad \Gamma_{+} .
$$

Applying corollary 6.3 we conclude that $v \equiv 0$ and therefore $U \equiv 0$.
Since the estimate (7) is equivalent to the first inequality in (5.20), we conclude by applying theorem 5.2. [

Theorem 6.4 leads to a strong improvement of theorem 4.1 concerning the exact controllability of the problem

$$
\begin{gather*}
y^{\prime \prime}-\Delta y+q y=0 \quad \text { in } \quad \Omega \times(0, T),  \tag{9}\\
y=v \quad \text { on } \quad \Gamma \times(0, T),  \tag{10}\\
y(0)=y^{0} \quad \text { and } \quad y^{\prime}(0)=y^{1} \tag{11}
\end{gather*}
$$

by eliminating the condition (4.10) on $q$ and by weakening the condition (4.11) on $T$ :

Theorem 6.5. - Let $T>2 R$. Given $\left(y^{0}, y^{1}\right),\left(y_{T}^{0}, y_{T}^{1}\right) \in L^{2}(\Omega) \times H^{-1}(\Omega)$ arbitrarily, there exists $v \in L^{2}\left(0, T ; L^{2}(\Gamma)\right)$ such that

$$
v=0 \quad \text { a.e. on } \quad \Gamma_{-} \times(0, T)
$$

and the solution of (9)-(11) satisfies

$$
y(T)=y_{T}^{0} \quad \text { and } \quad y^{\prime}(T)=y_{T}^{1}
$$

Proof. - We may repeat the proof of theorem 4.1 by applying theorem 6.4 instead of theorem 3.1.

Remark 6.6. - It is easy to give a formal recipe leading to the condition $T>2 R$ in theorem 6.5. Assume for simplicity that $n \geq 2$. Then the corresponding condition (4.11) of theorem 4.1 may be written explicitly as

$$
T>\frac{2 R}{1-2 R Q \lambda_{1}^{-1 / 2}}
$$

Letting $\lambda_{1} \rightarrow+\infty$ we obtain the weaker condition $T>2 R$.

### 6.3. The first Petrovsky system

Now consider the problem

$$
\begin{gather*}
u^{\prime \prime}+\Delta^{2} u=0 \quad \text { in } \quad \Omega \times \mathbb{R}  \tag{12}\\
u=\partial_{\nu} u=0 \quad \text { on } \quad \Gamma \times \mathbb{R}  \tag{13}\\
u(0)=u^{0} \quad \text { and } \quad u^{\prime}(0)=u^{1} \tag{14}
\end{gather*}
$$

We begin by establishing a variant of the uniqueness theorem 3.7 by eliminating hypothesis (3.28) on the length of the interval $I$.

Theorem 6.7. - For every interval I there exists a constant $c^{\prime}>0$ such that

$$
\begin{equation*}
\int_{I} \int_{\Gamma_{+}}|\Delta u|^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt} \geq c^{\prime} E, \forall\left(u^{0}, u^{1}\right) \in H_{0}^{2}(\Omega) \times L^{2}(\Omega) \tag{15}
\end{equation*}
$$

Proof. - Fix an interval $I$ arbitrarily. We are going to apply theorem 5.2 with

$$
H=L^{2}(\Omega), \quad V=H_{0}^{2}(\Omega) \text { and } p\left(v^{0}, v^{1}\right):=\left(\int_{\Gamma_{+}}\left|\Delta v^{0}\right|^{2} \mathrm{~d} \Gamma_{m}\right)^{1 / 2}
$$

(see example 5.1). Since (15) is equivalent to the first inequality in (5.20), it is sufficient to verify the properties (5.17), (5.18) and (5.28) (for $q=p$ ) with a suitable integer $k>1$ and three intervals $I_{1}, I_{2}, I^{\prime}$ such that $\left|I_{1}\right|<|I|$.

By theorem 2.6 condition (5.18) is satisfied with any integer $k \geq 1$ and with any interval $I_{2}$, even without the orthogonality assumption (5.19).

To prove (5.17) and (5.28) choose $l>1$ such that

$$
|I|>2 R \lambda_{l}^{-1 / 4}
$$

and put $k=2 l-1$; then

$$
\left\{\omega_{1}, \ldots, \omega_{k-1}\right\}=\left\{ \pm \sqrt{\lambda_{1}}, \ldots, \pm \sqrt{\lambda_{l-1}}\right\}
$$

It follows from theorem 3.7 and remark 3.13 that (5.17) and (5.28) are satisfied if

$$
\left|I_{1}\right|>2 R \lambda_{l}^{-1 / 4} \quad \text { and } \quad\left|I^{\prime}\right|>2 R \mu_{1}^{-1 / 2}
$$

Thus we may choose $\left|I_{1}\right|$ such that $\left|I_{1}\right|<|I|$ and we may conclude by applying theorem 5.2.

Let us apply this result to the exact controllability problem

$$
\begin{gather*}
y^{\prime \prime}+\Delta^{2} y=0 \quad \text { in } \quad \Omega \times(0, T),  \tag{16}\\
y=0 \quad \text { and } \quad \partial_{\nu} y=v \quad \text { on } \quad \Gamma \times(0, T),  \tag{17}\\
y(0)=y^{0} \quad \text { and } \quad y^{\prime}(0)=y^{1} . \tag{18}
\end{gather*}
$$

Repeating the proof of theorem 4.5 by using theorem 6.7 instead of theorem 3.7 at the end, we obtain the

Theorem 6.8. - Fix $T>0$ arbitrarily (arbitrarily small). Given $\left(y^{0}, y^{1}\right)$, $\left(y_{T}^{0}, y_{T}^{1}\right) \in L^{2}(\Omega) \times H^{-2}(\Omega)$ arbitrarily, there exists $v \in L^{2}\left(0, T ; L^{2}(\Gamma)\right)$ such that

$$
v=0 \quad \text { a.e. on } \quad \Gamma_{-} \times(0, T)
$$

and the solution of (16)-(18) satisfies

$$
y(T)=y_{T}^{0} \quad \text { and } \quad y^{\prime}(T)=y_{T}^{1}
$$

Theorems 6.7 and 6.8 are due to Zuazua [1]. He proved them by using an indirect compactness-uniqueness argument.

Remark 6.9. - We may formally obtain the condition $T>0$ from the stronger condition $T>2 R \lambda_{1}^{-1 / 4}$ of theorem 4.5 and remark 3.13 by letting $\lambda_{1} \rightarrow+\infty$.

### 6.4. The second Petrovsky system. Uniqueness theorems

We consider here the problem (cf. section 3.3)

$$
\begin{gather*}
u^{\prime \prime}+\Delta^{2} u=0 \quad \text { in } \quad \Omega \times \mathbb{R}  \tag{19}\\
u=\Delta u=0 \quad \text { on } \quad \Gamma \times \mathbb{R}  \tag{20}\\
u(0)=u^{0} \quad \text { and } \quad u^{\prime}(0)=u^{1} \tag{21}
\end{gather*}
$$

In the following variant of the uniqueness theorem 3.14 we do not need any hypothesis on the length of the interval $I$. If $\Gamma_{+}=\Gamma$, then the left-hand sides of (22) and (23) below are majorized by the left-hand side of (3.42).

Theorem 6.10. - For every interval I there exists a constant $c^{\prime}>0$ such that

$$
\begin{equation*}
\int_{I} \int_{\Gamma_{+}}\left|\partial_{\nu} u^{\prime}\right|^{2} \mathrm{~d} \Gamma \mathrm{dt} \geq c^{\prime} E_{1 / 4}, \forall\left(u^{0}, u^{1}\right) \in D_{3 / 4} \times D_{1 / 4} \tag{22}
\end{equation*}
$$

and

$$
\begin{equation*}
\int_{I} \int_{\Gamma_{+}}\left|\partial_{\nu} \Delta u\right|^{2} \mathrm{~d} \Gamma \mathrm{dt} \geq c^{\prime} E_{1 / 4}, \forall\left(u^{0}, u^{1}\right) \in D_{3 / 4} \times D_{1 / 4} \tag{23}
\end{equation*}
$$

Proof. - It is sufficient to prove the inequalities (22), (23) for $u^{0}, u^{1} \in Z$; the general case then follows by density, using theorem 2.13.

We are going to apply theorem 5.5 with

$$
\begin{gathered}
H=D_{1 / 4}\left(=H_{0}^{1}(\Omega)\right), \\
V=D_{3 / 4}\left(=\left\{v \in H^{3}(\Omega): v=\Delta v=0 \quad \text { on } \quad \Gamma\right\}\right), \\
p\left(v^{0}, v^{1}\right)=\left\|\partial_{\nu} v^{1}\right\|_{L^{2}\left(\Gamma_{+}\right)} \quad \text { for }(22), \\
p\left(v^{0}, v^{1}\right)=\left\|\partial_{\nu} \Delta v^{0}\right\|_{L^{2}\left(\Gamma_{+}\right)} \quad \text { for }(23)
\end{gathered}
$$

and with $q=p$.
It is clear that (22), (23) are equivalent to (5.36). Therefore it suffices to verify (5.28), (5.29), (5.30), (5.32) and (5.33) with a suitable integer $l>1$ and with suitable intervals $I^{\prime}, I_{1}^{\prime}, I_{2}$ such that $\left|I_{1}^{\prime}\right|<|I|$. (Recall that (5.28) implies (5.21) and note that (5.22) is obvious here.)

It follows from the standard trace theorems that condition (5.29) is fulfilled with any $\alpha>1$.

Condition (5.30) means that in some sense the eigenvalues $\lambda_{k}$ tend to $+\infty$ sufficiently quickly. y a well-known theorem of H . Weyl, see e.g. Agmon [1], in the present case we the estimate

$$
\lambda_{k}=(c+o(1)) k^{4 / n}, \quad k \rightarrow+\infty
$$

holds, which implies (5.30) for every $d>n / 2$.
It follows from theorem 2.13 that inequality (5.33) is fulfilled for all $U^{0} \in \mathcal{Z}$ (even in the absence of conditions (5.34) and (5.35)), for every interval $I_{2}$.

It remains to prove (5.28) and (5.32). These will be deduced from (the apparently weaker) theorem 3.14 and remark 3.17 by using a small "trick".

Observe that (by the special boundary conditions (20)) for $\left(u^{0}, u^{1}\right) \in \mathcal{Z}_{ \pm}$ the solution of (19)-(21) is also solution of the problem

$$
\begin{gather*}
u^{\prime} \pm i \Delta u=0 \quad \text { in } \quad \Omega \times \mathbb{R},  \tag{24}\\
u=0 \quad \text { on } \quad \Gamma \times \mathbb{R},  \tag{25}\\
u(0)=u^{0} . \tag{26}
\end{gather*}
$$

Since the solutions of (24)-(26) clearly satisfy

$$
\begin{equation*}
\left|\partial_{\nu} u^{\prime}\right|=\left|\partial_{\nu} \Delta u\right| \quad \text { on } \quad \Gamma \times \mathbb{R}, \tag{27}
\end{equation*}
$$

(3.42) reduces to

$$
\int_{J} \int_{\Gamma_{+}}\left|\partial_{\nu} u^{\prime}\right|^{2} \mathrm{~d} \Gamma \mathrm{dt}=\int_{J} \int_{\Gamma_{+}}\left|\partial_{\nu} \Delta u\right|^{2} \mathrm{~d} \Gamma \mathrm{dt} \geq \frac{c^{\prime}}{2} E_{1 / 4}
$$

i.e. to

$$
\int_{J} p(U(t))^{2} \mathrm{dt} \geq \frac{c^{\prime}}{2}\left\|U^{0}\right\|_{\mathcal{H}}^{2}
$$

if $|J|>R / \sqrt{\lambda_{1}}$. Hence (5.28) is satisfied if we choose the interval $I^{\prime}$ such that $\left|I^{\prime}\right|>R / \sqrt{\lambda_{1}}$.

The proof of (5.32) is similar. First we choose $k>1$ such that $|I|>R / \sqrt{\lambda_{k}}$ and then we choose an interval $I_{1}^{\prime}$ satisfying $|I|>\left|I_{1}^{\prime}\right|>R / \sqrt{\lambda_{k}}$. Using again (27), (5.32) follows from (3.42) and from remark 3.17 if we choose $l=2 k-1$.

We may apply theorem 5.5 and the proof is completed.
Next we deduce from theorem 6.10 the inverse inequality of the direct inequality obtained in theorem 2.10.

Theorem 6.11. - For every interval $I$ there exists a constant $c^{\prime}>0$ such that the solution of (19)-(21) satisfies

$$
\begin{equation*}
\int_{I} \int_{\Gamma_{+}}\left|\partial_{\nu} u\right|^{2} \mathrm{~d} \Gamma \mathrm{dt} \geq c^{\prime}\left(\left\|u^{0}\right\|_{H_{0}^{1}(\Omega)}^{2}+\left\|u^{1}\right\|_{H^{-1}(\Omega)}^{2}\right) \tag{28}
\end{equation*}
$$

for every $\left(u^{0}, u^{1}\right) \in H_{0}^{1}(\Omega) \times H^{-1}(\Omega)$.
Proof. - (Compare with that of theorem 2.10.) Using a density argument based on theorem 2.10, it is sufficient to prove (28) for $u^{0}, u^{1} \in Z$.

Applying the inequality (23) for $\left(A^{-1 / 2} u^{0}, A^{-1 / 2} u^{1}\right)$ instead of $\left(u^{0}, u^{1}\right)$ we obtain

$$
\begin{gathered}
\int_{I} \int_{\Gamma_{+}}\left|\partial_{\nu} \Delta A^{-1 / 2} u\right|^{2} \mathrm{~d} \Gamma \mathrm{dt} \geq \frac{c^{\prime}}{2}\left(\left\|A^{-1 / 2} u^{0}\right\|_{3 / 4}^{2}+\left\|A^{-1 / 2} u^{1}\right\|_{1 / 4}^{2}\right) \\
\quad=\frac{c^{\prime}}{2}\left(\left\|u^{0}\right\|_{1 / 4}^{2}+\left\|u^{1}\right\|_{-1 / 4}^{2}\right)=\frac{c^{\prime}}{2}\left(\left\|u^{0}\right\|_{H_{0}^{1}(\Omega)}^{2}+\left\|u^{1}\right\|_{H^{-1}(\Omega)}^{2}\right)
\end{gathered}
$$

Since $\Delta A^{-1 / 2} u \equiv-u,(28)$ hence follows (with $c^{\prime} / 2$ instead of $c^{\prime}$ ).

### 6.5. The second Petrovsky system. Exact controllability

Theorems 2.10 and 6.11 permit us to study the exact controllability of the problem

$$
\begin{gather*}
y^{\prime \prime}+\Delta^{2} y=0 \quad \text { in } \quad \Omega \times(0, T),  \tag{29}\\
y=0 \quad \text { and } \quad \Delta y=v \quad \text { on } \quad \Gamma \times(0, T),  \tag{30}\\
y(0)=y^{0}, \quad y^{\prime}(0)=y^{1} . \tag{31}
\end{gather*}
$$

It follows from theorem 2.14 that for any given $\left(y^{0}, y^{1}\right) \in H_{0}^{1}(\Omega) \times H^{-1}(\Omega)$ and $v \in L^{2}\left(0, T ; L^{2}(\Gamma)\right)$ this problem has a unique solution satisfying $\left(y(T), y^{\prime}(T)\right) \in H_{0}^{1}(\Omega) \times H^{-1}(\Omega)$.

Definition. - We say that the problem (29)-(31) is exactly controllable if for any given $\left(y^{0}, y^{1}\right),\left(y_{T}^{0}, y_{T}^{1}\right) \in H_{0}^{1}(\Omega) \times H^{-1}(\Omega)$ there exists $v \in$ $L^{2}\left(0, T ; L^{2}(\Gamma)\right)$ such that the solution of (29)-(31) satisfies $y(T)=y_{T}^{0}$ and $y^{\prime}(T)=y_{T}^{1}$.

Theorem 6.12. - The problem (29)-(31) is exactly controllable for every $T>0$ (arbitrarily small). Moreover, there exist controls $v$ satisfying

$$
\begin{equation*}
v=0 \quad \text { a.e. on } \quad \Gamma_{-} \times(0, T) . \tag{32}
\end{equation*}
$$

Proof. - We apply HUM as in chapter 4. Fixing $u^{0} \in H_{0}^{1}(\Omega)$ and $u^{1} \in H^{-1}$ ( $\Omega$ arbitrarily, we solve (29)-(31), then we solve the problem

$$
\begin{gather*}
y^{\prime \prime}+\Delta^{2} y=0 \quad \text { in } \quad \Omega \times(0, T),  \tag{33}\\
y=0 \quad \text { and } \quad \Delta y=-\partial_{\nu} u \quad \text { on } \quad \Gamma_{+} \times(0, T),  \tag{34}\\
y=\Delta y=0 \quad \text { on } \quad \Gamma_{-} \times(0, T),  \tag{32}\\
y(T)=y^{\prime}(T)=0, \tag{36}
\end{gather*}
$$

and we set

$$
\Lambda\left(u^{0}, u^{1}\right):=\left(y^{\prime}(0),-y(0)\right)
$$

It follows from theorems 2.10 and 2.14 that $\Lambda$ is a bounded linear map from $F:=H_{0}^{1}(\Omega) \times H^{-1}(\Omega)$ into $F^{\prime}=H^{-1}(\Omega) \times H_{0}^{1}(\Omega)$; it suffices to show that it is surjective. Using the Lax-Milgram theorem it is sufficient to prove the inequality

$$
\left\langle\Lambda\left(u^{0}, u^{1}\right),\left(u^{0}, u^{1}\right)\right\rangle_{F^{\prime}, F} \geq c\left\|\left(u^{0}, u^{1}\right)\right\|_{F}^{2}
$$

with a constant $c>0$, independent of $\left(u^{0}, u^{1}\right) \in H_{0}^{1}(\Omega) \times H^{-1}(\Omega)$. By theorem 6.11 it is sufficient to show that

$$
\begin{equation*}
\left\langle\Lambda\left(u^{0}, u^{1}\right),\left(u^{0}, u^{1}\right)\right\rangle_{F^{\prime}, F} \geq c \int_{I} \int_{\Gamma_{+}}\left|\partial_{\nu} u\right|^{2} \mathrm{~d} \Gamma \mathrm{dt} . \tag{37}
\end{equation*}
$$

As usual, we may restrict ourselves to the real case and we may asume that $\left(u^{0}, u^{1}\right) \in Z \times Z$. Multiplying the equation (33) by $u$ and integrating by parts we obtain that

$$
\begin{gathered}
0=\int_{0}^{T} \int_{\Omega} u\left(y^{\prime \prime}+\Delta^{2} y\right) \mathrm{dx} \mathrm{dt}= \\
\int_{0}^{T} \int_{\Omega}\left(u^{\prime \prime}+\Delta^{2} u\right) y \mathrm{dx} \mathrm{dt}+\left[\int_{\Omega} u y^{\prime}-u^{\prime} y \mathrm{dx}\right]_{0}^{T} \\
+\int_{0}^{T} \int_{\Gamma} u\left(\partial_{\nu} \Delta y\right)-\left(\partial_{\nu} u\right)(\Delta y)+(\Delta u)\left(\partial_{\nu} y\right)-\left(\partial_{\nu} \Delta u\right) y \mathrm{~d} \Gamma \mathrm{dt} \\
=\int_{\Omega} u^{1} y(0)-u^{0} y^{\prime}(0) \mathrm{dx}+\int_{0}^{T} \int_{\Gamma_{+}}\left(\partial_{\nu} u\right)^{2} \mathrm{~d} \Gamma \mathrm{dt} .
\end{gathered}
$$

Hence (37) follows.
Remark 6.13. - Improving a former result of Lions [5; p. 310], Zuazua [1] earlier obtained (with an indirect method) a theorem analogous to theorem 6.12, by using two controls. See also Lasiecka [1] for another result analogous to theorem 6.12.

Now it is easy to study the exact controllability of the problem by acting in the other boundary condition :

$$
\begin{gather*}
z^{\prime \prime}+\Delta^{2} z=0 \quad \text { in } \quad \Omega \times(0, T),  \tag{38}\\
z=v \quad \text { et } \Delta z=0 \quad \text { on } \quad \Gamma \times(0, T),  \tag{39}\\
z(0)=z^{0} \quad \text { and } \quad z^{\prime}(0)=z^{1} . \tag{40}
\end{gather*}
$$

We have the following result (we recall from lemma 1.7 that $D_{-1 / 4}=$ $\left.H^{-1}(\Omega) \subset D_{-3 / 4}\right)$.
Theorem 6.14. - Fix $T>0$ arbitrarily. Given

$$
\left(z^{0}, z^{1}\right),\left(z_{T}^{0}, z_{T}^{1}\right) \in D_{-1 / 4} \times D_{-3 / 4}
$$

arbitrarily, there exists $w \in L^{2}\left(0, T ; D_{-1 / 4}\right)$ such that

$$
w(t)=0 \quad \text { on } \quad \Gamma_{-} \quad \text { for almost every } \quad t \in(0, T)
$$

and the solution of (38)-(40) satisfies $z(T)=z_{T}^{0}, z^{\prime}(T)=z_{T}^{1}$.
Remark 6.15. - We define the solution of (38)-(40) in the following way : we solve (29)-(32) with $\left(y^{0}, y^{1}\right):=\left(A^{-1 / 2} z^{0}, A^{-1 / 2} z^{1}\right) \in D_{1 / 4} \times D_{-1 / 4}=$ $H_{0}^{1}(\Omega) \times H^{-1}(\Omega)$ and $w:=-v$, and then we set $z:=A^{1 / 2} y$. To justify this definition we observe that if $y$ is a sufficiently smooth solution of (29)-(32), then $z:=A^{1 / 2} y$ is given in fact by $z=-\Delta y$ and therefore (38)-(40) is satisfied in the usual sense :

$$
\begin{gathered}
z^{\prime \prime}+\Delta^{2} z=-\Delta\left(y^{\prime \prime}+\Delta y\right)=0 \quad \text { in } \quad \Omega \times(0, T), \\
z=-\Delta y=-v=w \quad \text { and } \quad \Delta z=-\Delta^{2} y=y^{\prime \prime}=0 \quad \text { on } \quad \Gamma \times(0, T), \\
z(0)=-\Delta y^{0}=-\Delta A^{1 / 2} z^{0}=z^{0} \\
z^{\prime}(0)=-\Delta y^{\prime}(0)=-\Delta A^{1 / 2} z^{1}=z^{1}
\end{gathered}
$$

Proof of theorem 6.14. - The theorem follows at once from theorem 6.13 and from remark 6.15.

### 6.6. The wave equation. Neumann or Robin condition

Let us return to the problem (cf. $\S \S 3.4$ and 4.3)

$$
\begin{gather*}
u^{\prime \prime}-\Delta u+q u=0 \quad \text { in } \quad \Omega \times \mathbb{R},  \tag{41}\\
u=0 \quad \text { on } \quad \Gamma_{0} \times \mathbb{R},  \tag{42}\\
\partial_{\nu} u+a u=0 \quad \text { on } \quad \Gamma_{1} \times \mathbb{R}  \tag{43}\\
u(0)=u^{0} \quad \text { and } \quad u^{\prime}(0)=u^{1} \tag{44}
\end{gather*}
$$

Assume that

$$
\begin{equation*}
m \cdot \nu \leq 0 \quad \text { on } \quad \Gamma_{0}, \quad m \cdot \nu \geq 0 \quad \text { on } \quad \Gamma_{1}, \tag{45}
\end{equation*}
$$

and that $a$ has the form

$$
\begin{equation*}
a=(m \cdot \nu) b, \quad b \in C^{1}\left(\Gamma_{1}\right), \quad b \geq 0 \text { on } \Gamma_{1} . \tag{46}
\end{equation*}
$$

First we improve proposition 3.21 by eliminating hypothesis (3.60) on $q$ and by weakening hypothesis (3.61) on the length of $I$ :

Proposition 6.16. - Assume (45), (46) and let I be an interval of length $>2 R$. Then there exists a constant $c^{\prime}>0$ such that

$$
\begin{equation*}
\int_{I} \int_{\Gamma_{1}}|u|^{2}+\left|u^{\prime}\right|^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt} \geq c^{\prime} E, \forall\left(u^{0}, u^{1}\right) \in D_{1} \times D_{1 / 2} \tag{47}
\end{equation*}
$$

Proof. - It suffices to show (47) for $\left(u^{0}, u^{1}\right) \in Z \times Z$; the general case then follows by density. We are going to apply theorem 5.3 with $V$ and $H$ defined as in section 1.4 (cf. example 5.1) and with

$$
p\left(v^{0}, v^{1}\right):=\left(\int_{\Gamma_{1}}\left|v^{0}\right|^{2}+\left|v^{1}\right|^{2} \mathrm{~d} \Gamma_{m}\right)^{1 / 2}, \quad q\left(v^{0}, v^{1}\right):=\left(\int_{\Gamma_{1}}\left|v^{0}\right|^{2} \mathrm{~d} \Gamma_{m}\right)^{1 / 2}
$$

Then (47) is equivalent to (5.26), and (5.22), (5.24) are obviously satisfied.
To prove (5.23) fix $l>1$ such that

$$
Q_{l}<1,
$$

and

$$
|I|>2 R_{l} /\left(1-Q_{l}\right)
$$

(this is possible because the right-hand side tends to $2 R$ as $l \rightarrow+\infty$, see remark 3.24), choose an interval $I_{1}$ such that

$$
|I|>\left|I_{1}\right|>2 R_{l} /\left(1-Q_{l}\right)
$$

Then putting $k=2 l-1$ property (5.23) follows from proposition 3.21 and remark 3.24.

It remains to verify (5.21). If $Q_{1}<1$, then we can show easily that the stronger condition (5.28) is also satisfied for every interval $I^{\prime}$ of length

$$
\left|I^{\prime}\right|>2 R_{1} /\left(1-Q_{1}\right)
$$

Indeed, given $j \geq 1$ and $U^{0} \in \mathcal{Z}_{j}$ arbitrarily, the solution of (5.12) satisfies

$$
p(U(t))^{2}=\left(1+\omega_{j}^{2}\right) q(U(t))^{2}
$$

using proposition 3.21 hence (5.28) follows.
If $Q_{1} \geq 1$, then we verify (5.21) directly. It is sufficient to show that if an eigenfunction $v \in Z$ satisfies

$$
-\Delta v+q v=\lambda v \quad \text { in } \quad \Omega
$$

$$
v=0 \quad \text { on } \quad \Gamma_{0},
$$

and

$$
v=\partial_{\nu} v+a v=0 \quad \text { on } \quad \Gamma_{1}
$$

with some real number $\lambda$, then in fact $v \equiv 0$. And this follows from corollary 6.3.

We are going to improve proposition 6.16 by eliminating the term $|u|^{2}$ from (47). We need a lemma. Fix an interval $I$ of length $>2 R$ and denote by $X=X(I)$ the completion of the vector space of the solutions of (41)-(44) corresponding to the initial data $\left(u^{0}, u^{1}\right) \in D_{1} \times D_{1 / 2}$, with respect to the norm

$$
u \mapsto\left(\int_{I} \int_{\Gamma_{1}}|u|^{2}+\left|u^{\prime}\right|^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt}\right)^{1 / 2} .
$$

By the preceding proposition $X$ is a Hilbert space and we have

$$
\begin{equation*}
\int_{I} \int_{\Gamma_{1}}|u|^{2}+\left|u^{\prime}\right|^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt} \geq c^{\prime} E, \quad \forall u \in X \tag{48}
\end{equation*}
$$

by density.
Lemma 6.17. - Let $|I|$ be an interval of length $>2 R$ and let $u \in X=X(I)$ be such that

$$
\begin{equation*}
u^{\prime}=0 \quad \text { on } \quad \Gamma_{1} \times I . \tag{49}
\end{equation*}
$$

Then

$$
\begin{equation*}
u \equiv 0 \tag{50}
\end{equation*}
$$

Proof. - Let us introduce the function

$$
U(t):=\int_{0}^{t} u^{\prime}(s) \mathrm{ds}, \quad t \in \mathbb{R}
$$

Then $U \in X$, and applying (48) for $U$ we obtain that $U \equiv 0$. Hence $u^{\prime} \equiv 0$ and therefore $u$ does not depend on $t \in \mathbb{R}: u(t)=u^{0}$ for all $t \in \mathbb{R}$. It remains to show that $u^{0} \equiv 0$. Since $u^{\prime} \equiv 0$, we deduce from (41)-(43) that $u^{0}$ belongs to the kernel of the operator $A$. Since $A$ is injective, the lemma follows. $\square$.

Proposition 6.18. - Assume (45), (46) and let I be an interval of length $>2 R$. Then there exists a constant $c^{\prime}>0$ such that

$$
\begin{equation*}
\int_{I} \int_{\Gamma_{1}}\left|u^{\prime}\right|^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt} \geq c^{\prime} E, \forall\left(u^{0}, u^{1}\right) \in D_{1} \times D_{1 / 2} \tag{52}
\end{equation*}
$$

Proof. - By proposition 6.16 it is sufficient to show the existence of a constant $c$ such that

$$
\int_{I} \int_{\Gamma_{1}}|u|^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt} \leq c \int_{I} \int_{\Gamma_{1}}\left|u^{\prime}\right|^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt}, \forall\left(u^{0}, u^{1}\right) \in D_{1} \times D_{1 / 2}
$$

Assume on the contrary that there is a sequence $\left(u_{n}^{0}, u_{n}^{1}\right)$ in $D_{1} \times D_{1 / 2}$ such that the corresponding solutions satisfy

$$
\begin{equation*}
\int_{I} \int_{\Gamma_{1}}\left|u_{n}\right|^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt}=1, \forall n \tag{53}
\end{equation*}
$$

and

$$
\begin{equation*}
\int_{I} \int_{\Gamma_{1}}\left|u_{n}^{\prime}\right|^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt} \rightarrow 0 \tag{54}
\end{equation*}
$$

Then the sequence $\left(u_{n}\right)$ is bounded in $X$. Extracting a subsequence if needed, we may assume that

$$
u_{n} \rightharpoonup u \quad \text { (weakly) in } \quad X
$$

for some $u \in X$. By (54) we have

$$
u^{\prime} \equiv 0 \quad \text { on } \quad \Gamma_{1} \times I
$$

therefore $u \equiv 0$ by lemma 6.17 .
On the other hand, since $\left(u_{n}\right)$ is bounded in $X$, we deduce from (48) that it is bounded in

$$
L^{\infty}\left(\mathbb{R} ; H^{1}(\Omega)\right) \cap W^{1, \infty}\left(\mathbb{R} ; L^{2}(\Omega)\right)
$$

and hence also in

$$
L^{2}\left(I ; H^{1}(\Omega)\right) \cap H^{1}\left(I ; L^{2}(\Omega)\right)
$$

Consequently, using an Ascoli type result (see Lions [1; theorem 1.5.1]) the sequence is precompact in $L^{2}\left(I ; H^{1-\varepsilon}(\Omega)\right)$ for all $\varepsilon>0$ and therefore its trace is precompact in $L^{2}(\Gamma \times I)$. Thus we deduce from (53) that

$$
\int_{I} \int_{\Gamma_{1}}|u|^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt}=1
$$

contradicting $u \equiv 0$.
Now we may improve theorem 3.26 by eliminating the hypotheses on $q, b$, and by weakening the hypothesis on the length of $I$ :

Theorem 6.19. - Assume (45), (46) and let I be an interval of length $>2 R$. Then there exists a constant $c^{\prime}>0$ such that

$$
\int_{I} \int_{\Gamma_{1}}|u|^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt} \geq c^{\prime} E_{-1 / 2}, \forall\left(u^{0}, u^{1}\right) \in V \times H
$$

Proof. - We repeat the proof of theorem 3.26, using proposition 6.18 instead of the estimate (3.68) in remark 3.23.

Finally, we can strongly improve theorem 4.8 on the exact controllability of the problem

$$
\begin{gather*}
y^{\prime \prime}-\Delta y+q y=0 \quad \text { in } \Omega \times(0, T),  \tag{55}\\
y=0 \quad \text { on } \quad \Gamma_{0} \times(0, T),  \tag{56}\\
\partial_{\nu} y+a y=v \quad \text { on } \quad \Gamma_{1} \times(0, T),  \tag{57}\\
y(0)=y^{0} \quad \text { and } \quad y^{\prime}(0)=y^{1} \quad \text { on } \quad \Omega . \tag{58}
\end{gather*}
$$

Theorem 6.20. - Assume (45), (46) and let $T>2 R$. Then, given

$$
\left(y^{0}, y^{1}\right),\left(y_{T}^{0}, y_{T}^{1}\right) \in V \times H
$$

arbitrarily, there exists a control $v \in L^{2}\left(0, T ; L^{2}\left(\Gamma_{1}\right)\right)$ such that the solution of (55)-(58) satisfies

$$
y(T)=y_{T}^{0} \quad \text { and } \quad y^{\prime}(T)=y_{T}^{1} \quad \text { on } \quad \Omega .
$$

For $q \equiv 0$ this theorem is due to Lions [5; p. 203]. He used unique continuation theorems for the wave equation. Let us note that the above proof does not use any unique continuation theorem if $q \equiv 0$.

Proof. - We repeat the proof of theorem 4.8, using the preceding theorem instead of theorem 3.26.

## 7. Dissipative evolutionary problems

In this chapter we introduce some dissipative systems governed by the wave equation or by a plate equation; their stabilization properties will be studied in the following chapters. In order to simplify the notation we limit ourselves to the real case.

### 7.1. Maximal monotone operators

In this section we recall some existence and perturbation results concerning the evolutionary problem

$$
\begin{equation*}
U^{\prime}+\mathcal{A} U=0 \quad \text { in } \quad \mathbb{R}_{+}:=[0,+\infty), \quad U(0)=U^{0} \tag{1}
\end{equation*}
$$

where $\mathcal{A}: D(\mathcal{A}) \subset \mathcal{H} \rightarrow \mathcal{H}$ is an operator, non necessarily linear, in a real Hilbert space $\mathcal{H}$. We refer e.g. to Barbu [1] or Brézis [1] for proof.

We say that $\mathcal{A}$ is maximal monotone if the following two properties are fulfilled :

$$
\begin{gather*}
\mathcal{A} \text { is monotone: } \quad(\mathcal{A} U-\mathcal{A} V, U-V)_{\mathcal{H}} \geq 0, \quad \forall U, V \in D(\mathcal{A}),  \tag{2}\\
I+\mathcal{A}  \tag{3}\\
\text { is surjective }: \quad R(I+\mathcal{A})=\mathcal{H} .
\end{gather*}
$$

Theorem 7.1. - Let $\mathcal{A}$ be a maximal monotone operator in a Hilbert space $\mathcal{H}$. Then for every $U^{0} \in \overline{D(\mathcal{A})}$ the problem (1) has a unique solution

$$
\begin{equation*}
U \in C\left(\mathbb{R}_{+} ; \mathcal{H}\right) \tag{4}
\end{equation*}
$$

(defined in some suitable sense). If $V^{0} \in \overline{D(\mathcal{A})}$ and if $V$ is the corresponding solution of (1), then the function

$$
\begin{equation*}
t \mapsto\|U(t)-V(t)\|_{\mathcal{H}} \quad \text { is non-increasing in } \quad \mathbb{R}_{+} . \tag{5}
\end{equation*}
$$

If $U^{0} \in D(\mathcal{A})$, then the solution is more regular :

$$
\begin{equation*}
U \in W^{1, \infty}\left(\mathbb{R}_{+} ; \mathcal{H}\right) \tag{6}
\end{equation*}
$$

and the function

$$
\begin{equation*}
t \mapsto\|\mathcal{A} U(t)\|_{\mathcal{H}} \quad \text { is defined everywhere and is non-increasing in } \mathbb{R}_{+} . \tag{7}
\end{equation*}
$$

Remark 7.2. - One can show that $\overline{D(\mathcal{A})}$ is always convex. Moreover, if $\mathcal{A}$ is linear, then $\overline{D(\mathcal{A})}=\mathcal{H}$. In the linear case the conclusion of the theorem is also stronger : for $U^{0} \in D(\mathcal{A})$ the solution satisfies

$$
U \in C^{1}\left(\mathbb{R}_{+} ; \mathcal{H}\right)
$$

instead of (6). We note that theorem 7.1 also remains valid for multivalued maximal monotone operators.

Let us also recall the following perturbation theorem (see e.g. Barbu [1; proposition 4.2.1 and theorem 4.2.1]).

Theorem 7.3. - Let $\mathcal{A}$ and $\mathcal{A}_{k}(k=1,2, \ldots)$ be maximal monotone operators in a Hilbert space $\mathcal{H}$ and assume that

$$
\left(I+\mathcal{A}_{k}\right)^{-1} W \rightarrow(I+\mathcal{A})^{-1} W \quad \text { in } \quad \mathcal{H}
$$

for every $W \in \mathcal{H}$, as $k \rightarrow+\infty$. Choose $U^{0}, U_{k}^{0} \in \mathcal{H}$ such that

$$
U_{k}^{0} \rightarrow U^{0} \quad \text { in } \quad \mathcal{H}
$$

Then the corresponding solutions of (1) and

$$
U_{k}^{\prime}+\mathcal{A}_{k} U_{k}=0 \quad \text { in } \quad \mathbb{R}_{+}, \quad U_{k}(0)=U_{k}^{0}
$$

satisfy

$$
U_{k}(t) \rightarrow U(t) \quad \text { in } \quad \mathcal{H}
$$

for every $t \in \mathbb{R}_{+}$.

### 7.2. The wave equation

Let $\Omega$ be a bounded domain of class $C^{2}$ in $\mathbb{R}^{n}$ and let $\left\{\Gamma_{0}, \Gamma_{1}\right\}$ be a partition of its boundary $\Gamma$. Fix three nonnegative functions $q: \Omega \rightarrow \mathbb{R}$, $a, l: \Gamma_{1} \rightarrow \mathbb{R}$ and let $g: \mathbb{R} \rightarrow \mathbb{R}$ be a continuous, non-decreasing function such that $g(0)=0$. Consider the problem

$$
\begin{gather*}
u^{\prime \prime}-\Delta u+q u=0 \quad \text { in } \quad \Omega \times \mathbb{R}_{+},  \tag{8}\\
u=0 \quad \text { on } \Gamma_{0} \times \mathbb{R}_{+},  \tag{9}\\
\partial_{\nu} u+a u+l g\left(u^{\prime}\right)=0 \quad \text { on } \quad \Gamma_{1} \times \mathbb{R}_{+},  \tag{10}\\
u(0)=u^{0} \quad \text { and } \quad u^{\prime}(0)=u^{1} \quad \text { on } \Omega . \tag{11}
\end{gather*}
$$

The particular case $l \equiv 0$ corresponds to the conservative system of section 1.3.

Assume that

$$
\begin{equation*}
q \in L^{\infty}(\Omega), \quad a, l \in C^{1}\left(\Gamma_{1}\right) \tag{12}
\end{equation*}
$$

Furthermore, in order to avoid some extra difficulties we shall always assume that

$$
\begin{gather*}
\overline{\Gamma_{0}} \cap \overline{\Gamma_{1}}=\emptyset,  \tag{13}\\
\Gamma_{0} \neq \emptyset \text { or } q \not \equiv 0 \text { or } a \not \equiv 0,  \tag{14}\\
n \geq 3, \tag{15}
\end{gather*}
$$

and that there exists a constant $c>0$ such that

$$
\begin{equation*}
|g(x)| \leq 1+c|x|^{n /(n-2)}, \quad \forall x \in \mathbb{R} \tag{16}
\end{equation*}
$$

(See Komornik and Zuazua [1], Baruce and Hanouzet [1], Komornik [10] and Tcheugoue [1] for the study of the problem if one of conditions (13)-(16) is not satisfied.)

Set

$$
\begin{gathered}
V=H_{\Gamma_{0}}^{1}(\Omega):=\left\{v \in H^{1}(\Omega): v=0 \text { on } \Gamma_{0}\right\} \\
\|v\|_{V}^{2}:=\int_{\Omega}|\nabla v|^{2}+q v^{2} \mathrm{dx}+\int_{\Gamma_{1}} a v^{2} \mathrm{~d} \Gamma .
\end{gathered}
$$

By hypothesis (14) the last expression defines a norm on $V$, which is equivalent to the norm induced by $H^{1}(\Omega)$; consequently, $V=H_{\Gamma_{0}}^{1}(\Omega)$ is a Hilbert space.

We shall prove the
Theorem 7.4. - Assume (12)-(16). Given $\left(u^{0}, u^{1}\right) \in H_{\Gamma_{0}}^{1}(\Omega) \times L^{2}(\Omega)$ arbitrarily, the problem (8)-(11) has a unique solution satisfying

$$
\begin{equation*}
u \in C\left(\mathbb{R}_{+} ; H_{\Gamma_{0}}^{1}(\Omega)\right) \cap C^{1}\left(\mathbb{R}_{+} ; L^{2}(\Omega)\right) \tag{17}
\end{equation*}
$$

The energy $E: \mathbb{R}_{+} \rightarrow \mathbb{R}_{+}$of the solution, defined by

$$
\begin{equation*}
E=E(u):=\frac{1}{2} \int_{\Omega}\left|u^{\prime}\right|^{2}+|\nabla u|^{2}+q u^{2} \mathrm{dx}+\frac{1}{2} \int_{\Gamma_{1}} a u^{2} \mathrm{~d} \Gamma \tag{18}
\end{equation*}
$$

is non-increasing. Moreover, if $u$ and $v$ are two solutions (corresponding to different initial data), then the function $E(u-v)$ is non-increasing; in particular,

$$
\begin{equation*}
\|E(u-v)\|_{L^{\infty}\left(\mathbb{R}_{+}\right)} \leq E(u-v)(0) \tag{19}
\end{equation*}
$$

Under a stronger growth assumption on $g$ we have a simple characterization of $D(\mathcal{A})$ :

Theorem 7.5. - Assume (12)-(15) and assume that $g$ is globally Lipschitz continuous : there exists a constant $c^{\prime}>0$ such that

$$
\begin{equation*}
\left|g\left(x_{1}\right)-g\left(x_{2}\right)\right| \leq c^{\prime} \mid x_{1}-x_{2}, \quad \forall x_{1}, x_{2} \in \mathbb{R} \tag{20}
\end{equation*}
$$

Then $D(\mathcal{A})$ consists of the couples

$$
\begin{equation*}
\left(u^{0}, u^{1}\right) \in\left(H^{2}(\Omega) \cap H_{\Gamma_{0}}^{1}(\Omega)\right) \times H_{\Gamma_{0}}^{1}(\Omega) \tag{21}
\end{equation*}
$$

which satisfy

$$
\begin{equation*}
\partial_{\nu} u^{0}+a u^{0}+\lg \left(u^{1}\right)=0 \quad \text { on } \quad \Gamma_{1} . \tag{22}
\end{equation*}
$$

Hence for $\left(u^{0}, u^{1}\right) \in D(\mathcal{A})$ the solution of (8)-(11) satisfies

$$
\begin{equation*}
u \in L^{\infty}\left(\mathbb{R}_{+} ; H^{2}(\Omega)\right), \quad u^{\prime} \in L^{\infty}\left(\mathbb{R}_{+} ; H_{\Gamma_{0}}^{1}(\Omega)\right) \tag{23}
\end{equation*}
$$

and

$$
\begin{equation*}
u^{\prime \prime} \in L^{\infty}\left(\mathbb{R}_{+} ; L^{2}(\Omega)\right) . \tag{24}
\end{equation*}
$$

The rest of this section is devoted to the proof of these theorems. Set

$$
H:=L^{2}(\Omega), \quad\|v\|_{H}^{2}:=\int_{\Omega} v^{2} \mathrm{dx}
$$

then $V$ and $H$ are separable Hilbert spaces with a dense and compact inclusion $V \subset H$. As usual, we introduce the duality mapping $A: V \rightarrow V^{\prime}$ and we identify $H$ with $H^{\prime}$.

By hypothesis (16) the formula

$$
\langle B z, v\rangle_{V^{\prime}, V}:=\int_{\Gamma_{1}} l g(z) v \mathrm{~d} \Gamma, \quad z, v \in V
$$

defines a map $B: V \rightarrow V^{\prime}$ (not linear in general). Indeed, using (16) and the trace and imbedding theorems

$$
\begin{equation*}
V \subset H^{1}(\Omega) \hookrightarrow H^{1 / 2}\left(\Gamma_{1}\right) \subset L^{(2 n-2) /(n-2)}\left(\Gamma_{1}\right) \tag{25}
\end{equation*}
$$

we have, using for brevity the notation $\|\cdot\|_{s}:=\|\cdot\|_{L^{s}\left(\Gamma_{1}\right)}$,

$$
\begin{gathered}
\left|\int_{\Gamma_{1}} l g(z) v \mathrm{~d} \Gamma\right| \leq c\|g(z)\|_{(2 n-2) / n}\|v\|_{(2 n-2) /(n-2)} \\
\leq c\left(1+\|z\|_{(2 n-2) /(n-2)}^{n /(n-2)}\right)\|v\|_{(2 n-2) /(n-2)} \leq c\left(1+\|z\|_{V}^{n /(n-2)}\right)\|v\|_{V}<+\infty
\end{gathered}
$$

Hence $B z \in V^{\prime}$ for every $z \in V$.
In order to find a reasonable definition of the (weak) solution of the problem (8)-(11) we multiply the equation (8) by $v \in V$ and we integrate by parts. Using the boundary conditions (9), (10) we obtain that

$$
\begin{gathered}
0=\int_{\Omega}\left(u^{\prime \prime}-\Delta u+q u\right) v \mathrm{dx} \\
=\int_{\Omega} u^{\prime \prime} v+\nabla u \cdot \nabla v+q u v \mathrm{dx}-\int_{\Gamma}\left(\partial_{\nu} u\right) v \mathrm{~d} \Gamma \\
=\int_{\Omega} u^{\prime \prime} v+\nabla u \cdot \nabla v+q u v \mathrm{dx}+\int_{\Gamma_{1}} a u v+l g\left(u^{\prime}\right) v \mathrm{~d} \Gamma \\
=\left\langle u^{\prime \prime}+A u+B u^{\prime}, v\right\rangle_{V^{\prime}, V},
\end{gathered}
$$

whence

$$
\begin{equation*}
u^{\prime \prime}+A u+B u^{\prime}=0 \quad \text { on } \quad \mathbb{R}_{+} \tag{26}
\end{equation*}
$$

Putting

$$
U=\left(U_{1}, U_{2}\right):=\left(u, u^{\prime}\right) \quad \text { and } \quad \mathcal{A} U:=\left(-U_{2}, A U_{1}+B U_{2}\right)
$$

we may write (26), (11) in the form

$$
\begin{equation*}
U^{\prime}+\mathcal{A} U=0 \quad \text { in } \quad \mathbb{R}_{+}, \quad U(0)=\left(u^{0}, u^{1}\right) \tag{27}
\end{equation*}
$$

If $u$ is a (sufficiently smooth) solution of (8)-(11), then

$$
U(t), \quad \mathcal{A} U(t) \in V \times H, \quad \forall t \in \mathbb{R}_{+} .
$$

This leads us to set

$$
\begin{equation*}
\mathcal{H}:=V \times H, \quad D(\mathcal{A}):=\left\{U \in V \times V: A U_{1}+B U_{2} \in H\right\} \tag{28}
\end{equation*}
$$

and to define the solution of (8)-(11) as that of (27). The definition is justified by

Proposition 7.6. - $\mathcal{A}$ is a maximal monotone operator in $\mathcal{H}$.
Proof. - The monotonicity of $\mathcal{A}$ follows from the nonnegativity of $l$ and from the non-decreasingness of $g$. Indeed, given $U, V \in D(\mathcal{A})$ arbitrarily, we have

$$
\begin{gathered}
(\mathcal{A} U-\mathcal{A} V, U-V)_{\mathcal{H}} \\
=\left(V_{2}-U_{2}, U_{1}-V_{1}\right)_{V}+\left(A U_{1}-A V_{1}+B U_{2}-B V_{2}, U_{2}-V_{2}\right)_{H} \\
=\left(V_{2}-U_{2}, U_{1}-V_{1}\right)_{V}+\left\langle A U_{1}-A V_{1}+B U_{2}-B V_{2}, U_{2}-V_{2}\right\rangle_{V^{\prime}, V} \\
=\left\langle B U_{2}-B V_{2}, U_{2}-V_{2}\right\rangle_{V^{\prime}, V}=\int_{\Gamma_{1}} l\left(g\left(U_{2}\right)-g\left(V_{2}\right)\right)\left(U_{2}-V_{2}\right) \mathrm{d} \Gamma \geq 0 .
\end{gathered}
$$

It remains to show that for $W=\left(W_{1}, W_{2}\right) \in \mathcal{H}=V \times H$ given arbitrarily, there exists $U=\left(U_{1}, U_{2}\right) \in D(\mathcal{A})$ such that $(I+\mathcal{A}) U=W$. It is sufficient to show that the map $I+A+B: V \rightarrow V^{\prime}$ is onto. Indeed, then there exists $U_{2} \in V$ satisfying

$$
(I+A+B) U_{2}=W_{2}-A W_{1}
$$

Choosing $U_{2}$ in this way and setting $U_{1}=U_{2}+W_{1}$ we clearly have $U \in V \times V$, $A U_{1}+B U_{2}=W_{2}-U_{2} \in H$ (hence $\left.U \in D(\mathcal{A})\right)$ and $(I+\mathcal{A}) U=W$.

To prove the surjectivity of the map $I+A+B: V \rightarrow V^{\prime}$ fix $f \in V^{\prime}$ arbitrarily, set

$$
G(t)=\int_{0}^{t} g(s) \mathrm{ds}, \quad t \in \mathbb{R}
$$

and consider the map $F: V \rightarrow \mathbb{R}$ defined by

$$
F(u)=\frac{1}{2}\|u\|_{H}^{2}+\frac{1}{2}\|u\|_{V}^{2}+\int_{\Gamma} G(u) \mathrm{d} \Gamma-\langle f, u\rangle_{V^{\prime}, V}
$$

Using the growth assumption (16) one may readily verify that the map $F$ is well-defined, continuously differentiable and that

$$
F^{\prime}(u) v=\langle(I+A+B) u-f, v\rangle_{V^{\prime}, V}, \quad \forall u, v \in V
$$

Furthermore, the monotonicity of $g$ implies the convexity of $F$. Finally, $F$ is coercive : $F(v) \rightarrow+\infty$ if $\|v\|_{V} \rightarrow+\infty$. This follows at once from the obvious inequality

$$
F(v) \geq\left(\frac{1}{2}\|v\|_{V}-\|f\|_{V^{\prime}}\right)\|v\|_{V}
$$

It follows that the infimum of $F$ is attained at some point $u \in V$. Then $F^{\prime}(u)=0$ i.e. $(I+A+B) u=f$.

Now set

$$
\begin{equation*}
D_{0}:=\left\{U \in V \times V: U_{1} \in H^{2}(\Omega), \quad \partial_{\nu} U_{1}+a U_{1}+l g\left(U_{2}\right)=0 \quad \text { on } \quad \Gamma_{1}\right\} \tag{29}
\end{equation*}
$$

and

$$
\begin{equation*}
D:=\left\{U \in V \times H_{0}^{1}(\Omega): U_{1} \in H^{2}(\Omega) \text { and } \partial_{\nu} U_{1}+a U_{1}=0 \quad \text { on } \quad \Gamma_{1}\right\} . \tag{30}
\end{equation*}
$$

Lemma 7.7. - We have $D \subset D_{0} \subset D(\mathcal{A})$, and $D$ is dense in $\mathcal{H}$. Consequently, $D(\mathcal{A})$ is dense in $\mathcal{H}$.

Proof. - It is clear that $D \subset D_{0}$ and that $D$ is dense in $\mathcal{H}=H_{\Gamma_{0}}^{1}(\Omega) \times L^{2}(\Omega)$. It remains to show that $D_{0} \subset D(\mathcal{A})$.

Fix $\left(U_{1}, U_{2}\right) \in D_{0}$ arbitrarily ; it suffices to prove the estimate

$$
\begin{equation*}
\left|\left\langle A U_{1}+B U_{2}, v\right\rangle_{V^{\prime}, V}\right| \leq c\|v\|_{H}, \quad \forall v \in V \tag{31}
\end{equation*}
$$

with a suitable constant $c$. Using the definition of $A$ and $B$ we have

$$
\begin{equation*}
\left\langle A U_{1}+B U_{2}, v\right\rangle_{V^{\prime}, V}=\int_{\Omega} \nabla U_{1} \cdot \nabla v+q U_{1} v \mathrm{dx}+\int_{\Gamma_{1}} a U_{1} v+l g\left(U_{2}\right) v \mathrm{~d} \Gamma \tag{32}
\end{equation*}
$$

Since $\left(U_{1}, U_{2}\right) \in D_{0}$ implies $U_{1} \in H^{2}(\Omega)$, we may apply Green's formula to the right-hand side. We obtain that

$$
\begin{gathered}
\left\langle A U_{1}+B U_{2}, v\right\rangle_{V^{\prime}, V} \\
=\int_{\Omega}\left(-\Delta U_{1}+q U_{1}\right) v \mathrm{dx}+\int_{\Gamma}\left(\partial_{\nu} U_{1}\right) v \mathrm{~d} \Gamma+\int_{\Gamma_{1}} a U_{1} v+\lg \left(U_{2}\right) v \mathrm{~d} \Gamma
\end{gathered}
$$

It follows from the definition of $D_{0}$ that the boundary integrals vanish. Hence

$$
\left\langle A U_{1}+B U_{2}, v\right\rangle_{V^{\prime}, V}=\int_{\Omega}\left(-\Delta U_{1}+q U_{1}\right) v \mathrm{dx}
$$

Since $-\Delta U_{1}+q U_{1} \in L^{2}(\Omega)=H$, hence (31) follows.
Proof of theorem 7.4. - Theorem 7.4 is an immediate consequence of theorem 7.1, proposition 7.6 and lemma 7.7 : observe that

$$
E(u)=\frac{1}{2}\|U(t)\|_{\mathcal{H}}^{2} \text { and } E(u-v)=\frac{1}{2}\|U(t)-V(t)\|_{\mathcal{H}}^{2} .
$$

For the proof of theorem 7.5 we need another lemma.
Lemma 7.8. - Let g be globally Lipschitz continuous. Then

$$
\begin{equation*}
D(\mathcal{A}) \subset D_{0} . \tag{33}
\end{equation*}
$$

Proof. - Fix $\left(U_{1}, U_{2}\right) \in D(\mathcal{A})$ arbitrarily and set $f:=A U_{1}+B U_{2}$, $h:=-a U_{1}-l g\left(U_{2}\right)$. We will show that $U_{1}$ is the weak solution of the problem

$$
\begin{gather*}
-\Delta U_{1}+q U_{1}=f \quad \text { in } \quad \Omega,  \tag{34}\\
U_{1}=0 \quad \text { on } \quad \Gamma_{0}  \tag{35}\\
\partial_{\nu} U_{1}=h \quad \text { on } \quad \Gamma_{1} . \tag{36}
\end{gather*}
$$

This will imply at once the boundary condition in the definition of $D_{0}$.
Let us recall that by definition $U_{1} \in V$ is the weak solution of $(34)-(36)$ if

$$
\int_{\Omega} \nabla U_{1} \cdot \nabla v+q U_{1} v \mathrm{dx}=\langle f, v\rangle_{V^{\prime}, V}+\int_{\Gamma_{1}} h v \mathrm{~d} \Gamma
$$

for all $v \in V$. We may rewrite this equation in the form

$$
\left\langle A U_{1}+B U_{2}, v\right\rangle_{V^{\prime}, V}=\int_{\Omega} \nabla U_{1} \cdot \nabla v+q U_{1} v \mathrm{dx}+\int_{\Gamma_{1}} a U_{1} v+l g\left(U_{2}\right) v \mathrm{~d} \Gamma
$$

and (as we have seen in the proof of the preceding lemma, cf. (32)) this is a direct consequence of the definitions of $A$ and $B$.

It remains to prove that $U_{1} \in H^{2}(\Omega)$; applying the elliptic regularity theory to the problem (34)-(36), it is sufficient to establish that

$$
\begin{equation*}
f \in L^{2}(\Omega) \quad \text { and } \quad h \in H^{1}(\Omega) \tag{37}
\end{equation*}
$$

The first relation in (37) follows from the definition of $D(\mathcal{A})$.
For the second relation it suffices to show that

$$
\begin{equation*}
g\left(U_{2}\right) \in L^{2}(\Omega) \quad \text { and } \quad \nabla g\left(U_{2}\right) \in L^{2}(\Omega) \tag{38}
\end{equation*}
$$

Both properties are obvious because the global Lipschitz continuity implies the inequality

$$
\begin{equation*}
|g(x)| \leq 1+c|x|, \quad \forall x \in \mathbb{R}: \tag{39}
\end{equation*}
$$

we have

$$
\int_{\Omega}\left|g\left(U_{2}\right)\right|^{2} \mathrm{dx} \leq c \int_{\Omega} 1+\left|U_{2}\right|^{2} \mathrm{dx} \leq c\left(1+\left\|U_{2}\right\|_{V}^{2}\right)<+\infty
$$

and

$$
\begin{aligned}
& \int_{\Omega}\left|\nabla g\left(U_{2}\right)\right|^{2} \mathrm{dx}=\int_{\Omega}\left|g^{\prime}\left(U_{2}\right) \nabla U_{2}\right|^{2} \mathrm{dx} \\
& \leq c \int_{\Omega}\left|\nabla U_{2}\right|^{2} \mathrm{dx} \leq\left\|U_{2}\right\|_{V}^{2}<+\infty
\end{aligned}
$$

Proof of theorem 7.5. - We apply the second half of theorem 7.1 and we use lemmas 7.7 and 7.8.

### 7.3. Kirchhoff plates

Let $\Omega$ be a bounded domain of class $C^{4}$ in $\mathbb{R}^{2}$ and let $\left\{\Gamma_{0}, \Gamma_{1}\right\}$ be a partition of its boundary $\Gamma$ such that

$$
\begin{equation*}
\Gamma_{0} \neq \emptyset \quad \text { and } \quad \overline{\Gamma_{0}} \cap \overline{\Gamma_{1}}=\emptyset \tag{40}
\end{equation*}
$$

Fix a number $\mu \in(0,1)$ and a nonnegative function $l \in C^{1}\left(\Gamma_{1}\right)$. Let $g: \mathbb{R} \rightarrow \mathbb{R}$ be a non-decreasing, continuous function such that $g(0)=0$ and consider the problem

$$
\begin{gather*}
u^{\prime \prime}+\Delta^{2} u=0 \quad \text { in } \quad \Omega \times \mathbb{R}_{+},  \tag{41}\\
u=u_{\nu}=0 \quad \text { on } \Gamma_{0} \times \mathbb{R}_{+},  \tag{42}\\
u_{\nu \nu}+\mu u_{\tau \tau}=0 \quad \text { on } \quad \Gamma_{1} \times \mathbb{R}_{+},  \tag{43}\\
u_{\nu \nu \nu}+(2-\mu) u_{\tau \tau \nu}=\lg \left(u^{\prime}\right) \quad \text { on } \quad \Gamma_{1} \times \mathbb{R}_{+},  \tag{44}\\
u(0)=u^{0} \quad \text { and } \quad u^{\prime}(0)=u^{1} \quad \text { on } \quad \Omega . \tag{45}
\end{gather*}
$$

This models the small transversal vibrations of a thin plate whose Poisson coefficient is equal to $\mu$; see e.g. Lagnese and Lions [1] or Lagnese [2]. Here and in the sequel the subscripts $\nu$ and $\tau$ stand for the normal and tangential derivatives; the unit normal and tangential vectors are given by $\nu=\left(\nu_{1}, \nu_{2}\right)$ and $\tau:=\left(-\nu_{2}, \nu_{1}\right)$, respectively. We shall use the notation $\nabla u=\left(u_{x}, u_{y}\right)$ for the gradient of $u$, and we introduce the quadratic form

$$
\begin{equation*}
Q(u)=u_{x x}^{2}+u_{y y}^{2}+2 \mu u_{x x} u_{y y}+2(1-\mu) u_{x y}^{2} . \tag{46}
\end{equation*}
$$

We shall write $d X:=d x d y$.
Setting

$$
\begin{equation*}
H_{\Gamma_{0}}^{2}(\Omega):=\left\{v \in H^{2}(\Omega): v=v_{\nu}=0 \quad \text { on } \quad \Gamma_{0}\right\} \tag{47}
\end{equation*}
$$

we have the following result :
Theorem 7.9. - Given $\left(u^{0}, u^{1}\right) \in H_{\Gamma_{0}}^{2}(\Omega) \times L^{2}(\Omega)$ arbitrarily, the problem (41)-(45) has a unique solution satisfying

$$
\begin{equation*}
u \in C\left(\mathbb{R}_{+} ; H_{\Gamma_{0}}^{2}(\Omega)\right) \cap C^{1}\left(\mathbb{R}_{+} ; L^{2}(\Omega)\right) \tag{48}
\end{equation*}
$$

The energy $E: \mathbb{R}_{+} \rightarrow \mathbb{R}_{+}$of the solution, defined by

$$
\begin{equation*}
E=E(u):=\frac{1}{2} \int_{\Omega}\left(u^{\prime}\right)^{2}+Q(u) \mathrm{dX} \tag{49}
\end{equation*}
$$

is non-increasing. Moreover, if $u$ and $v$ are two solutions (corresponding to different initial data), then the function $E(u-v)$ is non-increasing; in particular,

$$
\begin{equation*}
\|E(u-v)\|_{L^{\infty}\left(\mathbb{R}_{+}\right)} \leq E(u-v)(0) . \tag{50}
\end{equation*}
$$

If $g$ is globally Lipschitz and if

$$
\begin{equation*}
\left(u^{0}, u^{1}\right) \in\left(H^{4}(\Omega) \cap H_{\Gamma_{0}}^{2}(\Omega)\right) \times H_{\Gamma_{0}}^{2}(\Omega) \tag{51}
\end{equation*}
$$

and

$$
\begin{equation*}
u_{\nu \nu}^{0}+\mu u_{\tau \tau}^{0}=0 \text { and } u_{\nu \nu \nu}^{0}+(2-\mu) u_{\tau \tau \nu}^{0}=\lg \left(u^{1}\right) \operatorname{sur} \Gamma_{1} \times \mathbb{R}_{+}, \tag{52}
\end{equation*}
$$

then the solution is more regular :

$$
\begin{equation*}
u \in L^{\infty}\left(\mathbb{R}_{+} ; H^{4}(\Omega)\right), \quad u^{\prime} \in L^{\infty}\left(\mathbb{R}_{+} ; H_{\Gamma_{0}}^{2}(\Omega)\right) \tag{53}
\end{equation*}
$$

and

$$
\begin{equation*}
u^{\prime \prime} \in L^{\infty}\left(\mathbb{R}_{+} ; L^{2}(\Omega)\right) \tag{54}
\end{equation*}
$$

We begin by defining the solution of (41)-(45). Set

$$
V:=H_{\Gamma_{0}}^{2}(\Omega) \quad \text { and } \quad H:=L^{2}(\Omega)
$$

with

$$
\|v\|_{V}:=\left(\int_{\Omega} Q(v) \mathrm{dX}\right)^{1 / 2} \text { and }\|v\|_{H}:=\left(\int_{\Omega} v^{2} \mathrm{dX}\right)^{1 / 2}
$$

by hypothesis $\Gamma_{0} \neq \emptyset$ they are separable Hilbert spaces with a dense and compact imbedding $V \subset H$. We introduce the duality mapping $A: V \rightarrow V^{\prime}$ and we identify $H$ with $H^{\prime}$.

By the continuity of $g$ and to the Sobolev imbedding $H^{2}(\Omega) \subset C(\bar{\Omega})$ the formula

$$
\langle B z, v\rangle_{V^{\prime}, V}:=\int_{\Gamma_{1}} l g(z) v \mathrm{~d} \Gamma, \quad z, v \in V
$$

defines a (non-linear) map $B: V \rightarrow V^{\prime}$. As in the preceding section for the wave equation, one can readily verify that $B$ is monotone and hemicontinuous.

Next we need two lemmas.
Lemma 7.10. - Given $v \in H^{3}(\Omega)$ and $w \in H^{2}(\Omega)$ arbitrarily, the following identity holds true :

$$
\begin{equation*}
\int_{\Omega} v_{x x} w_{y y}+v_{y y} w_{x x}-2 v_{x y} w_{x y} \mathrm{dX}=\int_{\Gamma} v_{\tau \tau} w_{\nu}-v_{\nu \tau} w_{\tau} \mathrm{d} \Gamma . \tag{55}
\end{equation*}
$$

If $v \in H^{4}(\Omega)$, then we also have

$$
\begin{equation*}
\int_{\Omega} v_{x x} w_{y y}+v_{y y} w_{x x}-2 v_{x y} w_{x y} \mathrm{dX}=\int_{\Gamma} v_{\tau \tau} w_{\nu}+v_{\nu \tau \tau} w \mathrm{~d} \Gamma . \tag{56}
\end{equation*}
$$

Proof. - We apply Green's formula as follows :

$$
\begin{gathered}
\int_{\Omega} v_{x x} w_{y y} \mathrm{dX}=-\int_{\Omega} v_{x x y} w_{y} \mathrm{dX}+\int_{\Gamma} v_{x x} \nu_{2} w_{y} \mathrm{~d} \Gamma \\
=\int_{\Omega} v_{x y} w_{x y} \mathrm{dX}+\int_{\Gamma} w_{y}\left(\nu_{2} v_{x x}-\nu_{1} v_{x y}\right) \mathrm{d} \Gamma \\
=\int_{\Omega} v_{x y} w_{x y} \mathrm{~d} \mathrm{X}-\int_{\Gamma} v_{x \tau} w_{y} \mathrm{~d} \Gamma
\end{gathered}
$$

Using the obvious differential relations

$$
\begin{equation*}
u_{x}=\nu_{1} u_{\nu}-\nu_{2} u_{\tau}, \quad u_{y}=\nu_{2} u_{\nu}+\nu_{1} u_{\tau} \tag{57}
\end{equation*}
$$

hence we deduce that

$$
\begin{equation*}
\int_{\Omega} v_{x x} w_{y y}-v_{x y} w_{x y} \mathrm{dX}=-\int_{\Gamma}\left(\nu_{1} v_{\nu \tau}-\nu_{2} v_{\tau \tau}\right)\left(\nu_{2} w_{\nu}+\nu_{1} w_{\tau}\right) \mathrm{d} \Gamma . \tag{58}
\end{equation*}
$$

Analogously, we have

$$
\begin{equation*}
\int_{\Omega} v_{y y} w_{x x}-v_{x y} w_{x y} \mathrm{dX}=\int_{\Gamma}\left(\nu_{2} v_{\nu \tau}+\nu_{1} v_{\tau \tau}\right)\left(\nu_{1} w_{\nu}-\nu_{2} w_{\tau}\right) \mathrm{d} \Gamma \tag{59}
\end{equation*}
$$

Adding (58) to (59) we find (55).
If $v \in H^{4}(\Omega)$, then (56) follows from (55) by integration by parts.
Let us introduce the notation (compare with (46))

$$
\begin{equation*}
Q(u, v)=u_{x x} v_{x x}+u_{y y} v_{y y}+\mu\left(u_{x x} v_{y y}+u_{y y} v_{x x}\right)+2(1-\mu) u_{x y} v_{x y} . \tag{60}
\end{equation*}
$$

Lemma 7.11. - Given $v \in H^{4}(\Omega)$ and $w \in H^{2}(\Omega)$ arbitrarily, the following identity holds true :

$$
\begin{equation*}
\int_{\Omega}\left(\Delta^{2} v\right) w-Q(v, w) \mathrm{dX}=\int_{\Gamma}\left(v_{\nu \nu}+(2-\mu) v_{\tau \tau}\right)_{\nu} w-\left(v_{\nu \nu}+\mu v_{\tau \tau}\right) w_{\nu} \mathrm{d} \Gamma \tag{61}
\end{equation*}
$$

Proof. - We use (56) and (60) in the following way :

$$
\begin{gathered}
\int_{\Omega}\left(\Delta^{2} v\right) w-Q(v, w) \mathrm{dX}=\int_{\Omega} \Delta v \Delta w-Q(v, w) \mathrm{dX}+\int_{\Gamma}(\Delta v)_{\nu} w-(\Delta v) w_{\nu} \mathrm{d} \Gamma \\
=(1-\mu) \int_{\Omega} v_{x x} w_{y y}+v_{y y} w_{x x}-2 v_{x y} w_{x y} \mathrm{dX} \\
\quad+\int_{\Gamma}\left(v_{\nu \nu \nu}+v_{\tau \tau \nu}\right) w-\left(v_{\nu \nu}+v_{\tau \tau}\right) w_{\nu} \mathrm{d} \Gamma \\
=(1-\mu) \int_{\Gamma} v_{\tau \tau} w_{\nu}+v_{\nu \tau \tau} w \mathrm{~d} \Gamma+\int_{\Gamma}\left(v_{\nu \nu \nu}+v_{\tau \tau \nu}\right) w-\left(v_{\nu \nu}+v_{\tau \tau}\right) w_{\nu} \mathrm{d} \Gamma \\
=\int_{\Gamma}\left(v_{\nu \nu \nu}+(2-\mu) v_{\tau \tau \nu}\right) w-\left(v_{\nu \nu}+\mu v_{\tau \tau}\right) w_{\nu} \mathrm{d} \Gamma .
\end{gathered}
$$

Let $v \in V$ and multiply the equation (41) by $v$. Integrating on $\Omega$, applying the preceding lemma and using the boundary conditions (42)-(44) we obtain by a formal computation that

$$
\begin{gathered}
0=\int_{\Omega}\left(u^{\prime \prime}+\Delta^{2} u\right) v \mathrm{dX}=\int_{\Omega} u^{\prime \prime} v+Q(u, v) \mathrm{dX} \\
+\int_{\Gamma}\left(u_{\nu \nu}+(2-\mu) u_{\tau \tau}\right)_{\nu} v-\left(u_{\nu \nu}+\mu u_{\tau \tau}\right) v_{\nu} \mathrm{d} \Gamma \\
=\left\langle u^{\prime \prime}+A u, v\right\rangle_{V^{\prime}, V}+\int_{\Gamma} l g\left(u^{\prime}\right) v \mathrm{~d} \Gamma \\
=\left\langle u^{\prime \prime}+A u+B u^{\prime}, v\right\rangle_{V^{\prime}, V}
\end{gathered}
$$

whence

$$
u^{\prime \prime}+A u+B u^{\prime}=0 \quad \text { on } \quad \mathbb{R}_{+}
$$

Putting

$$
\begin{equation*}
U=\left(U_{1}, U_{2}\right):=\left(u, u^{\prime}\right) \quad \text { and } \quad \mathcal{A} U:=\left(-U_{2}, A U_{1}+B U_{2}\right) \tag{62}
\end{equation*}
$$

we may rewrite this equation as

$$
U^{\prime}+\mathcal{A} U=0 \quad \text { on } \quad \mathbb{R}_{+} .
$$

Taking into account the boundary and initial conditions it is natural to define $\mathcal{H}$ and $D(\mathcal{A})$ by

$$
\begin{align*}
\mathcal{H} & :=V \times H\left(=H_{\Gamma_{0}}^{2}(\Omega) \times L^{2}(\Omega)\right),  \tag{63}\\
D(\mathcal{A}) & :=\left\{U \in V \times V: A U_{1}+B U_{2} \in H\right\}, \tag{64}
\end{align*}
$$

and to define the solution of (41)-(45) as that of

$$
\begin{equation*}
U^{\prime}+\mathcal{A} U=0 \quad \text { on } \quad \mathbb{R}_{+}, \quad U(0)=\left(u^{0}, u^{1}\right) \tag{65}
\end{equation*}
$$

Proposition 7.12. - $\mathcal{A}$ is a maximal monotone operator in $\mathcal{H}$.
Proof. - We may repeat word by word the proof of proposition 7.6 in the preceding section.

Set

$$
\begin{align*}
& D_{0}:=\left\{U \in V \times V: U_{1} \in H^{4}(\Omega), \quad U_{1, \nu \nu}+\mu U_{1, \tau \tau}=0\right. \\
&\text { and } \left.\quad U_{1, \nu \nu \nu}+(2-\mu) U_{1, \tau \tau \nu}=\lg \left(U_{2}\right) \text { on } \Gamma_{1}\right\} \tag{66}
\end{align*}
$$

and

$$
\begin{gather*}
D:=\left\{U \in V \times H_{0}^{2}(\Omega): U_{1} \in H^{4}(\Omega)\right. \text { and } \\
\left.U_{1, \nu \nu}+\mu U_{1, \tau \tau}=U_{1, \nu \nu \nu}+(2-\mu) U_{1, \tau \tau \nu}=0 \text { on } \Gamma_{1}\right\} . \tag{67}
\end{gather*}
$$

Lemma 7.13. - We have $D \subset D_{0}=D(\mathcal{A})$ and $D$ is dense in $\mathcal{H}$. Consequently, $D(\mathcal{A})$ is dense in $\mathcal{H}$.

Proof. - Observe that for $U \in V \times V$ the equality $A U_{1}+B U_{2}=f$ means that $U_{1}$ is the weak solution of the problem

$$
\begin{gather*}
\Delta^{2} U_{1}=f \quad \text { in } \Omega,  \tag{68}\\
U_{1}=U_{1, \nu}=0 \quad \text { on } \Gamma_{0},  \tag{69}\\
U_{1, \nu \nu}+\mu U_{1, \tau \tau}=0 \quad \text { on } \quad \Gamma_{1},  \tag{70}\\
U_{1, \nu \nu \nu}+(2-\mu) U_{1, \tau \tau \nu}=\lg \left(U_{2}\right) \quad \text { on } \quad \Gamma_{1} . \tag{71}
\end{gather*}
$$

Indeed, if $A U_{1}+B U_{2}=f$, then for every fixed $v \in V$ we have, using lemma 7.11,

$$
\langle f, v\rangle_{V^{\prime}, V}=\int_{\Omega} Q\left(U_{1}, v\right) \mathrm{dX}+\int_{\Gamma_{1}} l g\left(U_{2}\right) v \mathrm{~d} \Gamma, \quad \forall v \in V
$$

and this is the usual definition of the weak solution $U_{1} \in V$ of (68)-(71). (It is easy to show, using lemma 7.11, that every regular solution is also a weak solution : adapt the formal computation leading to (65).)

Using this observation one can readily verify that $D \subset D_{0} \subset D(\mathcal{A})$ and that $D$ is dense in $\mathcal{H}$. To show the inclusion $D(\mathcal{A}) \subset D_{0}$ we verify that $U \in D(\mathcal{A})$ implies $g\left(U_{2}\right) \in H^{1}(\Omega)$; indeed, then $U_{1} \in H^{4}(\Omega)$ by the elliptic regularity theory applied to the problem (68)-(71).

Since $U_{2} \in V \subset H^{2}(\Omega) \subset L^{\infty}(\Omega)$ and $g, g^{\prime}$ are locally bounded, we have $g\left(U_{2}\right) \in L^{\infty}(\Omega)$ and $g^{\prime}\left(U_{2}\right) \in L^{\infty}(\Omega)$; in particular, $g\left(U_{2}\right) \in L^{2}(\Omega)$. It remains to show that $\nabla g\left(U_{2}\right) \in L^{2}(\Omega)$. This follows from the equality $\nabla g\left(U_{2}\right)=g^{\prime}\left(U_{2}\right) \nabla U_{2}$ because $g^{\prime}\left(U_{2}\right) \in L^{\infty}(\Omega)$ and $\nabla U_{2} \in H^{1}(\Omega) \subset L^{2}(\Omega)$. $\square$

Proof of theorem 7.9. - We apply theorem 7.1 and we use proposition 7.12 and lemma 7.13.

## 8. Linear stabilization

The aim of this chapter is to prove the exponential energy decay of the solutions of the wave equation under suitable linear boundary feedbacks (i.e. $g(x) \equiv x)$. We shall also apply a principle of RusSell [2] relating the problem of stabilization to that of exact controllability. In the last section we shall study Maxwell's equations, which are closely related to the wave equation.

As usual, for a given point $x^{0} \in \mathbb{R}^{n}$ we shall use the notation

$$
\begin{gathered}
m(x)=x-x^{0}, \quad x \in \mathbb{R}^{n}, \\
R=R\left(x^{0}\right)=\sup \left\{\left|x-x^{0}\right|: x \in \Omega\right\}, \\
\Gamma_{+}=\{x \in \Gamma: m(x) \cdot \nu(x)>0\}, \\
\mathrm{d} \Gamma_{m}=(m \cdot \nu) \mathrm{d} \Gamma .
\end{gathered}
$$

As in $\S 7.2$, we restrict ourselves for brevity to the case of dimension $n \geq 3$. We shall consider the real case only; the generalization to the complex case is obvious.

### 8.1. An integral inequality

In this section we recall the following simple result (used e.g. in Haraux [2] and Lagnese [2]) ; it will play an important role in this chapter.

Theorem 8.1. - Let $E: \mathbb{R}_{+} \rightarrow \mathbb{R}_{+}\left(\mathbb{R}_{+}:=[0,+\infty)\right)$ be a non-increasing function and assume that there exists a constant $T>0$ such that

$$
\begin{equation*}
\int_{t}^{\infty} E(s) \mathrm{ds} \leq T E(t), \quad \forall t \in \mathbb{R}_{+} . \tag{1}
\end{equation*}
$$

Then

$$
\begin{equation*}
E(t) \leq E(0) e^{1-t / T}, \quad \forall t \geq T \tag{2}
\end{equation*}
$$

Observe that the inequality (2) is also satisfied for $0 \leq t<T$ : indeed, then it is weaker than the trivial inequality $E(t) \leq E(0)$.

Proof of the theorem. - Define

$$
f(x):=e^{x / T} \int_{x}^{\infty} E(s) \mathrm{ds}, \quad x \in \mathbb{R}_{+}
$$

then $f$ is locally absolutely continuous and it is also non-increasing by (1) :

$$
f^{\prime}(x)=T^{-1} e^{x / T}\left(\int_{x}^{\infty} E(s) \mathrm{ds}-T E(x)\right) \leq 0
$$

almost everywhere in $\mathbb{R}_{+}$. Hence, using (1) again,

$$
f(x) \leq f(0)=\int_{0}^{\infty} E(s) \mathrm{ds} \leq T E(0), \quad \forall x \in \mathbb{R}_{+}
$$

i.e.

$$
\begin{equation*}
\int_{x}^{\infty} E(s) \mathrm{ds} \leq T E(0) e^{-x / T}, \quad \forall x \in \mathbb{R}_{+} \tag{3}
\end{equation*}
$$

Since $E$ is nonnegative and non-increasing, we have

$$
\int_{x}^{\infty} E(s) \mathrm{ds} \geq \int_{x}^{x+T} E(s) \mathrm{ds} \geq T E(x+T)
$$

Substituting into (3) we obtain that

$$
E(x+T) \leq E(0) e^{-x / T}, \quad \forall x \in \mathbb{R}_{+} ;
$$

setting $t:=x+T$ hence we conclude (2). $\quad \square$
Remark 8.2. - The theorem is optimal in the following sense : given $T>0$ and $t^{\prime} \geq T$ arbitrarily, there exists a non-increasing function $E: \mathbb{R}_{+} \rightarrow \mathbb{R}_{+}$, non identically zero, satisfying (1) and such that

$$
E\left(t^{\prime}\right)=E(0) e^{1-t^{\prime} / T}
$$

see the limit $\alpha \rightarrow 0$ of the example given by formula (9.5) in the next chapter.
Remark 8.3. - If the function $E$ is also continuous, then the inequalities (2) are strict; in particular, $E(T)<E(0)$. This result is also optimal, see Komornik [12].

### 8.2. Uniform stabilization of the wave equation I

Consider the linear case of the problem introduced in section 7.2 :

$$
\begin{gather*}
u^{\prime \prime}-\Delta u+q u=0 \quad \text { in } \quad \Omega \times \mathbb{R}_{+},  \tag{4}\\
u=0 \quad \text { on } \quad \Gamma_{0} \times \mathbb{R}_{+},  \tag{5}\\
\partial_{\nu} u+a u+l u^{\prime}=0 \quad \text { on } \quad \Gamma_{1} \times \mathbb{R}_{+},  \tag{6}\\
u(0)=u^{0} \quad \text { and } \quad u^{\prime}(0)=u^{1} \quad \text { sur } \quad \Omega \tag{7}
\end{gather*}
$$

where we assume (7.12)-(7.15).
Recall the definition of the energy :

$$
\begin{equation*}
E(t):=\frac{1}{2} \int_{\Omega}\left(u^{\prime}(t)\right)^{2}+|\nabla u(t)|^{2}+q(u(t))^{2} \mathrm{dx}+\frac{1}{2} \int_{\Gamma_{1}} a(u(t))^{2} \mathrm{~d} \Gamma \tag{8}
\end{equation*}
$$

The following result shows in particular that the energy is non-increasing.
Lemma 8.4. - Given $\left(u^{0}, u^{1}\right) \in D(\mathcal{A})$ arbitrarily, the solution of (4)-(7) satisfies the energy equalities

$$
\begin{equation*}
E(S)-E(T)=\int_{S}^{T} \int_{\Gamma_{1}} l\left(u^{\prime}(t)\right)^{2} \mathrm{~d} \Gamma, \quad 0 \leq S<T<\infty . \tag{9}
\end{equation*}
$$

Indeed, by the nonnegativity of $l$ the right-hand side of (9) is nonnegative whence $E(S) \geq E(T)$.

Proof. - We multiply equation (4) by $u^{\prime}$ and we integrate by parts in $\Omega \times(S, T)$. Using (5) and (6) we obtain that

$$
\begin{gathered}
0=\int_{S}^{T} \int_{\Omega} u^{\prime}\left(u^{\prime \prime}-\Delta u+q u\right) \mathrm{dx} \mathrm{dt} \\
=\left[\frac{1}{2} \int_{\Omega}\left(u^{\prime}\right)^{2}+|\nabla u|^{2}+q u^{2} \mathrm{dx}\right]_{S}^{T}-\int_{S}^{T} \int_{\Gamma} u^{\prime} \partial_{\nu} u \mathrm{~d} \Gamma \mathrm{dt} \\
=\left[\frac{1}{2} \int_{\Omega}\left(u^{\prime}\right)^{2}+|\nabla u|^{2}+q u^{2} \mathrm{dx}+\frac{1}{2} \int_{\Gamma_{1}} a u^{2} \mathrm{~d} \Gamma \mathrm{dt}\right]_{S}^{T}+\int_{S}^{T} \int_{\Gamma_{1}} l\left(u^{\prime}\right)^{2} \mathrm{~d} \Gamma \mathrm{dt}
\end{gathered}
$$

and (9) follows from the definition of the energy.
Remark 8.5. - Lemma 8.4 permits us to define $\sqrt{l} u^{\prime}$ by density as an element of $L^{2}\left(\mathbb{R}_{+} ; L^{2}\left(\Gamma_{1}\right)\right)$, for every $\left(u^{0}, u^{1}\right) \in H_{\Gamma_{0}}^{1}(\Omega) \times L^{2}(\Omega)$. Then we have in particular

$$
\begin{equation*}
l u^{\prime} \in L^{2}\left(\mathbb{R}_{+} ; L^{2}\left(\Gamma_{1}\right)\right) \tag{10}
\end{equation*}
$$

The purpose of this section is to show that a particular choice of the feedback (i.e. of $l$ and $a$ ) leads to fast energy decay.

Assume that there is a point $x^{0} \in \mathbb{R}^{n}$ such that

$$
\begin{equation*}
m \cdot \nu \leq 0 \quad \text { on } \quad \Gamma_{0} \quad \text { and } \quad m \cdot \nu \geq 0 \quad \text { on } \quad \Gamma_{1} . \tag{11}
\end{equation*}
$$

We define $Q_{1} \geq 0$ as in section 3.4 :

$$
\begin{equation*}
Q_{1}:=2 R \lambda_{1}^{-1 / 2} \sup _{\Omega} q \tag{12}
\end{equation*}
$$

where $\lambda_{1}$ is the biggest constant such that

$$
\int_{\Omega}|\nabla v|^{2}+q v^{2} \mathrm{dx}+\int_{\Gamma_{1}} a v^{2} \mathrm{~d} \Gamma \geq \lambda_{1} \int_{\Omega} v^{2} \mathrm{dx}
$$

for every $v \in H_{\Gamma_{0}}^{1}(\Omega)$.

Theorem 8.6. - Assume (7.12)-(7.15), (11),

$$
\begin{equation*}
Q_{1}<1 \tag{13}
\end{equation*}
$$

and choose

$$
\begin{equation*}
l:=(m \cdot \nu) / R \quad \text { and } \quad a:=(n-1)(m \cdot \nu) /\left(2 R^{2}\right) . \tag{14}
\end{equation*}
$$

Then for every given $\left(u^{0}, u^{1}\right) \in H_{\Gamma_{0}}^{1}(\Omega) \times L^{2}(\Omega)$ the solution of (4)-(7) satisfies the following estimate:

$$
\begin{equation*}
E(t) \leq E(0) \exp \left(1-\left(1-Q_{1}\right) t / 2 R\right), \quad \forall t \in \mathbb{R}_{+} \tag{15}
\end{equation*}
$$

Theorem 8.6 generalizes a former result in Komornik [4] which improved earlier theorems of Slemrod [1], Quinn and Russell [1], Russell [1], Chen [1], [2], Lagnese [1], Triggiani [1], Komornik and Zuazua [1]. As for the case $n=2$ we refer to Tcheugoue [1].

Remark 8.7. - More general feedbacks will be considered later in section 8.4 (but the decay estimates will be weaker).

The proof of the theorem will be based on the following identity where we set

$$
k:=1 / R, \quad b:=(n-1) /\left(2 R^{2}\right) \quad \text { and } \quad M u:=2 m \cdot \nabla u+(n-1) u
$$

for brevity.
Lemma 8.8. - Given $\left(u^{0}, u^{1}\right) \in D(\mathcal{A})$ and $0 \leq S<T<\infty$ arbitrarily, the solution of (4)-(7) satisfies the following identity:

$$
\begin{gather*}
2 \int_{S}^{T} E d t=\left[\int_{\Omega} u^{\prime} M u \mathrm{dx}\right]_{T}^{S} \\
-\int_{S}^{T} \int_{\Omega}(n-2) q u^{2}+2 q u m \cdot \nabla u \mathrm{dx} \mathrm{dt}+\int_{S}^{T} \int_{\Gamma_{0}}\left|\partial_{\nu} u\right|^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt}  \tag{16}\\
+\int_{S}^{T} \int_{\Gamma_{1}}\left(u^{\prime}\right)^{2}-|\nabla u|^{2}+b u^{2}-\left(k u^{\prime}+b u\right) M u \mathrm{~d} \Gamma_{m} \mathrm{dt} .
\end{gather*}
$$

Proof. - Recall the identity (3.18) which was proved for every sufficiently smooth function $u$ satisfying (4) :

$$
\begin{gathered}
\int_{S}^{T} \int_{\Gamma}\left(\partial_{\nu} u\right) M u+(m \cdot \nu)\left(\left(u^{\prime}\right)^{2}-|\nabla u|^{2}\right) \mathrm{d} \Gamma \mathrm{dt} \\
=\left[\int_{\Omega} u^{\prime} M u \mathrm{dx}\right]_{S}^{T}+\int_{S}^{T} \int_{\Omega}\left(u^{\prime}\right)^{2}+|\nabla u|^{2}+(n-1) q u^{2}+2 q u m \cdot \nabla u \mathrm{dx} \mathrm{dt}
\end{gathered}
$$

Using (8) hence we deduce that

$$
\begin{gather*}
\int_{S}^{T} \int_{\Gamma}\left(\partial_{\nu} u\right) M u+(m \cdot \nu)\left(\left(u^{\prime}\right)^{2}-|\nabla u|^{2}\right) \mathrm{d} \Gamma \mathrm{dt} \\
=\left[\int_{\Omega} u^{\prime} M u \mathrm{dx}\right]_{S}^{T}+2 \int_{S}^{T} E \mathrm{dt}  \tag{17}\\
+\int_{S}^{T} \int_{\Omega}(n-2) q u^{2}+2 q u m \cdot \nabla u \mathrm{dx} \mathrm{dt}-\int_{S}^{T} \int_{\Gamma_{1}} b u^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt} .
\end{gather*}
$$

Using the boundary conditions (5) et (6) we may replace the left-hand side of (17) by

$$
\int_{S}^{T} \int_{\Gamma_{0}}\left|\partial_{\nu} u\right|^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt}+\int_{S}^{T} \int_{\Gamma_{1}}\left(u^{\prime}\right)^{2}-|\nabla u|^{2}-\left(k u^{\prime}+b u\right) M u \mathrm{~d} \Gamma_{m} \mathrm{dt}
$$

and (16) follows.
Proof of theorem 8.6. - Introduce $R_{1}$ as in section 3.4, then $R_{1}=R$ by the particular choice of $b$ and by remark 3.20.

Let us first fix $\left(u^{0}, u^{1}\right) \in D(\mathcal{A})$ arbitrarily. Using lemma 3.22 (it remains valid for every $t \in \mathbb{R}_{+}$) the first term on the right-hand side of (16) is majorized by $2 R E(S)+2 R E(T)$. Using the definition of $Q_{1}$ and $E$, the second term is majorized by $2 Q_{1} \int_{S}^{T} E$ dt. By (11) the third term is $\leq 0$. Thus we deduce from (16) the following inequality :

$$
\begin{gather*}
2\left(1-Q_{1}\right) \int_{S}^{T} E \mathrm{dt} \leq 2 R E(S)+2 R E(T) \\
+\int_{S}^{T} \int_{\Gamma_{1}}\left(u^{\prime}\right)^{2}-|\nabla u|^{2}+b u^{2}-\left(k u^{\prime}+b u\right) M u \mathrm{~d} \Gamma_{m} \mathrm{dt} \tag{18}
\end{gather*}
$$

In order to majorize the last term of this inequality we note that, by the particular choice (cf. (14)) of the coefficients $k$ and $b$ we have

$$
\begin{gathered}
\left(u^{\prime}\right)^{2}-|\nabla u|^{2}+b u^{2}-\left(k u^{\prime}+b u\right) M u \\
\leq\left(u^{\prime}\right)^{2}-|\nabla u|^{2}+b u^{2}-2\left(k u^{\prime}+b u\right) m \cdot \nabla u+(1-n) u\left(k u^{\prime}+b u\right) \\
\leq\left(u^{\prime}\right)^{2}+b u^{2}+R^{2}\left(k u^{\prime}+b u\right)^{2}+(1-n) u\left(k u^{\prime}+b u\right) \\
=\left(u^{\prime}\right)^{2}+b\left(2-n+R^{2} b\right) u^{2}+R^{2} k^{2}\left(u^{\prime}\right)^{2}+\left(1-n+2 R^{2} b\right) k u u^{\prime} \\
=2\left(u^{\prime}\right)^{2}+((3-n) / 2) b u^{2} .
\end{gathered}
$$

In view of condition (8) we conclude from (18) that

$$
\begin{gathered}
2\left(1-Q_{1}\right) \int_{S}^{T} E \mathrm{dt} \leq 2 R E(S)+2 R E(T) \\
+2 \int_{S}^{T} \int_{\Gamma_{1}}\left(u^{\prime}\right)^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt}+\frac{3-n}{2} \int_{S}^{T} \int_{\Gamma_{1}} b u^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt}
\end{gathered}
$$

Applying lemma 8.4 and using the definition of $k$ we conclude that

$$
\begin{equation*}
2\left(1-Q_{1}\right) \int_{S}^{T} E \mathrm{dt} \leq 4 R E(S)+\frac{3-n}{2} \int_{S}^{T} \int_{\Gamma_{1}} b u^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt} . \tag{19}
\end{equation*}
$$

Since $n \geq 3$, the last term is $\leq 0$. Therefore we deduce from (19) that

$$
2\left(1-Q_{1}\right) \int_{S}^{T} E \mathrm{dt} \leq 4 R E(S)
$$

Letting $T \rightarrow+\infty$ we obtain for every fixed $S \in \mathbb{R}_{+}$the estimate

$$
2\left(1-Q_{1}\right) \int_{S}^{\infty} E \mathrm{dt} \leq 4 R E(S)
$$

Using (13) and applying theorem 8.1 hence (15) follows for every $\left(u^{0}, u^{1}\right) \in$ $D(\mathcal{A})$.

Now fix $\left(u^{0}, u^{1}\right) \in \mathcal{H}=H_{\Gamma_{0}}^{1}(\Omega) \times L^{2}(\Omega)$ arbitrarily. Choose (using lemma 7.7) a sequence of initial data $\left(u_{j}^{0}, u_{j}^{1}\right) \in D(\mathcal{A})$ converging to $\left(u^{0}, u^{1}\right)$ in $\mathcal{H}$ and apply the estimate (15) for each $\left(u_{j}^{0}, u_{j}^{1}\right)$ :

$$
E_{j}(t) \leq E_{j}(0) \exp \left(1-\left(1-Q_{1}\right) t /(2 R)\right), \quad \forall t \in \mathbb{R}_{+}
$$

It follows from property (7.19) in theorem 7.4 we have $E_{j}(t) \rightarrow E(t)$ as $j \rightarrow \infty$, for each fixed $t \in \mathbb{R}_{+}$. Passing to the limit in the above estimates hence (15) follows.

Remark 8.9. - We observed in remarks 3.19 and 3.25 that hypotheses (7.13) and (11) together are very restrictive. It was shown in Komornik and Zuazua [1] that in dimension $n \leq 3$ hypothesis (8) is not necessary. For $n>3$ its necessity remains an interesting open problem.

### 8.3. Application to the exact controllability. Russell's principle

By a general principle of Russell [2] the stabilizability of a linear reversible system implies its exact controllability. We apply here the corresponding
construction to deduce from theorem 8.6 the following generalization (for $n \geq 3$ ) of theorem 4.8 on the exact controllability of the problem

$$
\begin{gather*}
y^{\prime \prime}-\Delta y+q y=0 \quad \text { in } \quad \Omega \times(0, T)  \tag{20}\\
y=0 \quad \text { on } \quad \Gamma_{0} \times(0, T)  \tag{21}\\
\partial_{\nu} y+a y=v \quad \text { on } \quad \Gamma_{1} \times(0, T),  \tag{22}\\
y(0)=y^{0} \quad \text { and } \quad y^{\prime}(0)=y^{1} \quad \text { on } \quad \Omega \tag{23}
\end{gather*}
$$

where we eliminate hypotheses (4.39) and (4.40) sur $a$.
Theorem 8.10. - Assume (11), (13) and let

$$
\begin{equation*}
T>2 R /\left(1-Q_{1}\right) \tag{24}
\end{equation*}
$$

Then for every given

$$
\left(y^{0}, y^{1}\right),\left(y_{T}^{0}, y_{T}^{1}\right) \in H_{\Gamma_{0}}^{1}(\Omega) \times L^{2}(\Omega)
$$

there exists a control $v \in L^{2}\left(0, T ; L^{2}\left(\Gamma_{1}\right)\right)$ such that the solution of (20)-(23) satisfies

$$
\begin{equation*}
y(T)=y_{T}^{0} \quad \text { and } \quad y^{\prime}(T)=y_{T}^{1} \quad \text { on } \quad \Omega \tag{25}
\end{equation*}
$$

(Compare with theorem 6.20.)
Proof. - As in chapter 4, we may assume that $y_{T}^{0}=y_{T}^{1}=0$. First consider the case

$$
a \equiv(n-1)(m \cdot \nu) /\left(2 R^{2}\right)
$$

and set $l:=(m \cdot \nu) / R$ for brevity.
Given $\left(u^{0}, u^{1}\right) \in \mathcal{H}=H_{\Gamma_{0}}^{1}(\Omega) \times L^{2}(\Omega)$ arbitrarily, first solve the problem

$$
\begin{gathered}
u^{\prime \prime}-\Delta u+q u=0 \quad \text { in } \Omega \times \mathbb{R}_{+}, \\
u=0 \quad \text { on } \Gamma_{0} \times \mathbb{R}_{+}, \\
\partial_{\nu} u+a u+l u^{\prime}=0 \quad \text { on } \Gamma_{1} \times \mathbb{R}_{+}, \\
u(0)=u^{0} \quad \text { and } \quad u^{\prime}(0)=u^{1} \quad \text { on } \Omega,
\end{gathered}
$$

then the problem

$$
\begin{gathered}
z^{\prime \prime}-\Delta z+q z=0 \quad \text { in } \Omega \times \mathbb{R}_{+}, \\
z=0 \quad \text { on } \Gamma_{0} \times \mathbb{R}_{+}, \\
\partial_{\nu} z+a z+l z^{\prime}=0 \quad \text { on } \Gamma_{1} \times \mathbb{R}_{+}, \\
z(0)=-u(T) \text { and } z^{\prime}(0)=u^{\prime}(T) \text { on } \Omega,
\end{gathered}
$$

and define

$$
\begin{gathered}
y(t):=u(t)+z(T-t), \quad t \in[0, T], \\
v(t):=-l\left(\left(u^{\prime}(t)+z^{\prime}(T-t)\right), \quad t \in[0, T] .\right.
\end{gathered}
$$

Then $v \in L^{2}\left(0, T ; L^{2}\left(\Gamma_{1}\right)\right)$ by remark 8.5 and $y$ satisfies (20), (21), (22), (25) and

$$
y(0)=u^{0}+z(T) \quad \text { and } \quad y^{\prime}(0)=u^{1}-z^{\prime}(T) \quad \text { on } \quad \Omega .
$$

Therefore it is sufficient to show that for any given $\left(y^{0}, y^{1}\right) \in \mathcal{H}$ there exists $\left(u^{0}, u^{1}\right) \in \mathcal{H}$ such that

$$
y^{0}=u^{0}+z(T) \quad \text { and } \quad y^{1}=u^{1}-z^{\prime}(T) \quad \text { on } \quad \Omega .
$$

In other words, it suffices to show that the linear map

$$
L: \mathcal{H} \rightarrow \mathcal{H}
$$

defined by

$$
L\left(u^{0}, u^{1}\right):=\left(u^{0}+z(T), u^{1}-z^{\prime}(T)\right)
$$

is onto. Since $L=I-K$ where the linear map

$$
K: \mathcal{H} \rightarrow \mathcal{H}
$$

is defined by

$$
K\left(u^{0}, u^{1}\right):=\left(-z(T), z^{\prime}(T)\right),
$$

it is sufficient to verify that $\|K\|<1$. Indeed, then $L$ is invertible with

$$
L^{-1}=I+K+K^{2}+K^{3}+\cdots
$$

Now a twofold application of theorem 8.6 gives

$$
\begin{aligned}
& \left\|K\left(u^{0}, u^{1}\right)\right\| \leq \exp \left(1-\left(1-Q_{1}\right) T /(2 R)\right)^{1 / 2}\left\|\left(-u(T), u^{\prime}(T)\right)\right\| \\
& \quad \leq \exp \left(1-\left(1-Q_{1}\right) T /(2 R)\right)\left\|\left(u^{0}, u^{1}\right)\right\|=: \gamma\left\|\left(u^{0}, u^{1}\right)\right\|
\end{aligned}
$$

and the result follows because $\gamma<1$ by hypothesis (24).

Now consider the general case. Using the above special case, for arbitrarily given

$$
\left(y^{0}, y^{1}\right),\left(y_{T}^{0}, y_{T}^{1}\right) \in H_{\Gamma_{0}}^{1}(\Omega) \times L^{2}(\Omega)
$$

there exists $v \in L^{2}\left(0, T ; L^{2}\left(\Gamma_{1}\right)\right)$ such that the solution of problem (20)-(23), with $a$ replaced by $a^{\prime}:=(n-1)(m \cdot \nu) /\left(2 R^{2}\right)$, satisfies (25). It follows by an obvious algebraic manipulation that the solution of the original problem also has this property if we apply the control $v+\left(a-a^{\prime}\right) y$ instead of $v$. Since this new control also belongs to $L^{2}\left(0, T ; L^{2}\left(\Gamma_{1}\right)\right)$ by the regularity of the solutions of (20)-(23) (cf. theorem 7.4), the proof is completed.

### 8.4. Uniform stabilization of the wave equation II

Let us return to the dissipative problem of sections 7.2 and 8.2 :

$$
\begin{gather*}
u^{\prime \prime}-\Delta u+q u=0 \quad \text { in } \quad \Omega \times \mathbb{R}_{+},  \tag{26}\\
u=0 \quad \text { on } \quad \Gamma_{0} \times \mathbb{R}_{+},  \tag{27}\\
\partial_{\nu} u+a u+l u^{\prime}=0 \quad \text { on } \quad \Gamma_{1} \times \mathbb{R}_{+},  \tag{28}\\
u(0)=u^{0} \quad \text { and } \quad u^{\prime}(0)=u^{1} \quad \text { on } \quad \Omega . \tag{29}
\end{gather*}
$$

Assume that there is a point $x^{0} \in \mathbb{R}^{n}$ such that

$$
\begin{equation*}
m \cdot \nu \leq 0 \quad \text { on } \quad \Gamma_{0} \quad \text { and } \quad \min _{\Gamma_{1}} m \cdot \nu>0 \tag{30}
\end{equation*}
$$

and assume that

$$
\begin{equation*}
\min _{\Gamma_{1}} a>0 \tag{31}
\end{equation*}
$$

and

$$
\begin{equation*}
\min _{\Gamma_{1}} l>0 . \tag{32}
\end{equation*}
$$

Furthermore, we continue to assume that

$$
\begin{equation*}
Q_{1}<1 \tag{33}
\end{equation*}
$$

Theorem 8.11. - Assume (30)-(33). Then there exist two positive constants $C, \omega$ such that for any given $\left(u^{0}, u^{1}\right) \in H_{\Gamma_{0}}^{1}(\Omega) \times L^{2}(\Omega)$ the solution of (26)-(29) satisfies the estimate

$$
\begin{equation*}
E(t) \leq C E(0) e^{-\omega t}, \quad \forall t \in \mathbb{R}_{+} . \tag{34}
\end{equation*}
$$

Proof. - Repeating the density argument of the proof of theorem 8.6 we may assume that $\left(u^{0}, u^{1}\right) \in \mathcal{D}(\mathcal{A})$.

By (30)-(32) we may write $a=(m \cdot \nu) b$ and $l=(m \cdot \nu) k$ with suitable positive functions $b, k \in C^{\infty}\left(\Gamma_{1}\right)$. Then lemmas 8.4 and 8.8 remain valid. Applying lemma 3.22 we deduce from identity (16) the inequality

$$
\begin{gather*}
2\left(1-Q_{1}\right) \int_{S}^{T} E \mathrm{dt} \leq 2 R_{1} E(S)+2 R_{1} E(T)  \tag{35}\\
+\int_{S}^{T} \int_{\Gamma_{1}}\left(u^{\prime}\right)^{2}-|\nabla u|^{2}+b u^{2}-\left(k u^{\prime}+b u\right) M u \mathrm{~d} \Gamma_{m} \mathrm{dt}
\end{gather*}
$$

(Compare with (18).) Since

$$
2\left(k u^{\prime}+b u\right)(m \cdot \nabla u) \leq|\nabla u|^{2}+2 R^{2} k^{2}\left(u^{\prime}\right)^{2}+R^{2} b^{2} u^{2}
$$

and

$$
2\left(k u^{\prime}+b u\right) u \leq k^{2}\left|u^{\prime}\right|^{2}+(b+1)|u|^{2},
$$

hence we deduce that

$$
\begin{align*}
& 2\left(1-Q_{1}\right) \int_{S}^{T} E \mathrm{dt} \leq c E(S)+c E(T)  \tag{36}\\
& \quad+c_{1} \int_{S}^{T} \int_{\Gamma_{1}}\left(u^{\prime}\right)^{2}+u^{2} \mathrm{~d} \Gamma \mathrm{dt}
\end{align*}
$$

here and in the sequel we shall denote by $c$ diverse constants, independent of the initial data and of $S, T$.

By (9) and (32) we have $E(T) \leq E(S)$ and

$$
\int_{S}^{T} \int_{\Gamma_{1}}\left(u^{\prime}\right)^{2} \mathrm{~d} \Gamma \mathrm{dt} \leq c E(S)
$$

Therefore we deduce from (36) that

$$
\begin{equation*}
2\left(1-Q_{1}\right) \int_{S}^{T} E \mathrm{dt} \leq c E(S)+c_{1} \int_{S}^{T} \int_{\Gamma_{1}} u^{2} \mathrm{~d} \Gamma \mathrm{dt} \tag{37}
\end{equation*}
$$

We shall eliminate the term $u^{2}$ on the right-hand side by using a method of Conrad and Rao [1]:

Lemma 8.12. - There exists a constant $c>0$ such that for every $\varepsilon \in(0,1)$ the solutions of (26)-(29) satisfy the inequality

$$
\begin{equation*}
\int_{S}^{T} \int_{\Gamma_{1}} u^{2} \mathrm{~d} \Gamma \mathrm{dt} \leq \frac{c}{\varepsilon} E(S)+\varepsilon \int_{S}^{T} E \mathrm{dt} . \tag{38}
\end{equation*}
$$

Proof. - We define (for each fixed $t \in \mathbb{R}_{+}$) $z=z(t)$ by

$$
\begin{equation*}
\Delta z=0 \quad \text { in } \quad \Omega \quad \text { and } \quad z=u(t) \quad \text { on } \quad \Gamma ; \tag{39}
\end{equation*}
$$

then we have the estimates

$$
\begin{equation*}
\int_{\Omega} z^{2} \mathrm{dx} \leq c \int_{\Gamma} u^{2} \mathrm{~d} \Gamma \leq c E \tag{40}
\end{equation*}
$$

by the elliptic regularity theory. Since

$$
\int_{\Omega} \nabla z \cdot \nabla(u-z) \mathrm{dx}=-\int_{\Omega}(\Delta z)(u-z) \mathrm{dx}+\int_{\Gamma}\left(\partial_{\nu} z\right)(u-z) \mathrm{d} \Gamma=0
$$

by (39), we have

$$
\begin{equation*}
\int_{\Omega} \nabla z \cdot \nabla u \mathrm{dx}=\int_{\Omega}|\nabla z|^{2} \mathrm{dx} \geq 0 . \tag{41}
\end{equation*}
$$

Since $z^{\prime}$ satisfies (46) with $u$ replaced by $u^{\prime}$, we also have

$$
\begin{equation*}
\int_{\Omega}\left|z^{\prime}\right|^{2} \mathrm{dx} \leq \int_{\Gamma}\left|u^{\prime}\right|^{2} \mathrm{~d} \Gamma \leq-c E^{\prime} \tag{42}
\end{equation*}
$$

(Here we use (30), (32) and lemma 8.4.)
Now we have

$$
\begin{gathered}
0=\int_{S}^{T} \int_{\Omega} z\left(u^{\prime \prime}-\Delta u\right) \mathrm{dx} \mathrm{dt} \\
=\left[\int_{\Omega} z u^{\prime} \mathrm{dx}\right]_{S}^{T}+\int_{S}^{T} \int_{\Omega}-z^{\prime} u^{\prime}+\nabla z \cdot \nabla u \mathrm{dx} \mathrm{dt}-\int_{S}^{T} \int_{\Gamma} z \partial_{\nu} u \mathrm{~d} \Gamma \mathrm{dt} \\
=\left[\int_{\Omega} z u^{\prime} \mathrm{dx}\right]_{S}^{T}+\int_{S}^{T} \int_{\Omega}-z^{\prime} u^{\prime}+\nabla z \cdot \nabla u \mathrm{dx} \mathrm{dt}+\int_{S}^{T} \int_{\Gamma_{1}} u\left(a u+l u^{\prime}\right) \mathrm{d} \Gamma \mathrm{dt} .
\end{gathered}
$$

Using (40)-(42) hence we deduce that

$$
\begin{gathered}
\int_{S}^{T} \int_{\Gamma_{1}} a u^{2} \mathrm{~d} \Gamma \mathrm{dt} \leq\left[\int_{\Omega} z u^{\prime} \mathrm{dx}\right]_{T}^{S}+\int_{S}^{T} \int_{\Omega} z^{\prime} u^{\prime} \mathrm{dx} \mathrm{dt}-\int_{S}^{T} \int_{\Gamma_{1}} l u u^{\prime} \mathrm{d} \Gamma \mathrm{dt} \\
\leq\|z(S)\|_{H}\left\|u^{\prime}(S)\right\|_{H}+\|z(T)\|_{H}\left\|u^{\prime}(T)\right\|_{H} \\
+\int_{S}^{T}\left\|z^{\prime}(t)\right\|_{H}\left\|u^{\prime}(t)\right\|_{H} \mathrm{dt}+c \int_{S}^{T} \int_{\Gamma_{1}}\left|u u^{\prime}\right| \mathrm{d} \Gamma \mathrm{dt} \\
\leq c E(S)+c E(T)+c \int_{S}^{T}\left(-E^{\prime}\right)^{1 / 2} E^{1 / 2} \mathrm{dt}+c \int_{S}^{T} \int_{\Gamma_{1}}\left|u u^{\prime}\right| \mathrm{d} \Gamma \mathrm{dt} .
\end{gathered}
$$

Fix $\alpha>0$ such that $a \geq 2 \alpha$ on $\Gamma_{1}$ (cf. (31)). Since the energy is nonincreasing, using (9), (32) we obtain that

$$
\begin{gathered}
2 \alpha \int_{S}^{T} \int_{\Gamma_{1}} u^{2} \mathrm{~d} \Gamma \mathrm{dt} \leq \int_{S}^{T} \int_{\Gamma_{1}} a u^{2} \mathrm{~d} \Gamma \mathrm{dt} \\
\leq c E(S)+\int_{S}^{T} \alpha \varepsilon E-\frac{c}{\alpha \varepsilon} E^{\prime} \mathrm{dt}+\int_{S}^{T} \int_{\Gamma_{1}} \alpha u^{2}+\frac{c}{\alpha}\left(u^{\prime}\right)^{2} \mathrm{~d} \Gamma \mathrm{dt} \\
\leq
\end{gathered}
$$

which implies (38).
We conclude from (37) and (38) that

$$
2\left(1-Q_{1}-c_{1} \varepsilon\right) \int_{S}^{T} E \mathrm{dt} \leq \frac{c}{\varepsilon} E(S), \quad \forall 0 \leq S<T<+\infty
$$

for every $0<\varepsilon<1$. By hypothesis (33) we may choose $\varepsilon$ such that $1-Q_{1}-\varepsilon>0$. Then the estimate (34) follows by applying theorem 8.1. $\square$

### 8.5. Strong stabilization. LaSalle's principle

We consider in this section the problem

$$
\begin{gather*}
u^{\prime \prime}-\Delta u+q u=0 \quad \text { in } \Omega \times \mathbb{R},  \tag{43}\\
u=0 \quad \text { on } \quad \Gamma_{0} \times \mathbb{R},  \tag{44}\\
\partial_{\nu} u+a u+l u^{\prime}=0 \quad \text { on } \Gamma_{1} \times \mathbb{R},  \tag{45}\\
u(0)=u^{0} \quad \text { and } \quad u^{\prime}(0)=u^{1} \quad \text { on } \Omega \tag{46}
\end{gather*}
$$

under weaker hypotheses as before. Assume that there is a point $x^{0} \in \mathbb{R}^{n}$ such that

$$
\begin{equation*}
m \cdot \nu \leq 0 \quad \text { on } \quad \Gamma_{0}, \quad m \cdot \nu \geq 0 \quad \text { on } \quad \Gamma_{1}, \tag{47}
\end{equation*}
$$

and that

$$
\begin{equation*}
l>0 \quad \text { on } \quad \Gamma_{+} . \tag{48}
\end{equation*}
$$

Theorem 8.13. - Assume (47) and (48). Then the solution of (43)-(46) satisfies

$$
\begin{equation*}
E(t) \rightarrow 0, \quad \forall\left(u^{0}, u^{1}\right) \in H_{\Gamma_{0}}^{1}(\Omega) \times L^{2}(\Omega) \tag{49}
\end{equation*}
$$

The hypotheses of this theorem are much weaker than those of theorem 8.11, but the conclusion is also weaker : in general there is no exponential energy decay.

Note that there is no hypothesis involving $Q_{1}$.
The theorem remains valid (with a simple modification of the proof) for the nonlinear problem (7.8)-(7.11) if the function $g$ is locally Lipschitz and if it satisfies the following two conditions :

$$
\begin{gathered}
g(x) \neq 0 \text { if } x \neq 0 \\
\exists c^{\prime}>0: \quad\left|g^{\prime}(x)\right| \leq 1+c^{\prime}|x|^{n /(n-2)}, \quad \forall x \in \mathbb{R} .
\end{gathered}
$$

(For this we need a nonlinear generalization of lemma 8.4; cf. lemma 9.7 in the following chapter.)

More general results are obtained in Chen and Wang [1], Lasiecka [1], [3], Zuazua [6].

Proof. - Using the density of $D(\mathcal{A})$ in $\mathcal{H}$ and the inequality (7.19) we may assume that

$$
\left(u^{0}, u^{1}\right) \in D(\mathcal{A})
$$

Fix $\left(u^{0}, u^{1}\right) \in D(\mathcal{A})$ arbitrarily. Then the set

$$
\left\{\left(u(t), u^{\prime}(t)\right): t \in \mathbb{R}_{+}\right\}
$$

is bounded in $H^{2}(\Omega) \times H^{1}(\Omega)$ by theorem 7.5 and therefore it is precompact in $\mathcal{H}$ by Rellich's theorem. In order to prove the relation $E(t) \rightarrow 0$ it suffices to show that if for some increasing sequence $t_{n} \rightarrow \infty$ the sequence $U\left(t_{n}\right)=\left(u\left(t_{n}\right), u^{\prime}\left(t_{n}\right)\right)$ converges in $\mathcal{H}$ to some point $\left(z^{0}, z^{1}\right)$, then $z^{0}=z^{1} \equiv 0$.

Set

$$
z_{n}(t):=u\left(t_{n}+t\right), \quad t \in \mathbb{R}_{+}, \quad n=1,2, \ldots,
$$

then $z_{n}$ is the solution of (43)-(46) with $\left(u^{0}, u^{1}\right)$ replaced by $\left(u\left(t_{n}\right), u^{\prime}\left(t_{n}\right)\right)$. It follows from inequality (7.19) of theorem 7.4 that the sequence $\left(z_{n}, z_{n}^{\prime}\right)$ is precompact in $L^{\infty}\left(\mathbb{R}_{+} ; \mathcal{H}\right)$. Extracting a subsequence if needed, we may assume that

$$
\left(z_{n}, z_{n}^{\prime}\right) \rightarrow\left(z, z^{\prime}\right) \quad \text { in } \quad L^{\infty}\left(\mathbb{R}_{+} ; \mathcal{H}\right)
$$

Clearly, $z$ satisfies (43)-(46) with $\left(u^{0}, u^{1}\right)$ replaced by $\left(z^{0}, z^{1}\right)$. We are going to show that $z \equiv 0$.

Observe that $E(z)$ is constant. Indeed, we have

$$
E(z ; t)=\lim _{n \rightarrow \infty} E\left(z_{n} ; t\right)=\lim _{n \rightarrow \infty} E\left(u ; t_{n}+t\right)=\lim _{s \rightarrow \infty} E(u ; s), \quad \forall t \in \mathbb{R}_{+} ;
$$

the last limit exists by the non-increasingness of the energy and it is independent of $t$. Using (9) and (48) hence we deduce that

$$
\begin{equation*}
z^{\prime} \equiv 0 \quad \text { on } \quad \Gamma_{1} \times \mathbb{R}_{+} . \tag{50}
\end{equation*}
$$

Now we may apply lemma 6.16 from chapter 6 . Indeed, it follows from (50) that $z$ satisfies (6.41)-(6.44) with $\left(u^{0}, u^{1}\right)$ replaced by $\left(z^{0}, z^{1}\right)$, and condition (6.49) is satisfied with any interval $J$. Applying the lemma we obtain $z \equiv 0$ and hence $z^{0}=z^{1} \equiv 0$.

The method of the above proof is due to LaSalle [1]; also see Haraux [5] for various applications of this method to nonlinear partial differential equations.

Remark 8.14. - Alternatively, theorem 8.12 could have been obtained by applying a general theorem of Benchimol [1], based on the decomposition theory of semigroups developed by Sz.-Nagy and Foias [1] and Foguel [1]. This theorem permits one to reduce the problem, in the linear case, to the study of the eigenvalues of $\mathcal{A}$. Generalizing some standard finite-dimensional results, it is sufficient to prove that $\mathcal{A}$ does not have purely imaginary eigenvalues. (Since $\mathcal{A}$ is monotone, it is already known that no eigenvalue may have negative real part.) See e.g. Lagnese [2] for the application of Benchimol's theorem to the stabilization of some plate models.

### 8.6. Uniform stabilization of the wave equation III

We shall improve here theorem 8.10 of section 8.4 concerning the problem

$$
\begin{gather*}
u^{\prime \prime}-\Delta u+q u=0 \quad \text { in } \Omega \times \mathbb{R}_{+},  \tag{51}\\
u=0 \quad \text { on } \quad \Gamma_{0} \times \mathbb{R}_{+},  \tag{52}\\
\partial_{\nu} u+a u+l u^{\prime}=0 \quad \text { on } \quad \Gamma_{1} \times \mathbb{R}_{+},  \tag{53}\\
u(0)=u^{0} \quad \text { and } \quad u^{\prime}(0)=u^{1} \quad \text { on } \quad \Omega, \tag{54}
\end{gather*}
$$

by applying the strong stabilization result obtained in the preceding section; we shall eliminate hypothesis (33) on $Q_{1}$. The method of proof, introduced in Komornik and Rao [1], plays a similar role in stabilization problems to the method of chapter 5 in exact controllability problems.

Assume that there is a point $x^{0} \in \mathbb{R}^{n}$ such that

$$
\begin{equation*}
m \cdot \nu \leq 0 \quad \text { on } \quad \Gamma_{0}, \quad \min _{\Gamma_{1}} m \cdot \nu>0, \tag{55}
\end{equation*}
$$

and assume that

$$
\begin{equation*}
\min _{\Gamma_{1}} a>0, \quad \min _{\Gamma_{1}} l>0 \tag{56}
\end{equation*}
$$

Theorem 8.15. - Assume (55), (56). There exist two positive constants $C, \omega$ such that for every given $\left(u^{0}, u^{1}\right) \in H_{\Gamma_{0}}^{1}(\Omega) \times L^{2}(\Omega)$ the solution of (51)-(54) satisfies the estimate

$$
\begin{equation*}
E(t) \leq C E(0) e^{-\omega t}, \quad \forall t \in \mathbb{R}_{+} \tag{57}
\end{equation*}
$$

The proof, given in Komornik and Rao [1], is based on the following useful general result of Gibson [1] which is admitted here without proof.

Theorem 8.16. - Let $\mathcal{A}$ be a maximal monotone linear operator in a Hilbert space $\mathcal{H}$ and assume that the solutions of the problem

$$
\begin{equation*}
U^{\prime}+\mathcal{A} U=0 \quad \text { in } \quad \mathbb{R}_{+}, \quad U(0)=U^{0} \tag{58}
\end{equation*}
$$

are strongly stable :

$$
\begin{equation*}
\|U(t)\|_{\mathcal{H}} \rightarrow 0 \text { as } t \rightarrow+\infty, \quad \forall U^{0} \in \mathcal{H} \tag{59}
\end{equation*}
$$

Assume that there exists a compact linear operator $\mathcal{B}$ in $\mathcal{H}$ such that the solutions of the problem

$$
\begin{equation*}
V^{\prime}+\mathcal{A} V+\mathcal{B} V=0 \quad \text { in } \quad \mathbb{R}_{+}, \quad V(0)=V^{0} \tag{60}
\end{equation*}
$$

are uniformly exponentially stable : there exist two positive constants $C_{1}, \omega_{1}$ such that

$$
\begin{equation*}
\|V(t)\|_{\mathcal{H}} \leq C_{1}\left\|V^{0}\right\|_{\mathcal{H}} e^{-\omega_{1} t}, \quad \forall t \in \mathbb{R}_{+}, \quad \forall V^{0} \in \mathcal{H} \tag{61}
\end{equation*}
$$

Then the solutions of (58) are also uniformly exponentially stable : there exist two positive constants $C, \omega$ such that

$$
\begin{equation*}
\|U(t)\|_{\mathcal{H}} \leq C\left\|U^{0}\right\|_{\mathcal{H}} e^{-\omega t}, \quad \forall t \in \mathbb{R}_{+}, \quad \forall U^{0} \in \mathcal{H} \tag{62}
\end{equation*}
$$

Gibson's theorem generalized an earlier theorem of Russell [1].
Note that the operator $\mathcal{A}+\mathcal{B}$ is not necessarily maximal monotone; nevertheless, the problem (60) has a unique solution $V \in C\left(\mathbb{R}_{+} ; \mathcal{H}\right)$ for every
$V^{0} \in \mathcal{H}$. Indeed, the following more general holds true (see e.g. Brézis [1]) : if $\mathcal{A}$ is a (not necessarily linear) maximal monotone operator in a Hilbert space $\mathcal{H}$ and if $\mathcal{B}$ is a (non necessarily linear) Lipschitz continuous operator in $\mathcal{H}$, then problem (60) has a unique solution $V \in C\left(\mathbb{R}_{+} ; \mathcal{H}\right)$ for every $V^{0} \in \overline{D(\mathcal{A})}$.

Proof of theorem 8.15. - By hypotheses (55), (56) the conditions (47), (48) of theorem 8.13 are satisfied. Consequently, the solutions of problem (51)-(54) (i.e. of (58)) are strongly stable.

Furthermore, the formula

$$
\mathcal{B} V:=\left(0,-q V_{1}\right), \quad V=\left(V_{1}, V_{2}\right) \in \mathcal{H}
$$

defines a compact linear operator in $\mathcal{H}$ (because of the compactness of the imbedding $\left.H^{1}(\Omega) \subset L^{2}(\Omega)\right)$. The solutions of the corresponding problem (60) are uniformly exponentially stable. Indeed, (60) coincides with the problem (26)-(29) where we replace $q$ by zero and therefore we may apply theorem 8.11 : conditions (30)-(32) follow from hypotheses (55), (56), and (33) is satisfied because $Q_{1}=0$.

We conclude by applying theorem 8.16.

### 8.7. Uniform stabilization of Maxwell's equations

We consider here the problem

$$
\begin{gather*}
E^{\prime}-\operatorname{curl} H=H^{\prime}+\operatorname{curl} E \quad \text { in } \Omega \times \mathbb{R}_{+}  \tag{63}\\
\operatorname{div} E=\operatorname{div} H=0 \quad \text { in } \Omega \times \mathbb{R}_{+}  \tag{64}\\
\nu \times(E \times \nu+H)=0 \quad \text { on } \Gamma \times \mathbb{R}_{+}  \tag{65}\\
E(0)=E^{0} \quad \text { and } \quad H(0)=H^{0} \quad \text { in } \Omega \tag{66}
\end{gather*}
$$

where $\Omega$ is a bounded open domain in $\mathbb{R}^{3}$ having a boundary $\Gamma$ of class $C^{1}$ and $\nu$ denotes the outward unit normal vector to $\Gamma$.

Let us introduce the Hilbert space

$$
\begin{aligned}
\mathcal{H}:= & \left\{(E, H) \in L^{2}(\Omega)^{6} \mid \operatorname{div} E=\operatorname{div} H=0 \text { in } \Omega\right\}, \\
& \|(E, H)\|_{\mathcal{H}}:=\left(\frac{1}{2} \int_{\Omega}|E|^{2}+|H|^{2} \mathrm{dx}\right)^{1 / 2}
\end{aligned}
$$

and let us define in $\mathcal{H}$ a linear operator by setting

$$
D(\mathcal{A}):=\left\{(E, H) \in H^{1}(\Omega)^{6} \cap \mathcal{H} \mid \nu \times(E \times \nu+H)=0 \text { on } \Gamma\right\}
$$

and

$$
\mathcal{A}(E, H):=(-\operatorname{curl} H, \operatorname{curl} E) .
$$

Then we can rewrite the problem (63)-(66) in the following form :

$$
(E, H)^{\prime}+\mathcal{A}(E, H)=0 \quad \text { in } \quad \mathbb{R}_{+}, \quad(E, H)(0)=\left(E^{0}, H^{0}\right)
$$

It was proved by Barucq and Hanouzet [2] that the operator $\mathcal{A}$ is maximal monotone in $\mathcal{H}$. We admit this result here. Applying the Hille-Yosida theorem it follows that for every given $\left(E^{0}, H^{0}\right) \in \mathcal{H}$ the problem (63)-(66) has a unique (mild) solution

$$
(E, H) \in C\left(\mathbb{R}_{+} ; \mathcal{H}\right)
$$

moreover, for $\left(E^{0}, H^{0}\right) \in D(\mathcal{A})$ we have

$$
(E, H) \in C\left(\mathbb{R}_{+} ; D(\mathcal{A})\right) \cap C^{1}\left(\mathbb{R}_{+} ; \mathcal{H}\right)
$$

whence in particular

$$
\begin{equation*}
(E, H) \in C\left(\mathbb{R}_{+} ; H^{1}(\Omega)^{6}\right) \cap C^{1}\left(\mathbb{R}_{+} ; L^{2}(\Omega)^{6}\right) \tag{67}
\end{equation*}
$$

We define the energy of the solutions by

$$
\begin{equation*}
\mathcal{E}(t):=\frac{1}{2} \int_{\Omega}|E(t)|^{2}+|H(t)|^{2} \mathrm{dx}, \quad t \in \mathbb{R}_{+} . \tag{68}
\end{equation*}
$$

Now assume that $\Omega$ is strictly star-shaped with respect to the origin :

$$
\begin{equation*}
x \cdot \nu(x)>0 \text { for all } x \in \Gamma \tag{69}
\end{equation*}
$$

and set

$$
\begin{gather*}
R:=\sup _{x \in \Omega}|x|  \tag{70}\\
R_{1}:=\max _{x \in \Gamma} \frac{(x \cdot \nu(x))^{2}+|x|^{2}}{2 x \cdot \nu(x)} . \tag{71}
\end{gather*}
$$

Considering a point $x \in \Gamma$ with $|x|=R$ one can readily verifiy that $R_{1} \geq R$. If $\Omega$ is a ball centered at the origin, then $R$ is equal to its radious and $R_{1}=R$.

We shall prove the
Theorem 8.17. - Assume (69). Then for any given $\left(E^{0}, H^{0}\right) \in \mathcal{H}$ the solution of the problem (63)-(66) satisfies the energy estimates

$$
\begin{equation*}
\mathcal{E}(t) \leq \mathcal{E}(0) e^{1-\left(t /\left(R+R_{1}\right)\right)}, \quad \forall t \geq 0 \tag{72}
\end{equation*}
$$

Applying Russell's principle one can deduce from theorem 8.17 an exact controllability result which improves some earlier theorems of Russell [3] and Lagnese [3]. Moreover, one can show that in the special case of the ball $\Omega=\left\{x \in \mathbb{R}^{3}| | x \mid<R\right\}$ (then $T_{0}=2 R$ ) these results are optimal. We refer to Komornik [13] for the proof of these results.

Turning to the proof of the theorem first we note that it suffices to prove the estimates (72) for the case of smooth initial data $\left(E^{0}, H^{0}\right) \in D(\mathcal{A})$ : the general case then follows easily by density. Henceforth we only consider such solutions ; in this case the regularity (67) of the solution is sufficient to justify the computations of this section.

We need some lemmas. The first shows in particular that the energy is non-increasing.

Lemma 8.18. - The solution of the problem (63)-(66) satisfies the energy equalities

$$
\begin{equation*}
\mathcal{E}(S)-\mathcal{E}(T)=\int_{S}^{T} \int_{\Gamma}\left|E_{\tau}\right|^{2} \mathrm{~d} \Gamma \mathrm{dt}=\int_{S}^{T} \int_{\Gamma}\left|H_{\tau}\right|^{2} \mathrm{~d} \Gamma \mathrm{dt} \tag{73}
\end{equation*}
$$

for all $0 \leq S<T<+\infty$, where $E_{\tau}, H_{\tau}$ denote the tangential components of $E, \Gamma$.

Proof. - Applying Green's formula we deduce easily from (63) and (68) that

$$
\begin{equation*}
\mathcal{E}^{\prime}(t) \equiv-\int_{\Gamma}(E(t) \times H(t)) \cdot \nu \mathrm{d} \Gamma, \quad t \geq 0 \tag{74}
\end{equation*}
$$

It remains to show that

$$
\begin{equation*}
(E \times H) \cdot \nu=\left|E_{\tau}\right|^{2}=\left|H_{\tau}\right|^{2} \tag{75}
\end{equation*}
$$

at every fixed point $x \in \Gamma$. Let us choose the system of coordinates such that $\nu=(0,0,1)$ at this point. Then

$$
\begin{equation*}
\nu \times(E \times \nu+H)=\left(E_{1}-H_{2}, E_{2}+H_{1}, 0\right) \tag{76}
\end{equation*}
$$

Using (65) hence we conclude that $E_{1}=H_{2}, E_{2}=-H_{1}$ and therefore

$$
(E \times H) \cdot \nu=E_{1} H_{2}-E_{2} H_{1}=E_{1}^{2}+E_{2}^{2}=H_{1}^{2}+H_{2}^{2},
$$

which is just another form of (73).
Remark 8.19. - The formula (76) expresses the geometric meaning of the boundary condition (65) : $H_{\tau}$ is obtained from $E_{\tau}$ by a rotation of angle $\pi / 2$ in the positive direction in the tangent plane.

The main tool in our proof of theorem 8.17 is the following identity. In what follows we shall denote by $m$ the identity mapping $m(x) \equiv x, x \in \mathbb{R}^{3}$.

Lemma 8.20. - The solution of (63)-(66) satisfies the identity

$$
\begin{gather*}
\int_{S}^{T} \int_{\Omega}|E|^{2}+|H|^{2} \mathrm{dx} \mathrm{dt}=\left[2 \int_{\Omega}(E \times H) \cdot m \mathrm{dx}\right]_{S}^{T} \\
+\int_{S}^{T} \int_{\Gamma}(m \cdot \nu)\left(|E|^{2}+|H|^{2}\right)-2(m \cdot E)(\nu \cdot E)-2(m \cdot H)(\nu \cdot H) \mathrm{d} \Gamma \mathrm{dt} \tag{77}
\end{gather*}
$$

for all $0 \leq S<T<+\infty$.
Proof. - The identity (77) will be obtained by the multiplier method. Let us write the equations (63) explicitly. Putting $E=\left(E_{1}, E_{2}, E_{3}\right)$, $H=\left(H_{1}, H_{2}, H_{3}\right)$ and writing for brevity $f_{, j}=\partial_{j} f$ we have

$$
\begin{align*}
& E_{1}^{\prime}=H_{3,2}-H_{2,3} \quad \text { in } \quad \Omega \times(0,+\infty),  \tag{78}\\
& E_{2}^{\prime}=H_{1,3}-H_{3,1} \quad \text { in } \Omega \times(0,+\infty),  \tag{79}\\
& E_{3}^{\prime}=H_{2,1}-H_{1,2} \quad \text { in } \Omega \times(0,+\infty),  \tag{80}\\
& H_{1}^{\prime}=E_{2,3}-E_{3,2} \quad \text { in } \Omega \times(0,+\infty),  \tag{81}\\
& H_{2}^{\prime}=E_{3,1}-E_{1,3} \quad \text { in } \Omega \times(0,+\infty),  \tag{82}\\
& H_{3}^{\prime}=E_{1,2}-E_{2,1} \quad \text { in } \quad \Omega \times(0,+\infty) . \tag{83}
\end{align*}
$$

Using (78) and (82) we have

$$
\begin{aligned}
& 2\left(E_{1} H_{2} m_{3}\right)^{\prime}=2\left(H_{3,2}-H_{2,3}\right) H_{2} m_{3}+2 E_{1}\left(E_{3,1}-E_{1,3}\right) m_{3} \\
& \quad=2 H_{3,2} H_{2} m_{3}-m_{3}\left(H_{2}^{2}\right)_{3}+2 E_{1} E_{3,1} m_{3}-m_{3}\left(E_{1}^{2}\right)_{3} .
\end{aligned}
$$

Analogously, using (79) and (81) we have

$$
\begin{aligned}
& 2\left(E_{2} H_{1} m_{3}\right)^{\prime}=2\left(H_{1,3}-H_{3,1}\right) H_{1} m_{3}+2 E_{2}\left(E_{2,3}-E_{3,2}\right) m_{3} \\
& \quad=m_{3}\left(H_{1}^{2}\right)_{3}-2 H_{3,1} H_{1} m_{3}+m_{3}\left(E_{2}^{2}\right)_{3}-2 E_{2} E_{3,2} m_{3} .
\end{aligned}
$$

Integrating by parts their difference in $\Omega \times(S, T)$ and writing $\nu=\left(\nu_{1}, \nu_{2}, \nu_{3}\right)$ we obtain

$$
\begin{gathered}
{\left[2 \int_{\Omega} E_{1} H_{2} m_{3}-E_{2} H_{1} m_{3} \mathrm{dx}\right]_{S}^{T}} \\
=\int_{S}^{T} \int_{\Omega}\left(H_{1}^{2}+H_{2}^{2}+E_{1}^{2}+E_{2}^{2}\right)-2 m_{3} H_{3}\left(H_{1,1}+H_{2,2}\right) \\
-2 m_{3} E_{3}\left(E_{1,1}+E_{2,2}\right) \mathrm{dx} \mathrm{dt} \\
\int_{S}^{T} \int_{\Gamma}-m_{3} \nu_{3}\left(H_{1}^{2}+H_{2}^{2}+E_{1}^{2}+E_{2}^{2}\right)+2 m_{3} \nu_{2}\left(H_{2} H_{3}+E_{2} E_{3}\right) \\
+2 m_{3} \nu_{1}\left(H_{1} H_{3}+E_{1} E_{3}\right) \mathrm{d} \Gamma \mathrm{dt} .
\end{gathered}
$$

Next we observe that using (64) we have

$$
\begin{aligned}
& \int_{\Omega}-2 m_{3} H_{3}\left(H_{1,1}+H_{2,2}\right) \mathrm{dx}=\int_{\Omega} 2 m_{3} H_{3} H_{3,3} \mathrm{dx} \\
& =\int_{\Omega} m_{3}\left(H_{3}^{2}\right)_{3} \mathrm{dx}=\int_{\Gamma} m_{3} \nu_{3} H_{3}^{2} \mathrm{~d} \Gamma-\int_{\Omega} H_{3}^{2} \mathrm{dx}
\end{aligned}
$$

and analogously

$$
\int_{\Omega}-2 m_{3} E_{3}\left(E_{1,1}+E_{2,2}\right) \mathrm{dx}=\int_{\Gamma} m_{3} \nu_{3} E_{3}^{2} \mathrm{~d} \Gamma-\int_{\Omega} E_{3}^{2} \mathrm{dx}
$$

Using these equalities the above identity may be rewritten as

$$
\begin{aligned}
& {\left[2 \int_{\Omega} E_{1} H_{2} m_{3}-E_{2} H_{1} m_{3} \mathrm{dx}\right]_{S}^{T}=\int_{S}^{T} \int_{\Omega} H_{1}^{2}+H_{2}^{2}-H_{3}^{2}+E_{1}^{2}+E_{2}^{2}-E_{3}^{2} \mathrm{dx} \mathrm{dt}} \\
& \quad+\int_{S}^{T} \int_{\Gamma}-m_{3} \nu_{3}\left(|E|^{2}+|H|^{2}\right)+2 m_{3} H_{3}(\nu \cdot H)+2 m_{3} E_{3}(\nu \cdot E) \mathrm{d} \Gamma \mathrm{dt}
\end{aligned}
$$

Two analogous identities may be obtained by cyclical permutation of the indices $1,2,3$. Summing the three identities we obtain that

$$
\begin{gathered}
{\left[2 \int_{\Omega}(E \times H) \cdot m \mathrm{dx}\right]_{S}^{T}=\int_{S}^{T} \int_{\Omega}|E|^{2}+|H|^{2} \mathrm{dx} \mathrm{dt}} \\
+\int_{S}^{T} \int_{\Gamma}-(m \cdot \nu)\left(|E|^{2}+|H|^{2}\right)+2(m \cdot H)(\nu \cdot H)+2(m \cdot E)(\nu \cdot E) \mathrm{d} \Gamma \mathrm{dt}
\end{gathered}
$$

and this is equivalent to the identity (77).
Observe that we did not use the boundary condition (65) in the proof of the preceding lemma. Thus the identity (77) remains valid for every function $(E, H)$ satisfying (63), (64) and (67). Now we shall use the boundary condition (65) in order to majorize the boundary integral in (77).

Lemma 8.21. - We have

$$
\begin{gather*}
(m \cdot \nu)\left(|E|^{2}+|H|^{2}\right)-2(m \cdot E)(\nu \cdot E)-2(m \cdot H)(\nu \cdot H) \\
\leq R_{1}\left(\left|E_{\tau}\right|^{2}+\left|H_{\tau}\right|^{2}\right) \tag{84}
\end{gather*}
$$

on $\Gamma$.
Proof. - Putting for brevity $E_{\nu}:=E \cdot \nu$ and $H_{\nu}:=H \cdot \nu$ we have

$$
E=E_{\tau}+E_{\nu} \nu, \quad H=H_{\tau}+H_{\nu} \nu
$$

and the left-hand side of (2.15) may be written as

$$
(m \cdot \nu)\left(\left|E_{\tau}\right|^{2}+\left|H_{\tau}\right|^{2}-E_{\nu}^{2}-H_{\nu}^{2}\right)-2\left(m_{\tau} \cdot E_{\tau}\right) E_{\nu}-2\left(m_{\tau} \cdot H_{\tau}\right) H_{\nu}
$$

Since we have obviously

$$
-2\left(m_{\tau} \cdot E_{\tau}\right) E_{\nu} \leq(m \cdot \nu) E_{\nu}^{2}+\frac{\left(m_{\tau} \cdot E_{\tau}\right)^{2}}{m \cdot \nu}
$$

and

$$
-2\left(m_{\tau} \cdot H_{\tau}\right) H_{\nu} \leq(m \cdot \nu) H_{\nu}^{2}+\frac{\left(m_{\tau} \cdot H_{\tau}\right)^{2}}{m \cdot \nu}
$$

the left-hand side of (84) is less than or equal to

$$
\begin{equation*}
(m \cdot \nu)\left(\left|E_{\tau}\right|^{2}+\left|H_{\tau}\right|^{2}\right)+\frac{\left(m_{\tau} \cdot E_{\tau}\right)^{2}+\left(m_{\tau} \cdot H_{\tau}\right)^{2}}{m \cdot \nu} . \tag{85}
\end{equation*}
$$

It follows from remark 8.19 that $E_{\tau}, H_{\tau}$ are orthogonal and they have the same length. Therefore the expression (85) is equal to

$$
\left((m \cdot \nu)+\frac{\left|m_{\tau}\right|^{2}}{2 m \cdot \nu}\right)\left(\left|E_{\tau}\right|^{2}+\left|H_{\tau}\right|^{2}\right)=\frac{(m \cdot \nu)^{2}+|m|^{2}}{2 m \cdot \nu}\left(\left|E_{\tau}\right|^{2}+\left|H_{\tau}\right|^{2}\right)
$$

Using (71) hence (84) follows.
Now we are ready to complete the proof of the theorem. Using (73) and (84) the last integral in the identity (77) is less than or equal to

$$
2 R_{1}(\mathcal{E}(S)-\mathcal{E}(T))
$$

Furthermore, we have obviously

$$
\left|2 \int_{\Omega} E \cdot(m \times H) \mathrm{dx}\right| \leq 2 \int_{\Omega}|m||E||H| \mathrm{dx} \leq 2 R \mathcal{E}
$$

Therefore we deduce from (77) the following inequality :

$$
\int_{S}^{T} \mathcal{E}(t) \mathrm{dt} \leq R(\mathcal{E}(S)+\mathcal{E}(T))+R_{1}(\mathcal{E}(S)-\mathcal{E}(T))
$$

Since $R_{1} \geq R$, the right-hand side is less than or equal to $\left(R+R_{1}\right) \mathcal{E}(S)$. Letting $T \rightarrow+\infty$ we obtain that

$$
\begin{equation*}
\int_{S}^{+\infty} \mathcal{E}(t) \mathrm{dt} \leq\left(R+R_{1}\right) \mathcal{E}(S), \quad \forall S \geq 0 \tag{86}
\end{equation*}
$$

Applying theorem 8.1 hence the estimate (72) follows.

## 9. Nonlinear stabilization

The aim of this chapter is to obtain polynomial energy decay estimates for the solutions of the wave equation and of a plate model by the application of suitable boundary feedbacks.

As before, we shall use for any given $x^{0} \in \mathbb{R}^{n}$ the notation

$$
\begin{gathered}
m(x)=x-x^{0}, \quad x \in \mathbb{R}^{n}, \\
R=R\left(x^{0}\right)=\sup \left\{\left|x-x^{0}\right|: x \in \Omega\right\}, \\
\mathrm{d} \Gamma_{m}=(m \cdot \nu) \mathrm{d} \Gamma .
\end{gathered}
$$

### 9.1. A nonlinear integral inequality

In this section we give a nonlinear generalization of theorem 8.1, which improves some earlier results of Haraux [2] and Lagnese [1].

Theorem 9.1. - Let $E: \mathbb{R}_{+} \rightarrow \mathbb{R}_{+}\left(\mathbb{R}_{+}:=[0,+\infty)\right)$ be a non-increasing function and assume that there are two constants $\alpha>0$ and $T>0$ such that

$$
\begin{equation*}
\int_{t}^{\infty} E^{\alpha+1}(s) \mathrm{ds} \leq T E(0)^{\alpha} E(t), \quad \forall t \in \mathbb{R}_{+} \tag{1}
\end{equation*}
$$

Then we have

$$
\begin{equation*}
E(t) \leq E(0)\left(\frac{T+\alpha t}{T+\alpha T}\right)^{-1 / \alpha}, \quad \forall t \geq T \tag{2}
\end{equation*}
$$

Note that inequality (2) is also satisfied for $0 \leq t<T$ : then it follows from the obvious inequality $E(t) \leq E(0)$.

Observe that letting $\alpha \rightarrow 0$ in this theorem we obtain theorem 8.1.
Proof. - If $E(0)=0$, then $E \equiv 0$ and there is nothing to prove. Otherwise, replacing the function $E$ by the function $E / E(0)$ we may assume that $E(0)=1$ and we have to prove the following estimation holds :

$$
\begin{equation*}
E(t) \leq\left(\frac{T+\alpha t}{T+\alpha T}\right)^{-1 / \alpha}, \quad \forall t \geq T \tag{3}
\end{equation*}
$$

Introduce the function

$$
F: \mathbb{R}_{+} \rightarrow \mathbb{R}_{+}, \quad F(t)=\int_{t}^{\infty} E^{\alpha+1} d s
$$

It is non-increasing and locally absolutely continuous. Differentiating and using (1) we find that

$$
-F^{\prime} \geq T^{-\alpha-1} F^{\alpha+1} \quad \text { a.e. in } \quad(0,+\infty)
$$

whence

$$
\left(F^{-\alpha}\right)^{\prime} \geq \alpha T^{-\alpha-1} \quad \text { a.e. in } \quad(0, B), \quad B:=\sup \{t: E(t)>0\} .
$$

(Observe that $F^{-\alpha}(t)$ is defined for $t<B$.) Integrating in $[0, s]$ we obtain that

$$
\left(F(s)^{-\alpha}-F(0)^{-\alpha}\right) \geq \alpha T^{-\alpha-1} s \quad \text { for every } \quad s \in[0, B),
$$

whence

$$
\begin{equation*}
F(s) \leq\left(F(0)^{-\alpha}+\alpha T^{-\alpha-1} s\right)^{-1 / \alpha} \quad \text { for every } \quad s \in[0, B) \tag{4}
\end{equation*}
$$

Since $F(s)=0$ if $s \geq B$, this inequality holds in fact for every $s \in \mathbb{R}_{+}$. Since $F(0) \leq T E(0)^{\alpha+1}=T$ by (1), the right-hand side of (4) is less than equal to

$$
\left(T^{-\alpha}+\alpha T^{-\alpha-1} s\right)^{-1 / \alpha}=T^{(\alpha+1) / \alpha}(T+\alpha s)^{-1 / \alpha} .
$$

On the other hand, $E$ being nonnegative and non-increasing, the left-hand side of (4) may be estimated as follows :
$F(s)=\int_{s}^{+\infty} E^{\alpha+1} \mathrm{dt} \geq \int_{s}^{T+(\alpha+1) s} E^{\alpha+1} \mathrm{dt} \geq(T+\alpha s) E(T+(\alpha+1) s)^{\alpha+1}$.
Therefore we deduce from (4) the estimate

$$
(T+\alpha s) E(T+(\alpha+1) s)^{\alpha+1} \leq T^{(\alpha+1) / \alpha}(T+\alpha s)^{-1 / \alpha},
$$

whence

$$
E(T+(\alpha+1) s) \leq\left(1+\frac{\alpha s}{T}\right)^{-1 / \alpha}, \quad \forall s \geq 0
$$

Choosing $t:=T+(\alpha+1) s$ hence the inequality (3) follows.
Remark 9.2. - The theorem is optimal in the following sense : given $\alpha>0, T>0, C>0$ and $t^{\prime} \geq T$ arbitrarily, there exists a non-increasing function $E: \mathbb{R}_{+} \rightarrow \mathbb{R}_{+}$satisfying (1) and such that

$$
E(0)=C \quad \text { and } \quad E\left(t^{\prime}\right)=E(0)\left(\frac{T+\alpha t^{\prime}}{T+\alpha T}\right)^{-1 / \alpha}
$$

We leave to the reader to verify that the following example has these properties :

$$
E(t):= \begin{cases}C\left(1+\alpha C^{-\alpha} t / T\right)^{-1 / \alpha}, & \text { if } 0 \leq t \leq t^{\prime \prime}  \tag{5}\\ C(1+\alpha)^{1 / \alpha}\left(1+\alpha C^{-\alpha} t^{\prime} / T\right)^{-1 / \alpha}, & \text { if } t^{\prime \prime} \leq t \leq t^{\prime} \\ 0, & \text { if } t>t^{\prime}\end{cases}
$$

where $t^{\prime \prime}=\left(t^{\prime}-T C^{\alpha}\right) /(\alpha+1)$. Let us also note that for $t<T$ we cannot state more than the trivial estimate $E(t) \leq E(0)$. Indeed, for any given $\alpha>0$, $T>0, C>0$ and $t^{\prime}<T$ the function

$$
E(t):= \begin{cases}C, & \text { if } 0 \leq t \leq T  \tag{6}\\ 0, & \text { if } t>T\end{cases}
$$

satisfies (1) and $E\left(t^{\prime}\right)=E(0)=C$.
Remark 9.3. - Assume that $E$ is also continuous. Then the inequalities (3) are strict ; in particular, $E(T)<E(0)$. See Komornik [9], [12] for this result, for a detailed study of integral inequalities of type (1) (also for $\alpha<0$ ) and for the study of closely related differential inequalities.

### 9.2. Uniform stabilization of the wave equation I

Fix a point $x^{0} \in \mathbb{R}^{n}$ and consider the nonlinear problem

$$
\begin{gather*}
u^{\prime \prime}-\Delta u=0 \quad \text { in } \Omega \times \mathbb{R}_{+},  \tag{7}\\
u=0 \quad \text { on } \Gamma_{0} \times \mathbb{R}_{+},  \tag{8}\\
\partial_{\nu} u+(m \cdot \nu) g\left(u^{\prime}\right)=0 \quad \text { on } \quad \Gamma_{1} \times \mathbb{R}_{+},  \tag{9}\\
u(0)=u^{0} \quad \text { and } \quad u^{\prime}(0)=u^{1} \quad \text { on } \quad \Omega, \tag{10}
\end{gather*}
$$

a particular case of the problem in section 7.2.
We assume that $g: \mathbb{R} \rightarrow \mathbb{R}$ is a non-decreasing, continuous function satisfying $g(0)=0$ and we assume that

$$
\begin{gather*}
n \geq 3  \tag{11}\\
\Gamma_{0} \neq \emptyset \quad \text { and } \overline{\Gamma_{0}} \cap \overline{\Gamma_{1}}=\emptyset  \tag{12}\\
m \cdot \nu \leq 0 \quad \text { on } \quad \Gamma_{0} \quad \text { and } \quad m \cdot \nu \geq 0 \quad \text { on } \quad \Gamma_{1} . \tag{13}
\end{gather*}
$$

(More general situations are considered in Zuazua [5], Conrad and Rao [1] and in Komornik [10].)

We recall that the energy is defined by

$$
E:=\frac{1}{2} \int_{\Omega}\left|u^{\prime}\right|^{2}+|\nabla u|^{2} \mathrm{dx}
$$

In this section we prove the
Theorem 9.4. - Assume (11)-(13). Assume that there exist $p>1$ and four constants $c_{1}, c_{2}, c_{3}, c_{4}>0$ such that

$$
\begin{equation*}
c_{1}|x|^{p} \leq|g(x)| \leq c_{2}|x|^{1 / p} \quad \text { if } \quad|x| \leq 1 \tag{14}
\end{equation*}
$$

and

$$
\begin{equation*}
c_{3}|x| \leq|g(x)| \leq c_{4}|x| \quad \text { if } \quad|x|>1 \tag{15}
\end{equation*}
$$

Then for any given $\left(u_{0}, u_{1}\right) \in H_{\Gamma_{0}}^{1}(\Omega) \times L^{2}(\Omega)$ the solution of $(7)-(10)$ satisfies the estimates

$$
\begin{equation*}
E(t) \leq C t^{2 /(1-p)}, \forall t>0 \tag{16}
\end{equation*}
$$

with a constant $C$ only depending on the initial energy $E(0)$ (and in a continuous way).

Theorem 9.4 improves some earlier results of Zuazua [5] and Conrad and Rao [1]. Let us note that the dependence of the constant $C$ on $E(0)$ was studied in detail by Carpio [1] and later by Souplet [1] and Kouémou [1].

Remark 9.5. - A similar result holds for $p=1$ : then (16) is replaced by

$$
E(t) \leq C E(0) e^{-\omega t}, \forall t>0
$$

with two positive constants $C, \omega$, independent of the initial data. Thus we may recover certain results of the preceding chapter.

First we shall prove the theorem under the additional hypothesis that

$$
g \quad \text { is globally Lipschitz continuous. }
$$

This assumption will be removed at the end of this section.
By lemma 7.7 and by the inequality (7.19) in theorem 7.4 it is sufficient to prove the estimate (16) for smooth initial data $\left(u_{0}, u_{1}\right) \in D(\mathcal{A})$ : the general case then follows by an obvious density argument.

In this case by theorem 7.5 the solution of (7)-(10) is sufficiently smooth to justify all computations that follow.

We begin by generalizing lemma 8.4 :

Lemma 9.6. - The function $E: \mathbb{R}_{+} \rightarrow \mathbb{R}_{+}$is non-increasing, locally absolutely continuous and

$$
\begin{equation*}
E^{\prime}=-\int_{\Gamma_{1}}(m \cdot \nu) u^{\prime} g\left(u^{\prime}\right) \mathrm{d} \Gamma \quad \text { a.e. in } \quad \mathbb{R}_{+} . \tag{17}
\end{equation*}
$$

Proof. - Given $0 \leq S<T<\infty$ arbitrarily, we have

$$
\begin{gathered}
0=\int_{S}^{T} \int_{\Omega} u^{\prime}\left(u^{\prime \prime}-\Delta u\right) \mathrm{dx} \mathrm{dt} \\
=\int_{S}^{T} \int_{\Omega} u^{\prime} u^{\prime \prime}+\nabla u \cdot \nabla u^{\prime} \mathrm{dx} \mathrm{dt}-\int_{S}^{T} \int_{\Gamma} u^{\prime} \partial_{\nu} u \mathrm{~d} \Gamma \mathrm{dt} \\
=\int_{S}^{T} \int_{\Omega} u^{\prime} u^{\prime \prime}+\nabla u \cdot \nabla u^{\prime} \mathrm{dx} \mathrm{dt}+\int_{S}^{T} \int_{\Gamma_{1}}(m \cdot \nu) u^{\prime} g\left(u^{\prime}\right) \mathrm{d} \Gamma \mathrm{dt}
\end{gathered}
$$

whence

$$
\begin{equation*}
E(S)-E(T)=\int_{S}^{T} \int_{\Gamma_{1}}(m \cdot \nu) u^{\prime} g\left(u^{\prime}\right) \mathrm{d} \Gamma \mathrm{dt} \tag{18}
\end{equation*}
$$

Since $m \cdot \nu \geq 0$ on $\Gamma_{1}$ and $x g(x) \geq 0, \quad \forall x \in \mathbb{R}$, the right-hand side of (18) is nonnegative; hence $E$ is non-increasing. Furthermore, (18) implies that E is locally absolutely continuous and that (17) is satisfied.

Next we generalize the identity of lemma 8.8 :
Lemma 9.7. - Putting for brevity

$$
\begin{equation*}
M u:=2 m \cdot \nabla u+(n-1) u \tag{19}
\end{equation*}
$$

for any fixed $0 \leq S<T<\infty$ we have

$$
\begin{gather*}
2 \int_{S}^{T} E^{(p+1) / 2} \mathrm{dt}-\int_{S}^{T} E^{(p-1) / 2} \int_{\Gamma_{0}}\left(\partial_{\nu} u\right)^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt} \\
=\left[E^{(p-1) / 2} \int_{\Omega} u^{\prime} M u \mathrm{dx}\right]_{T}^{S}+\frac{p-1}{2} \int_{S}^{T} E^{(p-3) / 2} E^{\prime} \int_{\Omega} u^{\prime} M u \mathrm{dx} \mathrm{dt}  \tag{20}\\
+\int_{S}^{T} E^{(p-1) / 2} \int_{\Gamma_{1}}\left(u^{\prime}\right)^{2}-|\nabla u|^{2}-g\left(u^{\prime}\right) M u \mathrm{~d} \Gamma_{m} \mathrm{dt}
\end{gather*}
$$

Proof. - We have

$$
\begin{gather*}
0=\int_{S}^{T} E^{(p-1) / 2} \int_{\Omega}(M u)\left(u^{\prime \prime}-\Delta u\right) \mathrm{dx} \mathrm{dt} \\
=\left[E^{(p-1) / 2} \int_{\Omega} u^{\prime} M u \mathrm{dx}\right]_{S}^{T}-\frac{p-1}{2} \int_{S}^{T} E^{(p-3) / 2} E^{\prime} \int_{\Omega} u^{\prime} M u \mathrm{dx} \mathrm{dt}  \tag{21}\\
-\int_{S}^{T} E^{(p-1) / 2} \int_{\Omega} u^{\prime} M u^{\prime}+(M u)(\Delta u) \mathrm{dx} \mathrm{dt}
\end{gather*}
$$

Integrating by parts and using the relation div $m=n$ we transform the interior integral in the last term as follows :

$$
\begin{gathered}
\int_{\Omega} u^{\prime} M u^{\prime}+(M u)(\Delta u) \mathrm{dx} \\
=\int_{\Omega} m \cdot \nabla\left(u^{\prime}\right)^{2}+(n-1)\left(u^{\prime}\right)^{2}-\nabla u \cdot \nabla(M u) \mathrm{dx}+\int_{\Gamma}(M u) \partial_{\nu} u \mathrm{~d} \Gamma \\
=\int_{\Omega} m \cdot \nabla\left(u^{\prime}\right)^{2}+(n-1)\left(u^{\prime}\right)^{2}-2|\nabla u|^{2}-m \cdot \nabla|\nabla u|^{2} \\
-(n-1)|\nabla u|^{2} \mathrm{dx}+\int_{\Gamma}(M u) \partial_{\nu} u \mathrm{~d} \Gamma \\
=-\int_{\Omega}\left(u^{\prime}\right)^{2}+|\nabla u|^{2} \mathrm{dx}+\int_{\Gamma}(m \cdot \nu)\left(\left(u^{\prime}\right)^{2}-|\nabla u|^{2}\right)+(M u) \partial_{\nu} u \mathrm{~d} \Gamma \\
=-\int_{\Omega}\left(u^{\prime}\right)^{2}+|\nabla u|^{2} \mathrm{dx}+\int_{\Gamma_{0}}-(m \cdot \nu)|\nabla u|^{2}+(2 m \cdot \nabla u) \partial_{\nu} u \mathrm{~d} \Gamma \\
+\int_{\Gamma_{1}}\left(u^{\prime}\right)^{2}-|\nabla u|^{2}-(M u) g\left(u^{\prime}\right) \mathrm{d} \Gamma_{m} .
\end{gathered}
$$

Since (8) implies that $\nabla u=\nu \partial_{\nu} u$ on $\Gamma_{0}$, hence we conclude that

$$
\begin{gathered}
\int_{\Omega} u^{\prime} M u^{\prime}+(M u)(\Delta u) \mathrm{dx}=-2 E+\int_{\Gamma_{0}}\left(\partial_{\nu} u\right)^{2} \mathrm{~d} \Gamma_{m} \\
+\int_{\Gamma_{1}}\left(u^{\prime}\right)^{2}-|\nabla u|^{2}-(M u) g\left(u^{\prime}\right) \mathrm{d} \Gamma_{m} .
\end{gathered}
$$

Substituting into (21) we obtain (20).
The following lemma is an immediate consequence of the definition of the energy and of the hypothesis $\Gamma_{0} \neq \emptyset$. Here and in the sequel $c$ will denote diverse positive constants only depending on $E(0)$.

Lemma 9.8. - We have

$$
\begin{equation*}
\left|\int_{\Omega} u^{\prime} M u d x\right| \leq c E \tag{22}
\end{equation*}
$$

We deduce from lemma 9.8 and from the non-increasingness of the energy that

$$
\left|E^{(p-1) / 2} \int_{\Omega} u^{\prime} M u \mathrm{dx}\right| \leq c E^{(p+1) / 2} \leq c E
$$

and

$$
\left|E^{(p-3) / 2} E^{\prime} \int_{\Omega} u^{\prime} M u \mathrm{dx}\right| \leq-c E^{(p-1) / 2} E^{\prime} \leq-c\left(E^{(p+1) / 2}\right)^{\prime}
$$

hence the first and second terms on the right-hand side of the identity (20) are majorized by $c E(S)$. Since the second integral on the left-hand side of (20) is $\leq$ by (12), we deduce from (20) that
$2 \int_{S}^{T} E^{(p+1) / 2} \mathrm{dt} \leq c E(S)+\int_{S}^{T} E^{(p-1) / 2} \int_{\Gamma_{1}}\left(u^{\prime}\right)^{2}-|\nabla u|^{2}-g\left(u^{\prime}\right) M u \mathrm{~d} \Gamma_{m} \mathrm{dt} ;$
Using the definition of $M u$ and the relation $|m| \leq R$, we obtain for any fixed $\varepsilon>0$ the inequality

$$
\begin{gathered}
2 \int_{S}^{T} E^{(p+1) / 2} \mathrm{dt} \leq c E(S) \\
+\int_{S}^{T} E^{(p-1) / 2} \int_{\Gamma_{1}}\left(u^{\prime}\right)^{2}+R^{2} g\left(u^{\prime}\right)^{2}+\varepsilon u^{2}+\frac{(n-1)^{2}}{4 \varepsilon} g\left(u^{\prime}\right)^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt}
\end{gathered}
$$

Choosing $\varepsilon=\varepsilon(\Omega)$ such that

$$
\varepsilon \int_{\Gamma_{1}} u^{2} \mathrm{~d} \Gamma_{m} \leq \frac{1}{2} \int_{\Omega}|\nabla u|^{2} \mathrm{dx}(\leq E)
$$

hence we conclude that

$$
\begin{equation*}
\int_{S}^{T} E^{(p+1) / 2} d t \leq c E(S)+c \int_{S}^{T} E^{(p-1) / 2} \int_{\Gamma_{1}}\left(u^{\prime}\right)^{2}+g\left(u^{\prime}\right)^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt} . \tag{23}
\end{equation*}
$$

We are going to majorize the last integral in (23). Set

$$
\begin{equation*}
\Gamma_{2}=\left\{x \in \Gamma_{1}:\left|u^{\prime}(x)\right| \leq 1\right\} \quad \text { and } \quad \Gamma_{3}=\left\{x \in \Gamma_{1}:\left|u^{\prime}(x)\right|>1\right\} . \tag{24}
\end{equation*}
$$

(Note that $\Gamma_{2}$ and $\Gamma_{3}$ depend on $t \in \mathbb{R}_{+}$.) Using (15) and (17) we obtain that

$$
\begin{gather*}
\int_{S}^{T} E^{(p-1) / 2} \int_{\Gamma_{3}}\left(u^{\prime}\right)^{2}+g\left(u^{\prime}\right)^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt} \\
\leq-c \int_{S}^{T} E^{(p-1) / 2} E^{\prime} \mathrm{dt} \leq c E^{(p+1) / 2}(S) \leq c E(S) . \tag{25}
\end{gather*}
$$

Furthermore, using (14) we also have

$$
\begin{gathered}
\int_{\Gamma_{2}}\left(u^{\prime}\right)^{2}+g\left(u^{\prime}\right)^{2} \mathrm{~d} \Gamma_{m} \leq c \int_{\Gamma_{2}}\left(u^{\prime} g\left(u^{\prime}\right)\right)^{2 /(p+1)} \mathrm{d} \Gamma_{m} \\
\leq c\left(\int_{\Gamma_{2}} u^{\prime} g\left(u^{\prime}\right) \mathrm{d} \Gamma_{m}\right)^{2 /(p+1)} \leq c\left(-E^{\prime}\right)^{2 /(p+1)}
\end{gathered}
$$

hence, using the Young inequality we obtain for every $\varepsilon>0$ the estimate

$$
\begin{gather*}
\int_{S}^{T} E^{(p-1) / 2} \int_{\Gamma_{2}}\left(u^{\prime}\right)^{2}+g\left(u^{\prime}\right)^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt} \\
\leq c \int_{S}^{T} E^{(p-1) / 2}\left(-E^{\prime}\right)^{2 /(p+1)} \mathrm{dt}  \tag{26}\\
\leq \int_{S}^{T} \varepsilon E^{(p+1) / 2}-c(\varepsilon) E^{\prime} \mathrm{dt} \leq \varepsilon \int_{S}^{T} E^{(p+1) / 2} \mathrm{dt}+c(\varepsilon) E(S) .
\end{gather*}
$$

Combining (23), (25) and (26) we find that

$$
\int_{S}^{T} E^{(p+1) / 2} \mathrm{dt} \leq c(\varepsilon) E(S)+\varepsilon c \int_{S}^{T} E^{(p+1) / 2} \mathrm{dt}
$$

Choosing $\varepsilon$ such that $\varepsilon c<1$, hence we conclude that

$$
\begin{equation*}
\int_{S}^{T} E^{(p+1) / 2} \mathrm{dt} \leq c E(S) \tag{27}
\end{equation*}
$$

Letting $T \rightarrow+\infty$ we obtain that

$$
\int_{S}^{\infty} E^{(p+1) / 2} \mathrm{dt} \leq c E(S), \quad \forall S \geq 0
$$

and we conclude by applying theorem 9.1 with $\alpha=(p-1) / 2$ and $T=c E(0)^{-\alpha}$.

Thus the proof of the theorem is completed in the case of globally Lipschitz continuous functions $g$. Turning to the general case, let us admit for a moment the following

Lemma 9.9. - Let $g: \mathbb{R} \rightarrow \mathbb{R}$ be a non-decreasing, continuous function satisfying the inequalities (14) and (15) for some $p \geq 1$ and for some positive constants $c_{i}$. Then there exists a sequence of non-decreasing, globally Lipschitz continuous functions $g_{k}: \mathbb{R} \rightarrow \mathbb{R}$ satisfying the inequalities

$$
\begin{equation*}
c_{1}^{\prime}|x|^{p} \leq\left|g_{k}(x)\right| \leq c_{2}^{\prime}|x|^{1 / p} \quad \text { if } \quad|x| \leq 1 \tag{28}
\end{equation*}
$$

and

$$
\begin{equation*}
c_{3}^{\prime}|x| \leq\left|g_{k}(x)\right| \leq c_{4}^{\prime}|x| \quad \text { if } \quad|x|>1 . \tag{29}
\end{equation*}
$$

with suitable positive constants $c_{i}^{\prime}$, independent of $k$, and such that

$$
\begin{equation*}
\left|g_{k}(x)\right| \leq|g(x)|, \quad \forall x \in \mathbb{R}, \quad k=1,2, \ldots \tag{30}
\end{equation*}
$$

and

$$
\begin{equation*}
g_{k}(x) \rightarrow g(x), \quad \forall x \in \mathbb{R} \tag{31}
\end{equation*}
$$

If we replace the function $g$ in the equation (9) by one of the functions $g_{k}$, the already proved part of the theorem applies. Using theorem 7.3 the general case of the theorem will follow if we prove that

$$
\begin{equation*}
\left(I+\mathcal{A}_{k}\right)^{-1} W \rightarrow(I+\mathcal{A})^{-1} W \quad \text { in } \quad \mathcal{H}, \quad \forall W \in \mathcal{H} \tag{32}
\end{equation*}
$$

where $\mathcal{A}_{k}$ denotes the generator of the semigroup associated with the problem (7)-(10) where $g$ is replaced by $g_{k}$. Putting $f=W_{2}-A W_{1}$ and denoting by $u, u_{k}$ the solutions of

$$
(I+A+B) u=f \quad \text { and } \quad\left(I+A+B_{k}\right) u_{k}=f
$$

where we use the same notations as in the proof of proposition 7.6 (and $B_{k}$ denotes the operator analogous to $B$ but corresponding to the function $g_{k}$ ), it is sufficient to prove that $u_{k} \rightarrow u$ in $H$ as $k \rightarrow+\infty$.

We have

$$
\left\langle\left(I+A+B_{k}\right) u_{k}-\left(I+A+B_{k}\right) u, u_{k}-u\right\rangle_{V^{\prime}, V}=\left\langle\left(B-B_{k}\right) u, u_{k}-u\right\rangle_{V^{\prime}, V}
$$

whence

$$
\left\|u_{k}-u\right\|_{V}^{2} \leq\left\|\left(B-B_{k}\right) u\right\|_{V^{\prime}}\left\|u_{k}-u\right\|_{V}
$$

and therefore

$$
\left\|u_{k}-u\right\|_{V} \leq\left\|g(u)-g_{k}(u)\right\|_{L^{(2 n-2) / n}\left(\Gamma_{1}\right)}
$$

By (28) and (29) the expression on the right-hand side tends to zero by Lebesgue's dominated convergence theorem. Hence $u_{k} \rightarrow u$ in $V$ and therefore also in $H$. Since

$$
\left(I+\mathcal{A}_{k}\right)^{-1} W=\left(u_{k}+W_{1}, u_{k}\right) \quad \text { and } \quad(I+\mathcal{A})^{-1} W=\left(u+W_{1}, u\right)
$$

hence (32) follows.
Proof of lemma 9.9. - Set

$$
g_{k}(x):=g\left(\left(\mathrm{id}_{\mathbb{R}}+k^{-1} g\right)^{-1}(x)\right), \quad x \in \mathbb{R}, k=1,2, \ldots
$$

One may readily verify that the functions $g_{k}$ are well-defined, continuous, non-decreasing and that properties (30), (31) are satisfied. It follows from
(14), (15) and (30) that the right-hand side inequalities in (28), (29) are satisfied with $c_{2}^{\prime}=c_{2}$ and $c_{4}^{\prime}=c_{4}$.

To prove their Lipschitz continuity choose $x_{1}, x_{2} \in \mathbb{R}$ arbitrarily and set

$$
y_{i}:=\left(\mathrm{id}_{\mathbb{R}}+k^{-1} g\right)^{-1}\left(x_{i}\right), \quad i=1,2 .
$$

Then we have $g_{k}\left(x_{i}\right)=g\left(y_{i}\right)$ and therefore

$$
x_{1}-x_{2}=\left(y_{1}-y_{2}\right)+k^{-1}\left(g_{k}\left(x_{1}\right)-g_{k}\left(x_{2}\right)\right) .
$$

Using the monotonicity of $g$ hence we conclude that

$$
\left|g_{k}\left(x_{1}\right)-g_{k}\left(x_{2}\right)\right| \leq k\left|x_{1}-x_{2}\right|
$$

To prove the first inequality in (28) fix $0 \leq x \leq 1$ arbitrarily (the case $-1 \leq x \leq 0$ is similar) and set $y:=\left(\operatorname{id}_{\mathbb{R}}+k^{-1} g\right)^{-1}(x)$. Then $y+k^{-1} g(y)=x$ and $0 \leq y \leq x$. If $y \geq x / 2$, then

$$
g_{k}(x)=g(y) \geq g(x / 2) \geq c_{1}(x / 2)^{p} .
$$

If $y<x / 2$, then $k^{-1} g(y) \geq x / 2$ and therefore

$$
g_{k}(x)=g(y) \geq x / 2 \geq(x / 2)^{p} .
$$

Hence the first inequality in (28) is satisfied with $c_{1}^{\prime}:=2^{-p} \min \left\{c_{1}, 1\right\}$.
To prove the first inequality in (29) let us choose a number $0<\varepsilon \leq 1 / 2$ such that $\varepsilon+g(\varepsilon)<1$ and $-\varepsilon+g(-\varepsilon)>-1$. Observe that (14) and (15) imply the existence of a positive constant $c_{3}^{\prime \prime}$ such that

$$
|g(y)| \geq c_{3}^{\prime \prime}|y| \quad \text { if } \quad|y| \geq \varepsilon
$$

Now fix $x \geq 1$ arbitrarily (the case $x \leq-1$ is analogous) and set $y:=\left(\operatorname{id}_{\mathbb{R}}+k^{-1} g\right)^{-1}(x)$ as above. We have $y+k^{-1} g(y)=x$ and $\varepsilon \leq y \leq x$. If $y \geq x / 2$, then

$$
g_{k}(x)=g(y) \geq g(x / 2) \geq\left(c_{3}^{\prime \prime} / 2\right) x
$$

because $x / 2 \geq 1 / 2 \geq \varepsilon$. If $y<x / 2$, then $k^{-1} g(y) \geq x / 2$ and therefore

$$
g_{k}(x)=g(y) \geq x / 2
$$

Hence the first inequality in (29) is satisfied with $c_{3}^{\prime}:=2^{-1} \min \left\{c_{3}^{\prime \prime}, 1\right\}$.

### 9.3. Uniform stabilization of the wave equation II

The theorem 9.4 proved in the preceding section has a serious drawback : it never applies for bounded functions $g$ (because of $c_{3}>0$ in (15)). The purpose of this section is to obtain a variant of this theorem for bounded feedback functions.

Theorem 9.10. - Assume (11)-(13) and assume that the function $g$ is bounded, globally Lipschitz continuous and that the inequalities (14) are satisfied with some positive constants $c_{1}, c_{2}$ and with a number $p$ satisfying

$$
\begin{equation*}
p \geq n-1 \tag{30}
\end{equation*}
$$

Then for every

$$
\begin{equation*}
\left(u^{0}, u^{1}\right) \in\left(H^{2}(\Omega) \cap H_{\Gamma_{0}}^{1}(\Omega)\right) \times H_{\Gamma_{0}}^{1}(\Omega) \tag{31}
\end{equation*}
$$

satisfying

$$
\begin{equation*}
\frac{\partial u^{0}}{\partial \nu}+(m \cdot \nu) g\left(u^{1}\right)=0 \quad \text { on } \quad \Gamma_{1} \tag{32}
\end{equation*}
$$

the solution of (7)-(10) satisfies the estimates

$$
\begin{equation*}
E(t) \leq C t^{2 /(1-p)}, \forall t>0 \tag{33}
\end{equation*}
$$

with a constant $C$ depending on the initial data.

Remark 9.11. - Observe that if the condition (30) is not satisfied initially, it will be satisfied if we replace $p$ by $n-1$; at the same time (14) continues to hold with the same constants $c_{1}, c_{2}$.

Proof. - Repeating the proof of the preceding theorem, except that part where the first inequality of (15) involving $c_{3}$ is applied, now we obtain the following inequality :

$$
\begin{equation*}
\int_{S}^{T} E^{(p+1) / 2} \mathrm{dt} \leq c E(S)+c \int_{S}^{T} E^{(p-1) / 2} \int_{\Gamma_{3}}\left(u^{\prime}\right)^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt} . \tag{34}
\end{equation*}
$$

It remains to establish the estimate

$$
\begin{equation*}
E^{(p-1) / 2} \int_{\Gamma_{3}}\left(u^{\prime}\right)^{2} \mathrm{~d} \Gamma_{m} \leq \varepsilon E^{(p+1) / 2}-c(\varepsilon) E^{\prime} \tag{35}
\end{equation*}
$$

for every $\varepsilon>0$. Then the theorem will follow. Indeed, choosing a sufficiently small $\varepsilon$, (34)-(35) imply (27) and the proof may be completed by applying theorem 9.1 as in the preceding section.

For brevity we shall denote the norm of $L^{\beta}(\Gamma)$ by $\left\|\|_{\beta}\right.$. Set

$$
s:=2 /(p+1) \quad \text { and } \quad \alpha:=(2-s) /(1-s) ;
$$

we have $0<s<1$ and $\alpha=2 p /(p-1)>2$. We establish for every $\varepsilon>0$ the inequality

$$
\begin{equation*}
E^{(p-1) / 2} \int_{\Gamma_{3}}\left(u^{\prime}\right)^{2} \mathrm{~d} \Gamma_{m} \leq \varepsilon E^{(p+1) / 2}\left\|u^{\prime}\right\|_{\alpha}^{\alpha}-c(\varepsilon) E^{\prime} . \tag{36}
\end{equation*}
$$

Indeed, we have

$$
\begin{gathered}
E^{(p-1) / 2} \int_{\Gamma_{3}}\left(u^{\prime}\right)^{2} \mathrm{~d} \Gamma_{m} \leq c E^{(p-1) / 2} \int_{\Gamma_{3}}\left|u^{\prime}\right|^{2-s}\left(u^{\prime} g\left(u^{\prime}\right)\right)^{s} \mathrm{~d} \Gamma_{m} \\
\leq c E^{(p-1) / 2}\left\|\left|u^{\prime}\right|^{2-s}\right\|_{1 /(1-s)}\left\|\left(u^{\prime} g\left(u^{\prime}\right)\right)^{s}\right\|_{1 / s} \\
=c E^{(p-1) / 2}\left\|u^{\prime}\right\|_{\alpha}^{(1-s) \alpha}\left\|u^{\prime} g\left(u^{\prime}\right)\right\|_{1}^{s} \\
=c E^{(p-1) / 2}\left\|u^{\prime}\right\|_{\alpha}^{(1-s) \alpha}\left(-E^{\prime}\right)^{s} \leq \varepsilon E^{(p-1) /(2(1-s))}\left\|u^{\prime}\right\|_{\alpha}^{\alpha}-c(\varepsilon) E^{\prime} \\
=\varepsilon E^{(p+1) / 2}\left\|u^{\prime}\right\|_{\alpha}^{\alpha}-c(\varepsilon) E^{\prime} .
\end{gathered}
$$

Using the trace theorem

$$
H^{1}(\Omega) \hookrightarrow L^{2 p /(p-1)}(\Gamma)=L^{\alpha}(\Gamma)
$$

(following from (30)) and the regularity property (7.24) we deduce from (36) that

$$
E^{(p-1) / 2} \int_{\Gamma_{3}}\left(u^{\prime}\right)^{2} \mathrm{~d} \Gamma_{m} \leq c \varepsilon E^{(p+1) / 2}-c(\varepsilon) E^{\prime} ;
$$

hence (35) follows (with another $\varepsilon$ ).

### 9.4. Uniform stabilization of Kirchhoff plates

Consider the problem of section 7.3 :

$$
\begin{gather*}
u^{\prime \prime}+\Delta^{2} u=0 \quad \text { in } \Omega \times \mathbb{R}_{+},  \tag{37}\\
u=u_{\nu}=0 \quad \text { on } \Gamma_{0} \times \mathbb{R}_{+},  \tag{38}\\
u_{\nu \nu}+\mu u_{\tau \tau}=0 \quad \text { on } \Gamma_{1} \times \mathbb{R}_{+},  \tag{39}\\
u_{\nu \nu \nu}+(2-\mu) u_{\tau \tau \nu}=(m \cdot \nu) g\left(u^{\prime}\right) \quad \text { on } \Gamma_{1} \times \mathbb{R}_{+},  \tag{40}\\
u(0)=u^{0} \quad \text { and } u^{\prime}(0)=u^{1} \quad \text { on } \quad \Omega . \tag{41}
\end{gather*}
$$

Assume that there is a point $X^{0} \in \mathbb{R}^{2}$ such that

$$
\begin{equation*}
m \cdot \nu \geq 0 \text { on } \Gamma_{1} \quad \text { and } \quad m \cdot \nu \leq 0 \text { on } \quad \Gamma_{0} \tag{42}
\end{equation*}
$$

We shall prove the
Theorem 9.12. - Assume (42) and assume that there exist two numbers $p, q>1$ and two positive constants $c_{1}, c_{2}$ such that

$$
\begin{equation*}
c_{1}|x|^{p} \leq|g(x)| \leq c_{2}|x|^{1 / p} \quad \text { if } \quad|x| \leq 1 \tag{43}
\end{equation*}
$$

and

$$
\begin{equation*}
c_{1}|x| \leq|g(x)| \leq c_{2}|x|^{q} \quad \text { if } \quad|x|>1 . \tag{44}
\end{equation*}
$$

Then for any given $\left(u^{0}, u^{1}\right) \in H_{\Gamma_{0}}^{2} \times L^{2}(\Omega)$ the solution of (37)-(40) satisfies the estimates

$$
\begin{equation*}
E(t) \leq C t^{2 /(1-p)}, \quad \forall t>0 \tag{45}
\end{equation*}
$$

with a constant $C$ only depending on the initial energy $E(0)$ (and in a continuous way). $\quad \square$

Remark 9.13. - A similar result holds for $p=1$ : then the estimate (45) is replaced by

$$
E(t) \leq C e^{-\omega t}, \quad \forall t>0
$$

with two positive constants $C, \omega$ which do not depend on the initial data.
Remark 9.14. - Theorem 9.12 improves some earlier results of Lagnese [2] (whose method is followed here) and of Rao [1] by weakening their growth assumptions and by using one feedback only. Also see Komornik [11] for a more general result.

There are many other theorems concerning the strong or uniform stabilization of different plate and beam models; see e.g. Bartolomeo and Triggiani [1], Horn [1], Horn and I. Lasiecka [1], Lagnese [2], Lagnese and Leugering [1], I. Lasiecka [1], [2], Lasiecka and Tataru [1], Lasiecka and Triggiani [3], Lebeau [1], W. Littmann and L. Markus [1], Littmann and Taylor [1], Puel and Tucsnak [1], Tataru [1].

The rest of this section is devoted to the proof of theorem 9.12. Using lemma 9.9 and the argument following it in section 9.2 we may assume without loss of generality that $g$ is globally Lipschitz. Furthermore, using theorem 7.9 and a density argument we may assume that the solutions verify
the regularity conditions (7.53) and (7.54) ; then all computations that follow are justified.

Lemma 9.15. - The function $E: \mathbb{R}_{+} \rightarrow \mathbb{R}_{+}$is non-increasing, locally absolutely continuous and

$$
\begin{equation*}
E^{\prime}=-\int_{\Gamma_{1}} u^{\prime} g\left(u^{\prime}\right) \mathrm{d} \Gamma_{m} \quad \text { a.e. on } \quad \mathbb{R}_{+} \tag{46}
\end{equation*}
$$

Proof. - Given $0 \leq S<T<+\infty$ arbitrarily, we apply lemma 7.12 ; we obtain that

$$
\begin{aligned}
0= & \int_{S}^{T} \int_{\Omega} u^{\prime}\left(u^{\prime \prime}+\Delta^{2} u\right) \mathrm{dX} \mathrm{dt}=\int_{S}^{T} \int_{\Omega} u^{\prime} u^{\prime \prime}+Q\left(u, u^{\prime}\right) \mathrm{dX} \mathrm{dt} \\
& +\int_{S}^{T} \int_{\Gamma}\left(u_{\nu \nu}+(2-\mu) u_{\tau \tau}\right)_{\nu} u^{\prime}-\left(u_{\nu \nu}+\mu u_{\tau \tau}\right) u_{\nu}^{\prime} \mathrm{d} \Gamma \mathrm{dt} \\
= & {\left[\frac{1}{2} \int_{\Omega}\left(u^{\prime}\right)^{2}+Q(u) \mathrm{dX}\right]_{S}^{T}+\int_{S}^{T} \int_{\Gamma_{1}}(m \cdot \nu) u^{\prime} g\left(u^{\prime}\right) \mathrm{d} \Gamma \mathrm{dt} } \\
= & {\left[\frac{1}{2} \int_{\Omega}\left(u^{\prime}\right)^{2}+Q(u) \mathrm{dX}\right]_{S}^{T}+\int_{S}^{T} \int_{\Gamma_{1}}(m \cdot \nu) u^{\prime} g\left(u^{\prime}\right) \mathrm{d} \Gamma \mathrm{dt}, }
\end{aligned}
$$

i.e.

$$
\begin{equation*}
E(S)-E(T)=\int_{S}^{T} \int_{\Gamma_{1}} u^{\prime} g\left(u^{\prime}\right) \mathrm{d} \Gamma_{m} \mathrm{dt} \tag{47}
\end{equation*}
$$

This implies that $E$ is locally absolutely continuous; differentiating (47) we find (46). Finally, $E$ is non-increasing because the right-hand side of (47) is nonnegative by the increasingness of $g$ and by hypothesis (42).
Lemma 9.16. - Given $v \in H^{3}(\Omega)$ arbitrarily, the following identity holds true:

$$
\begin{equation*}
\int_{\Omega} Q(v, m \cdot \nabla v) \mathrm{dX}=\int_{\Omega} Q(v) d X+\frac{1}{2} \int_{\Gamma} Q(v) \mathrm{d} \Gamma_{m} \tag{48}
\end{equation*}
$$

Proof. - Integrating by parts we obtain the following five simple identities :

$$
\begin{gathered}
\int_{\Omega} v_{x x}\left(m_{1} v_{x}+m_{2} v_{y}\right)_{x x} \mathrm{dX}=\int_{\Omega} 2 v_{x x}^{2}+m_{1} v_{x x} v_{x x x}+m_{2} v_{x x} v_{x x y} \mathrm{dX} \\
=\int_{\Omega} 2 v_{x x}^{2}+\frac{1}{2} m \cdot \nabla\left(v_{x x}^{2}\right) \mathrm{dX} \\
\begin{array}{c}
\int_{\Omega} v_{y y}\left(m_{1} v_{x}+m_{2} v_{y}\right)_{y y} \mathrm{dX}=\int_{\Omega} 2 v_{y y}^{2}+m_{1} v_{y y} v_{y y x}+m_{2} v_{y y} v_{y y y} \mathrm{dX} \\
=\int_{\Omega} 2 v_{y y}^{2}+\frac{1}{2} m \cdot \nabla\left(v_{y y}^{2}\right) \mathrm{dX}
\end{array} .
\end{gathered}
$$

$$
\begin{gathered}
\int_{\Omega} v_{x x}\left(m_{1} v_{x}+m_{2} v_{y}\right)_{y y} \mathrm{dX}=\int_{\Omega} 2 v_{x x} v_{y y}+m_{1} v_{x x} v_{x y y}+m_{2} v_{x x} v_{y y y} \mathrm{dX} \\
\int_{\Omega} v_{y y}\left(m_{1} v_{x}+m_{2} v_{y}\right)_{x x} \mathrm{dX}=\int_{\Omega} 2 v_{x x} v_{y y}+m_{1} v_{y y} v_{x x x}+m_{2} v_{y y} v_{y x x} \mathrm{dX} \\
\int_{\Omega} v_{x y}\left(m_{1} v_{x}+m_{2} v_{y}\right)_{x y} \mathrm{dX}=\int_{\Omega} 2 v_{x y}^{2}+m_{1} v_{x y} v_{x x y}+m_{2} v_{x y} v_{x y y} \mathrm{dX} \\
=\int_{\Omega} 2 v_{x y}^{2}+\frac{1}{2} m \cdot \nabla\left(v_{x y}^{2}\right) \mathrm{dX} .
\end{gathered}
$$

Combining these identities, using the definitions of $Q(v), Q(u, v)$ and finally applying the divergence theorem we obtain that

$$
\begin{gathered}
\int_{\Omega} Q(v, m \cdot \nabla v) \mathrm{dX}=\int_{\Omega} 2 Q(v)+\frac{1}{2} m \cdot \nabla Q(v) \mathrm{dX}=\int_{\Omega}\left(2-\frac{1}{2} \operatorname{div} m\right) Q(v) \mathrm{dX} \\
\quad+\frac{1}{2} \int_{\Gamma}(m \cdot \nu) Q(v) \mathrm{d} \Gamma=\int_{\Omega} Q(v) \mathrm{dX}+\frac{1}{2} \int_{\Gamma} Q(v) \mathrm{d} \Gamma_{m}
\end{gathered}
$$

Now we prove our basic identity :
Lemma 9.17. - Given $0 \leq S<T<+\infty$ arbitrarily, the following identity holds true :

$$
\begin{gather*}
4 \int_{S}^{T} E^{(p+1) / 2} \mathrm{dt}-\int_{S}^{T} E^{(p-1) / 2} \int_{\Gamma_{0}}(\Delta u)^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt} \\
=\left[E^{(p-1) / 2} \int_{\Omega} 2 u^{\prime} m \cdot \nabla u \mathrm{dX}\right]_{T}^{S} \\
\quad+(p-1) \int_{S}^{T} E^{(p-3) / 2} E^{\prime} \int_{\Omega} u^{\prime} m \cdot \nabla u \mathrm{dX} \mathrm{dt}  \tag{49}\\
+\int_{S}^{T} E^{(p-1) / 2} \int_{\Gamma_{1}}\left(u^{\prime}\right)^{2}-Q(u)-2(m \cdot \nabla u) g\left(u^{\prime}\right) \mathrm{d} \Gamma_{m} \mathrm{dt} .
\end{gather*}
$$

Proof. - We multiply the equation (39) by $2 E^{(p-1) / 2} m \cdot \nabla u$ and we integrate by parts in $\Omega \times(S, T)$. Using lemmas 7.12 and 9.16 we obtain that

$$
\begin{gathered}
0=\int_{S}^{T} E^{(p-1) / 2} \int_{\Omega} 2 m \cdot \nabla u\left(u^{\prime \prime}+\Delta^{2} u\right) \mathrm{dX} \mathrm{dt} \\
=\left[E^{(p-1) / 2} \int_{\Omega} 2 u^{\prime} m \cdot \nabla u \mathrm{dX}\right]_{S}^{T}-(p-1) \int_{S}^{T} E^{(p-3) / 2} E^{\prime} \int_{\Omega} u^{\prime} m \cdot \nabla u d X d t
\end{gathered}
$$

$$
\begin{gathered}
+\int_{S}^{T} E^{(p-1) / 2} \int_{\Omega}-2 u^{\prime} m \cdot \nabla u^{\prime}+2 Q(u, m \cdot \nabla u) \mathrm{dX} \mathrm{dt} \\
+2 \int_{S}^{T} E^{(p-1) / 2} \int_{\Gamma}\left(u_{\nu \nu}+(2-\mu) u_{\tau \tau}\right)_{\nu} m \cdot \nabla u-\left(u_{\nu \nu}+\mu u_{\tau \tau}\right)(m \cdot \nabla u)_{\nu} \mathrm{d} \Gamma \mathrm{dt} \\
=\left[E^{(p-1) / 2} \int_{\Omega} 2 u^{\prime} m \cdot \nabla u \mathrm{dX}\right]_{S}^{T}-(p-1) \int_{S}^{T} E^{(p-3) / 2} E^{\prime} \int_{\Omega} u^{\prime} m \cdot \nabla u \mathrm{dX} \mathrm{dt} \\
\quad+\int_{S}^{T} E^{(p-1) / 2} \int_{\Omega}-m \cdot \nabla\left(u^{\prime}\right)^{2}+2 Q(u) \mathrm{dX} \mathrm{dt} \\
+\int_{S}^{T} E^{(p-1) / 2} \int_{\Gamma}(m \cdot \nu) Q(u)+2\left(u_{\nu \nu}+(2-\mu) u_{\tau \tau}\right)_{\nu} m \cdot \nabla u \\
=\left[E^{(p-1) / 2} \int_{\Omega} 2 u^{\prime} m \cdot \nabla u \mathrm{dX}\right]_{S}^{T}-(p-1) \int_{S}^{T} E^{(p-3) / 2} E^{\prime} \int_{\Omega} u^{\prime} m \cdot \nabla u \mathrm{dX} \mathrm{dt} \\
\quad-2\left(u_{\nu \nu}+\mu u_{\tau \tau}\right)(m \cdot \nabla u)_{\nu} \mathrm{d} \Gamma \mathrm{dt} \\
\quad+2 \int_{S}^{T} E^{(p-1) / 2} \int_{\Omega}\left(u^{\prime}\right)^{2}+Q(u) \mathrm{dX} \mathrm{dt} \\
\quad+\int_{S}^{T} E^{(p-1) / 2} \int_{\Gamma}(m \cdot \nu)\left(Q(u)-\left(u^{\prime}\right)^{2}\right)+2\left(u_{\nu \nu}+(2-\mu) u_{\tau \tau}\right)_{\nu} m \cdot \nabla u \\
-2\left(u_{\nu \nu}+\mu u_{\tau \tau}\right)(m \cdot \nabla u)_{\nu} \mathrm{d} \Gamma \mathrm{dt} .
\end{gathered}
$$

Recalling the definition of the energy we conclude that

$$
\begin{gathered}
4 \int_{S}^{T} E^{(p+1) / 2} \mathrm{dt} \\
=\left[E^{(p-1) / 2} \int_{\Omega} 2 u^{\prime} m \cdot \nabla u \mathrm{dX}\right]_{T}^{S}+(p-1) \int_{S}^{T} E^{(p-3) / 2} E^{\prime} \int_{\Omega} u^{\prime} m \cdot \nabla u \mathrm{dX} \mathrm{dt} \\
\quad+\int_{S}^{T} E^{(p-1) / 2} \int_{\Gamma}(m \cdot \nu)\left(\left(u^{\prime}\right)^{2}-Q(u)\right)-2\left(u_{\nu \nu}+(2-\mu) u_{\tau \tau}\right)_{\nu} m \cdot \nabla u \\
+2\left(u_{\nu \nu}+\mu u_{\tau \tau}\right)(m \cdot \nabla u)_{\nu} \mathrm{d} \Gamma \mathrm{dt} .
\end{gathered}
$$

It remains to show that

$$
\begin{gathered}
(m \cdot \nu)\left(\left(u^{\prime}\right)^{2}-Q(u)\right)-2\left(u_{\nu \nu}+(2-\mu) u_{\tau \tau}\right)_{\nu} m \cdot \nabla u \\
+2\left(u_{\nu \nu}+\mu u_{\tau \tau}\right)(m \cdot \nabla u)_{\nu}=(m \cdot \nu)\left\{\left(u^{\prime}\right)^{2}-Q(u)-2(m \cdot \nabla u)\left(\alpha u+g\left(u^{\prime}\right)\right)\right\}
\end{gathered}
$$

on $\Gamma_{1}$ and

$$
\begin{gather*}
(m \cdot \nu)\left(\left(u^{\prime}\right)^{2}-Q(u)\right)-2\left(u_{\nu \nu}+(2-\mu) u_{\tau \tau}\right)_{\nu} m \cdot \nabla u \\
+2\left(u_{\nu \nu}+\mu u_{\tau \tau}\right)(m \cdot \nabla u)_{\nu}=(\Delta u)^{2} \tag{50}
\end{gather*}
$$

on $\Gamma_{0}$.

The first equality follows at once from the boundary conditions (39) and (40).

For the proof of (50) first observe that the boundary conditions (38) imply $u^{\prime}=u_{x}=u_{y}=0$ on $\Gamma_{0}$. Consequently, the vectors $\nabla u_{x}=\left(u_{x x}, u_{x y}\right)$ and $\nabla u_{y}=\left(u_{x y}, u_{y y}\right)$ are orthogonal to $\Gamma_{0}$ and therefore they are parallel. Hence $u_{x x} u_{y y}=u_{x y}^{2}$ and $Q(u)=(\Delta u)^{2}$ on $\Gamma_{0}$. Furthermore, $m \cdot \nabla u=(m \cdot \nu) u_{\nu}$ and

$$
(m \cdot \nabla u)_{\nu}=\left(m_{1} u_{x}+m_{2} u_{y}\right)_{\nu}=(m \cdot \nu) u_{\nu \nu}+(m \cdot \tau) u_{\nu \tau}=(m \cdot \nu) u_{\nu \nu}
$$

on $\Gamma_{0}$ because $u_{\nu}=0$ on $\Gamma_{0}$.
Combining these relations we conclude that the left-hand side of (50) is equal to

$$
-(m \cdot \nu)(\Delta u)^{2}+2(m \cdot \nu) u_{\nu \nu}^{2}
$$

Since the boundary conditions (38) imply that $u_{\tau \tau}=0$ on $\Gamma_{0}$, we have

$$
-(m \cdot \nu)(\Delta u)^{2}+2(m \cdot \nu) u_{\nu \nu}^{2}=(m \cdot \nu)(\Delta u)^{2} .
$$

This completes the proof of (50) and that of the lemma.
In order to majorize the right-hand side of (49) we need the
Lemma 9.18. - The semi-norm $p: H^{2}(\Omega) \rightarrow \mathbb{R}_{+}$defined by

$$
\begin{equation*}
p(v):=\left(\int_{\Omega} v_{x x}^{2}+v_{x y}^{2}+v_{y y}^{2} \mathrm{dX}\right)^{1 / 2} \tag{51}
\end{equation*}
$$

is equivalent to the norm $\left\|\|_{H^{2}(\Omega)}\right.$ on the subspace $H_{\Gamma_{0}}^{2}(\Omega)$.
Proof. - First we show that $p$ is a norm on $H_{\Gamma_{0}}^{2}(\Omega)$. Indeed, let $v \in H_{\Gamma_{0}}^{2}(\Omega)$ such that $p(v)=0$. Then the second derivative of $v$ vanishes. Since $\Omega$ is connexe, hence we conclude that $v$ is affine. Since $\Gamma_{0} \neq \emptyset$ and $v=v_{x}=v_{y}=0$ on $\Gamma_{0}$ by hypotheses $\Gamma_{0} \neq \emptyset$ and (40), it follows that $v=0$.

The estimate

$$
\begin{equation*}
p \leq c\| \|_{H^{2}(\Omega)} \tag{52}
\end{equation*}
$$

is obvious. Assume that the inverse inequality is false. Then there exists a sequence $v^{k}$ in $H_{\Gamma_{0}}^{2}(\Omega)$ such that

$$
\begin{equation*}
\left\|v^{k}\right\|_{H^{2}(\Omega)}=1, \quad k=1,2, \ldots \tag{53}
\end{equation*}
$$

and

$$
\begin{equation*}
p\left(v^{k}\right) \rightarrow 0 \tag{54}
\end{equation*}
$$

By the compactness of the imbedding $H^{2}(\Omega) \subset H^{1}(\Omega)$ we may assume, by extracting a subsequence if needed, that

$$
v^{k} \rightarrow v \quad \text { in } \quad H^{1}(\Omega)
$$

for some $v \in H^{1}(\Omega)$. Using (51) and (54) we conclude that

$$
v^{k} \rightarrow v \quad \text { in } \quad H^{2}(\Omega) ;
$$

consequently, using (52) and (54) we obtain that $v=0$ and hence

$$
v^{k} \rightarrow 0 \quad \text { in } \quad H^{2}(\Omega)
$$

But this contradicts (53).
In what follows $c$ will denote diverse positive constants depending on $E(0)$ only.

Corollary 9.19. - We have

$$
\begin{equation*}
\left|\int_{\Omega} u^{\prime} m \cdot \nabla u \mathrm{dX}\right| \leq c E . \tag{55}
\end{equation*}
$$

Proof. - Using the trivial inequality

$$
\left|\int_{\Omega} u^{\prime} m \cdot \nabla u \mathrm{dX}\right| \leq c\left(\left\|u^{\prime}\right\|_{L^{2}(\Omega)}^{2}+\|u\|_{H^{1}(\Omega)}^{2}\right)
$$

and the definition of the energy, (55) follows from the preceding lemma and from the hypothesis $0<\mu<1$.

Lemma 9.20. - Given $0 \leq S<T<+\infty$ arbitrarily, we have

$$
\begin{gather*}
\int_{S}^{T} E^{(p+1) / 2} \mathrm{dt} \leq c E^{(p+1) / 2}(S)  \tag{56}\\
+c \int_{S}^{T} E^{(p-1) / 2} \int_{\Gamma_{1}}\left(u^{\prime}\right)^{2}+|\nabla u|\left|g\left(u^{\prime}\right)\right| \mathrm{d} \Gamma_{m} \mathrm{dt}
\end{gather*}
$$

Proof. - Since $m \cdot \nu \leq 0$ on $\Gamma_{0}$, the left-hand side of (49) is minorized by its first term. Applying corollary 9.19 and using the non-increasingness of $E$ we may estimate the first and second term on the right-hand side of (49) as follows :

$$
\begin{gathered}
{\left[E^{(p-1) / 2} \int_{\Omega} 2 u^{\prime} m \cdot \nabla u \mathrm{dX}\right]_{T}^{S}+(p-1) \int_{S}^{T} E^{(p-3) / 2} E^{\prime} \int_{\Omega} u^{\prime} m \cdot \nabla u \mathrm{dX} \mathrm{dt}} \\
\leq c E^{(p+1) / 2}(S)+c E^{(p+1) / 2}(T)-c \int_{S}^{T} E^{(p-1) / 2} E^{\prime} \mathrm{dt} \leq c E^{(p+1) / 2}(S)
\end{gathered}
$$

Furthermore, using hypothesis $0<\mu<1$ we have

$$
Q(u)=\mu\left(u_{x x}+u_{y y}\right)^{2}+(1-\mu)\left(u_{x x}^{2}+u_{y y}^{2}+2 u_{x y}^{2}\right) \geq 0
$$

and therefore we deduce from (49) that

$$
\begin{gathered}
4 \int_{S}^{T} E^{(p+1) / 2} \mathrm{dt} \leq c E^{(p+1) / 2}(S) \\
+\int_{S}^{T} E^{(p-1) / 2} \int_{\Gamma_{1}}\left(u^{\prime}\right)^{2}-2(m \cdot \nabla u) g\left(u^{\prime}\right) \mathrm{d} \Gamma_{m} \mathrm{dt}
\end{gathered}
$$

which implies (56).
Now we introduce the notation

$$
\Gamma_{2}=\left\{x \in \Gamma_{1}:\left|u^{\prime}(x)\right| \leq 1\right\} \quad \text { and } \quad \Gamma_{3}=\left\{x \in \Gamma_{1}:\left|u^{\prime}(x)\right|>1\right\} .
$$

Lemma 9.21. - Given $0 \leq S<T<+\infty$ arbitrarily, we have

$$
\begin{gather*}
\int_{S}^{T} E^{(p+1) / 2} \mathrm{dt} \leq c E^{(p+1) / 2}(S) \\
+c \int_{S}^{T} E^{(p-1) / 2} \int_{\Gamma_{2}}\left(u^{\prime}\right)^{2}+\left|g\left(u^{\prime}\right)\right|^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt}  \tag{57}\\
+c \int_{S}^{T} E^{(p-1) / 2} \int_{\Gamma_{3}}\left(u^{\prime}\right)^{2}+\left|g\left(u^{\prime}\right)\right|^{1+q^{-1}} \mathrm{~d} \Gamma_{m} \mathrm{dt} .
\end{gather*}
$$

Proof. - For any fixed $\varepsilon>0$ we deduce from (56) that

$$
\begin{gather*}
\int_{S}^{T} E^{(p+1) / 2} \mathrm{dt} \leq c E^{(p+1) / 2}(S) \\
+c \int_{S}^{T} E^{(p-1) / 2} \int_{\Gamma_{1}}\left(u^{\prime}\right)^{2}+\varepsilon|\nabla u|^{2}+c(\varepsilon)\left|g\left(u^{\prime}\right)\right|^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt}  \tag{58}\\
+c \int_{S}^{T} E^{(p-1) / 2} \int_{\Gamma_{2}}\left(u^{\prime}\right)^{2}+\varepsilon|\nabla u|^{2}+\varepsilon|\nabla u|^{1+q}+c(\varepsilon)\left|g\left(u^{\prime}\right)\right|^{1+q^{-1}} \mathrm{~d} \Gamma_{m} \mathrm{dt} .
\end{gather*}
$$

Using the Sobolev inequality

$$
\|u\|_{W^{1, \gamma}(\Gamma)} \leq c E^{1 / 2}, \quad \forall \gamma \geq 1
$$

with $\gamma=2$ and $\gamma=1+q$, we conclude from (58) that

$$
\begin{aligned}
& \int_{S}^{T} E^{(p+1) / 2} \mathrm{dt} \leq c E^{(p+1) / 2}(S)+c \varepsilon \int_{S}^{T} E^{(p+1) / 2} \mathrm{dt} \\
& \quad+c \int_{S}^{T} E^{(p-1) / 2} \int_{\Gamma_{2}}\left(u^{\prime}\right)^{2}+c(\varepsilon)\left|g\left(u^{\prime}\right)\right|^{2} \mathrm{~d} \Gamma_{m} \mathrm{dt} \\
& +c \int_{S}^{T} E^{(p-1) / 2} \int_{\Gamma_{3}}\left(u^{\prime}\right)^{2}+c(\varepsilon)\left|g\left(u^{\prime}\right)\right|^{1+q^{-1}} \mathrm{~d} \Gamma_{m} \mathrm{dt}
\end{aligned}
$$

The lemma follows if we choose $\varepsilon>0$ such that $c \varepsilon<1$.
Now we are ready to complete the proof of the theorem. Applying (43) and (44) we have

$$
\begin{aligned}
& \int_{\Gamma_{2}}\left|g\left(u^{\prime}\right)\right|^{2}+\left(u^{\prime}\right)^{2} \mathrm{~d} \Gamma \leq c \int_{\Gamma_{2}}\left(u^{\prime} g\left(u^{\prime}\right)\right)^{2 /(p+1)} \mathrm{d} \Gamma \\
& \quad \leq c\left(\int_{\Gamma_{2}} u^{\prime} g\left(u^{\prime}\right) \mathrm{d} \Gamma\right)^{2 /(p+1)} \leq c\left(-E^{\prime}\right)^{2 /(p+1)}
\end{aligned}
$$

and

$$
\int_{\Gamma_{3}}\left|g\left(u^{\prime}\right)\right|^{1+q^{-1}}+\left(u^{\prime}\right)^{2} \mathrm{~d} \Gamma \leq c \int_{\Gamma_{3}} u^{\prime} g\left(u^{\prime}\right) \mathrm{d} \Gamma \leq-c E^{\prime}
$$

Substituting into (57) we obtain for every $\varepsilon>0$ the estimate

$$
\begin{gathered}
\int_{S}^{T} E^{(p+1) / 2} \mathrm{dt} \leq c E^{(p+1) / 2}(S) \\
+c \int_{S}^{T} E^{(p-1) / 2}\left(-E^{\prime}\right)^{2 /(p+1)} \mathrm{dt}-c \int_{S}^{T} E^{(p-1) / 2} E^{\prime} \mathrm{dt} \\
\leq c E^{(p+1) / 2}(S)+\int_{S}^{T} \varepsilon E^{(p+1) / 2}-c(\varepsilon) E^{\prime} \mathrm{dt} \\
\leq c(\varepsilon) E(S)+\varepsilon \int_{S}^{T} E^{(p+1) / 2} \mathrm{dt} .
\end{gathered}
$$

Choosing $\varepsilon<1$ we conclude that $E$ satisfies an inequality of the form

$$
\int_{S}^{T} E^{(p+1) / 2} \mathrm{dt} \leq A E(S) \quad \text { for every } \quad 0 \leq S<T<+\infty
$$

with a suitable positive constant $A$. We may conclude by applying theorem 9.1.

## 10. Stabilization of the Korteweg-de Vries equation

In this chapter we shall study, following Komornik, Russell and Zhang [1], [2] the stabilizability of the non-linear Korteweg-de Vries equation by linear distributed feedbacks. The proof is based on a remarkable property of this equation : the existence of an infinite sequence of conservation laws corresponding to an infinite sequence of useful multipliers.

### 10.1. Formulation of the results

Let $\Omega=(0,1), k>0$ and consider the problem

$$
\begin{gather*}
u^{\prime}+u u_{x}+u_{x x x}=-k(u-[u]) \quad \text { in } \quad \Omega \times \mathbb{R}_{+},  \tag{1}\\
u(0, t)=u(1, t), \quad \forall t \in \mathbb{R}_{+},  \tag{2}\\
u_{x}(0, t)=u_{x}(1, t), \quad \forall t \in \mathbb{R}_{+},  \tag{3}\\
u_{x x}(0, t)=u_{x x}(1, t), \quad \forall t \in \mathbb{R}_{+},  \tag{4}\\
u(0)=u^{0} \quad \text { on } \quad \Omega \tag{5}
\end{gather*}
$$

where $[u]$ denotes the mean-value of $u$ defined by

$$
\begin{equation*}
[u]:=\int_{\Omega} u \mathrm{dx} . \tag{6}
\end{equation*}
$$

For $k=0$ the equation (1) is a good model of shallow water : $u(x, t)$ denotes the depth of water at a point $x$ at time $t$; see Miura [1], Temam [1]. The periodic boundary conditions correspond to a circular movement. In this model $[u]$ denotes the total volume of water.

For $k>0$ the action of the "feedback" $-k(u-[u])$ consists in balancing the level of water, conserving at the same time its total volume. Indeed, the latter property follows, at least formally, from (1), (2) and (5) :

$$
\begin{aligned}
{[u]^{\prime} } & =\int_{\Omega} u^{\prime} \mathrm{dx}=-\int_{\Omega} u u_{x}+u_{x x x}+k(u-[u]) \mathrm{dx} \\
& =-\int_{\Omega}\left(u^{2} / 2+u_{x x}\right)_{x} \mathrm{dx}+k[u]-k \int_{\Omega} u \mathrm{dx} \\
& =-\int_{\Omega}\left(u^{2} / 2+u_{x x}\right)_{x} \mathrm{dx}=\left[u^{2} / 2+u_{x x}\right]_{0}^{1}=0
\end{aligned}
$$

whence

$$
\begin{equation*}
[u(t)]=\left[u^{0}\right], \quad \forall t \in \mathbb{R}_{+} . \tag{7}
\end{equation*}
$$

The following formal computation shows that $u(t)$ converges exponentially to the constant $M:=\left[u^{0}\right]=[u]$ in $L^{2}(\Omega)$ as $t \rightarrow+\infty$ :

$$
\begin{gathered}
\left(\int(u-[u])^{2} \mathrm{dx}\right)^{\prime}=\int 2(u-M) u^{\prime} \mathrm{dx} \\
=\int-2(u-M)\left(u u_{x}+u_{x x x}+k(u-M)\right) \mathrm{dx} \\
=\int-2 u^{2} u_{x}+2 M u u_{x}-2 u u_{x x x}-2 k(u-M)^{2} \mathrm{dx} \\
=\left[-(2 / 3) u^{3}+M u^{2}-2 u u_{x x}+u_{x}^{2}\right]_{0}^{1}-2 k \int(u-M)^{2} \mathrm{dx} \\
=-2 k \int(u-M)^{2} \mathrm{dx}
\end{gathered}
$$

whence

$$
\begin{equation*}
\left\|u(t)-\left[u^{0}\right]\right\|_{L^{2}(\Omega)}=\left\|u^{0}-\left[u^{0}\right]\right\|_{L^{2}(\Omega)} e^{-k t}, \quad \forall t \geq 0 \tag{8}
\end{equation*}
$$

Let us introduce the Hilbert spaces

$$
\begin{gathered}
H_{p}^{m}:=\left\{w \in H^{m}(\Omega): w^{j}(0)=w^{j}(1), j=0, \ldots, m-1\right\}, m=1,2, \ldots \\
H_{p}^{0}:=L^{2}(\Omega) \text { and } H_{p}^{-1}:=\left(H_{p}^{1}\right)^{\prime} .
\end{gathered}
$$

Identifying $\left(L^{2}(\Omega)\right)^{\prime}$ with $L^{2}(\Omega)$ we obtain the algebraical and topological inclusions

$$
\cdots \subset H_{p}^{2} \subset H_{p}^{1} \subset H_{p}^{0} \subset H_{p}^{-1}
$$

The problem (1)-(5) is well-posed in the following sense :
Theorem 10.1. - Let $m \geq 2$ and $u^{0} \in H_{p}^{m}$. Then the problem (1)-(5) has a unique solution

$$
\begin{equation*}
u \in C\left(\mathbb{R}_{+} ; H_{p}^{m}\right) \cap C^{1}\left(\mathbb{R}_{+} ; H_{p}^{m-3}\right) \tag{9}
\end{equation*}
$$

Furthermore, the mapping $u^{0} \mapsto\left(u, u^{\prime}\right)$ is continuous from $H_{p}^{m}$ into $H^{m} \times H^{m-3}$ 。

The solution $u(t)$ converges quickly to the constant $\left[u^{0}\right]$ as $t \rightarrow+\infty$ :

Theorem 10.2. - Let $m \geq 2$ and $u^{0} \in H_{p}^{m}$. Then for every fixed $\left.k^{\prime} \in\right] 0, k[$ there exists a constant $C=C\left(u^{0}, k^{\prime}\right)$ such that the solution of (1)-(5) satisfies the estimate

$$
\begin{equation*}
\left\|\left(u(t)-\left[u^{0}\right], u^{\prime}(t)\right)\right\|_{H_{p}^{m} \times H_{p}^{m-3}} \leq C e^{-k^{\prime} t}, \quad \forall t \in \mathbb{R}_{+} \tag{10}
\end{equation*}
$$

Theorems 10.1 et 10.2 were proved in Komornik, Russell and Zhang [2]. Here we admit theorem 10.1 without proof and we only prove theorem 10.2 in the particular case $m=2$.

### 10.2. Uniform stabilization by linear feedbacks

We shall often use the equality (7) and therefore we shall write $[u]$ instead of $\left[u^{0}\right]$. For brevity we shall write $\int$ instead of $\int_{\Omega}$.

Applying a usual density argument it is sufficient to prove the estimates (10) for $u^{0} \in H_{p}^{5}$. According to theorem 10.1 thus we may assume that

$$
\begin{equation*}
u \in C\left(\mathbb{R}_{+} ; H_{p}^{5}\right) \cap C^{1}\left(\mathbb{R}_{+} ; H_{p}^{2}\right) \tag{11}
\end{equation*}
$$

This regularity property is sufficient to justify all computations which follow.
It is convenient to introduce the notations

$$
\begin{equation*}
M:=\left[u^{0}\right], \quad v:=u-M, \quad v^{0}:=u^{0}-M \tag{12}
\end{equation*}
$$

then we deduce from $(1),(5),(7)$ and (11) that

$$
\begin{gather*}
v \in C\left(\mathbb{R}_{+} ; H_{p}^{5}\right),  \tag{13}\\
v \in C^{1}\left(\mathbb{R}_{+} ; H_{p}^{2}\right),  \tag{14}\\
{[v(t)]=0, \forall t \in \mathbb{R}_{+},}  \tag{15}\\
v^{\prime}+v v_{x}+M v_{x}+v_{x x x}+k v=0 \quad \text { in } \Omega \times \mathbb{R}_{+},  \tag{16}\\
v(0)=v^{0} \quad \text { on } \Omega, \tag{17}
\end{gather*}
$$

and the estimates (10) take the following form :

$$
\begin{equation*}
\left\|\left(v(t), v^{\prime}(t)\right)\right\|_{H_{p}^{2} \times H_{p}^{-1}} \leq C e^{-k^{\prime} t}, \quad \forall t \in \mathbb{R}_{+} \tag{18}
\end{equation*}
$$

Lemma 10.3. - The function

$$
\begin{equation*}
t \mapsto \int v(t)^{2} \mathrm{dx}, \quad t \in \mathbb{R}_{+} \tag{19}
\end{equation*}
$$

is continuously differentiable and

$$
\begin{equation*}
\left(\int v^{2} \mathrm{dx}\right)^{\prime} \equiv-2 k \int v^{2} \mathrm{dx} \tag{20}
\end{equation*}
$$

Proof. - Since $H_{p}^{2}$ is a Banach algebra, it follows from (14) that the function $v^{2}$ is continuously differentiable. Hence the function (19), being the composition of two function of class $C^{2}$, is also continuously differentiable.

Using (16) and the periodicity of $v$ (cf. (13)) we easily obtain the identity (20) :

$$
\begin{gathered}
\left(\int v^{2} \mathrm{dx}\right)^{\prime}=\int 2 v v^{\prime} \mathrm{dx}=\int-2 v\left(v v_{x}+M v_{x}+v_{x x x}+k v\right) \mathrm{dx} \\
=\int-2 v^{2} v_{x}-2 M v v_{x}-2 v v_{x x x}-2 k v^{2} \mathrm{dx} \\
=\left[-(2 / 3) v^{3}-M v^{2}-2 v v_{x x}+v_{x}^{2}\right]_{0}^{1}-2 k \int v^{2} \mathrm{dx}=-2 k \int v^{2} \mathrm{dx} .
\end{gathered}
$$

Lemma 10.4. - The function

$$
\begin{equation*}
t \mapsto \int v_{x}(t)^{2}-(1 / 3) v^{3}(t)^{2} \mathrm{dx}, \quad t \in \mathbb{R}_{+} \tag{21}
\end{equation*}
$$

is continuously differentiable and

$$
\begin{equation*}
\left(\int v_{x}(t)^{2}-(1 / 3) v^{3}(t)^{2} \mathrm{dx}\right)^{\prime} \equiv-2 k \int v_{x}(t)^{2}-(1 / 2) v^{3}(t)^{2} \mathrm{dx} \tag{22}
\end{equation*}
$$

Proof. - It follows easily from (14) that the function (21) is continuously differentiable. Using (16) and the periodicity of $v$ hence the identity (22) follows :

$$
\begin{gathered}
\left(\int v_{x}^{2}-(1 / 3) v^{3} \mathrm{dx}\right)^{\prime}=\int 2 v_{x} v_{x}^{\prime}-v^{2} v^{\prime} \mathrm{dx} \\
=\left[2 v_{x} v^{\prime}\right]_{0}^{1}+\int-v^{\prime}\left(2 v_{x x}+v^{2}\right) \mathrm{dx}=\int\left(v v_{x}+M v_{x}+v_{x x x}+k v\right)\left(2 v_{x x}+v^{2}\right) \mathrm{dx} \\
=\int\left(v^{2}\right)_{x} v_{x x}+\left(M v_{x}^{2}+v_{x x}^{2}+(1 / 4) v^{4}+(M / 3) v^{3}\right)_{x}+v^{2} v_{x x x}-2 k v_{x}^{2}+k v^{3} \mathrm{dx} \\
=\left[v^{2} v_{x x}+M v_{x}^{2}+v_{x x}^{2}+(1 / 4) v^{4}+(M / 3) v^{3}\right]_{0}^{1}+\int-2 k v_{x}^{2}+k v^{3} \mathrm{dx} \\
=-2 k \int-(1 / 2) v^{3} \mathrm{dx} .
\end{gathered}
$$

## Lemma 10.5. - The function

$$
\begin{equation*}
t \mapsto \int v_{x x}(t)^{2}-(5 / 3) v_{x}^{2} v(t)^{2}+(5 / 36) v(t)^{4} \mathrm{dx}, \quad t \in \mathbb{R}_{+} \tag{23}
\end{equation*}
$$

is continuously differentiable and

$$
\begin{align*}
& \left(\int v_{x x}^{2}-(5 / 3) v_{x}^{2} v+(5 / 36) v^{4} \mathrm{dx}\right)^{\prime} \\
\equiv & -2 k \int v_{x x}^{2}-(5 / 2) v_{x}^{2} v+(5 / 18) v^{4} \mathrm{dx} \tag{24}
\end{align*}
$$

Proof. - By (14) the function (23) is continuously differentiable. To show the identity (24) first we deduce from (16), using the periodicity of $v$, the following identity :

$$
\begin{gathered}
\left(\int v_{x x}^{2}-(5 / 3) v_{x}^{2} v+(5 / 36) v^{4} \mathrm{dx}\right)^{\prime} \\
=\int 2 v_{x x} v_{x x}^{\prime}-(10 / 3) v_{x} v_{x}^{\prime} v-(5 / 3) v_{x}^{2} v^{\prime}+(5 / 9) v^{3} v^{\prime} \mathrm{dx} \\
=\left[2 v_{x x} v_{x}^{\prime}-2 v_{x x x} v^{\prime}-(10 / 3) v v_{x} v^{\prime}\right]_{0}^{1} \\
+\int v^{\prime}\left(2 v_{x x x x}+(5 / 3) v_{x}^{2}+(10 / 3) v v_{x x}+(5 / 9) v^{3}\right) \mathrm{dx} \\
=\int v^{\prime}\left(2 v_{x x x x}+(5 / 3) v_{x}^{2}+(10 / 3) v v_{x x}+(5 / 9) v^{3}\right) \mathrm{dx} \\
=-\int\left(v v_{x}+M v_{x}+v_{x x x}+k v\right)\left(2 v_{x x x x}+(5 / 3) v_{x}^{2}+(10 / 3) v v_{x x}+(5 / 9) v^{3}\right) \mathrm{dx} \\
=-k \int 2 v v_{x x x x}+(5 / 3) v v_{x}^{2}+(10 / 3) v^{2} v_{x x}+(5 / 9) v^{4} \mathrm{dx} \\
-\int 2 v v_{x} v_{x x x x}+(5 / 3) v v_{x}^{3}+(10 / 3) v^{2} v_{x} v_{x x}+(5 / 9) v^{4} v_{x}+2 v_{x x x} v_{x x x x} \\
+(5 / 3) v_{x}^{2} v_{x x x}+(10 / 3) v v_{x x} v_{x x x}+(5 / 9) v^{3} v_{x x x} \mathrm{dx} \\
=:-k I_{1}-M I_{2}-I_{3} .
\end{gathered}
$$

It suffices to show that

$$
I_{1}=\int 2 v_{x x}^{2}-5 v v_{x}^{2}+(5 / 9) v^{4} \mathrm{dx} \quad \text { and } \quad I_{2}=I_{3}=0
$$

We have

$$
\begin{gathered}
I_{1}=\left[2 v v_{x x x}-2 v_{x} v_{x x}+(10 / 3) v^{2} v_{x}\right]_{0}^{1} \\
+\int 2 v_{x x}^{2}+(5 / 3) v v_{x}^{2}-(20 / 3) v v_{x}^{2}+(5 / 9) v^{4} \mathrm{dx}=\int 2 v_{x x}^{2}-5 v v_{x}^{2}+(5 / 9) v^{4} \mathrm{dx}
\end{gathered}
$$

and
$I_{2}=\left[2 v_{x} v_{x x x}-v_{x x}^{2}+(5 / 3) v v_{x}^{2}+(5 / 36) v^{4}\right]_{0}^{1}+\int(5 / 3) v_{x}^{3}-(5 / 3) v_{x}^{3} \mathrm{dx}=0$.
Finally, we have

$$
\begin{aligned}
& I_{3}=\int 2 v v_{x} v_{x x x x}+(5 / 3) v v_{x}^{3}+(5 / 3) v_{x}^{2} v_{x x x}+(10 / 3) v^{2} v_{x} v_{x x}+(10 / 3) v v_{x x} v_{x x x} \\
& +(5 / 9) v^{4} v_{x}+(5 / 9) v^{3} v_{x x x}+2 v_{x x x} v_{x x x x} \mathrm{dx} \\
& =\left[2 v v_{x} v_{x x x}+(1 / 9) v^{5}\right]_{0}^{1}+\int-2 v_{x}^{2} v_{x x x}-2 v v_{x x} v_{x x x}+(5 / 3) v v_{x}^{3} \\
& + \\
& =\int-2 v_{x}^{2} v_{x x x}-2 v v_{x x}^{2} v_{x x x}+(10 / 3) v^{2} v_{x x x} v_{x x}+(10 / 3) v v_{x x} v_{x x x}+(5 / 9) v^{3} v_{x x x} \mathrm{dx} \\
& \\
& \quad+(10 / 3) v v_{x x} v_{x x x}+(5 / 9) v^{3} v_{x x x}^{3}+(5 / 3) v_{x}^{2} v_{x x x}+(5 / 9) v^{3} v_{x x} \mathrm{dx} \\
& =\left[-2 v_{x}^{2} v_{x x}-v v_{x x}^{2}+(5 / 3) v v_{x}^{2}+(5 / 3) v v_{x x}^{2} v_{0}^{1}+\int 4 v_{x} v_{x x}^{2}+v_{x} v_{x x}^{2}\right. \\
& +(5 / 3) v v_{x}^{3}-(10 / 3) v_{x} v_{x x}^{2}+(10 / 3) v^{2} v_{x} v_{x x}-(5 / 3) v_{x} v_{x x}^{2}-(5 / 3) v^{2} v_{x} v_{x x} \mathrm{dx} \\
& \quad=\int(5 / 3) v v_{x}^{3}+(5 / 3) v^{2} v_{x} v_{x x} \mathrm{dx}=0 .
\end{aligned}
$$

In order to simplify the notation we shall write $\|\cdot\|_{p}$ for the norm of $L^{p}(\Omega)$, $1 \leq p \leq \infty$. Since $\Omega$ is the unit interval, the Hölder inequality is particularly simple :

$$
\begin{equation*}
\|v\|_{p} \leq\|v\|_{q}, \quad \forall v \in L^{q}(\Omega), \quad 1 \leq p \leq q \leq \infty . \tag{25}
\end{equation*}
$$

We shall also use the Poincaré-Wirtinger inequality :

$$
\begin{equation*}
\|v\|_{\infty} \leq\left\|v_{x}\right\|_{1}, \quad \text { if } \quad v \in H^{1}(\Omega) \quad \text { and } \quad[v]=0 \tag{26}
\end{equation*}
$$

The proof is simple : since $v$ is continuous, there exists $a \in \Omega$ such that $v(a)=0$. Then for any $y \in \Omega$ we have

$$
|v(y)|=|v(y)-v(a)|=\left|\int_{a}^{y} v_{x} \mathrm{dx}\right| \leq \int_{\Omega}\left|v_{x}\right| \mathrm{dx}=\left\|v_{x}\right\|_{1} .
$$

Noe that lemma 10.3 implies that

$$
\begin{equation*}
\|v(t)\|_{2}=\left\|v^{0}\right\|_{2} e^{-k t}, \quad \forall t \in \mathbb{R}_{+} \tag{27}
\end{equation*}
$$

Now let us show that for each fixed $k^{\prime} \in(0, k)$ there exists a positive constant $C^{\prime}$ such that

$$
\begin{equation*}
\left\|v_{x}(t)\right\|_{2}=C^{\prime} e^{-k^{\prime} t}, \quad \forall t \geq 0 \tag{28}
\end{equation*}
$$

Using (25)-(27) we have

$$
\left\|v^{3}(t)\right\|_{1} \leq\|v(t)\|_{\infty}^{2}\|v(t)\|_{1} \leq\left\|v_{x}(t)\right\|_{1}^{2}\|v(t)\|_{2} \leq\left\|v_{x}(t)\right\|_{2}^{2}\left\|v^{0}\right\|_{2} e^{-k t}
$$

consequently, for any fixed $\varepsilon>0$ (to be chosen later) there exists $T^{\prime}>0$ such that

$$
\begin{equation*}
\int v^{3}(t) \mathrm{dx} \leq \varepsilon \int v_{x}^{2} \mathrm{dx}, \quad \forall t>T^{\prime} \tag{29}
\end{equation*}
$$

If $\varepsilon \leq 2$, then we deduce from (29) the inequalities

$$
\begin{equation*}
\int\left(v_{x}^{2}-(1 / 3) v^{3}\right)(t) \mathrm{dx} \geq(1 / 3) \int v_{x}^{2}(t) \mathrm{dx} \geq 0, \quad \forall t>T^{\prime} \tag{30}
\end{equation*}
$$

If $\varepsilon$ is sufficiently small, then we also deduce from (29) that

$$
\begin{equation*}
-2 k \int\left(v_{x}^{2}-(1 / 2) v^{3}\right)(t) \mathrm{dx} \leq-2 k^{\prime} \int\left(v_{x}^{2}-(1 / 3) v^{3}\right)(t) \mathrm{dx}, \quad \forall t>T^{\prime} \tag{31}
\end{equation*}
$$

(It suffices to choose $\varepsilon \leq\left(6 k-6 k^{\prime}\right) /\left(3 k-2 k^{\prime}\right)$.)
Thus, choosing a sufficiently small $\varepsilon$ we deduce from (21), (30) and (31) that

$$
\begin{gathered}
(1 / 3) \int v_{x}^{2}(t) \mathrm{dx} \leq \int\left(v_{x}^{2}-(1 / 3) v^{3}\right)(t) \mathrm{dx} \\
\leq \int\left(v_{x}^{2}-(1 / 3) v^{3}\right)(0) \mathrm{dx} e^{-k^{\prime}\left(t-T^{\prime}\right)}=: C^{\prime} e^{-k^{\prime} t}, \quad \forall t>T^{\prime}
\end{gathered}
$$

which implies (28) for all $t>T^{\prime}$. The left-hand side of (28) being continuous, the estimate (28) remains valid for all $t \geq 0$ with some bigger constant $C^{\prime}$.

Next we show similarly that for any fixed $k^{\prime}<k$ there exists a positive constant $C^{\prime \prime}$ such that

$$
\begin{equation*}
\left\|v_{x x}\right\|_{2}=C^{\prime \prime} e^{-k^{\prime} t}, \quad \forall t \geq 0 \tag{32}
\end{equation*}
$$

Using (25)-(28) we have

$$
\left\|v^{4}(t)\right\|_{1} \leq\|v(t)\|_{\infty}^{2}\|v(t)\|_{2}^{2} \leq\left\|v_{x}(t)\right\|_{2}^{2}\|v(t)\|_{2}^{2} \leq\left\|v_{x x}(t)\right\|_{2}^{2}\left\|v^{0}\right\|_{2}^{2} e^{-2 k t}
$$

and

$$
\left|\int\left(v_{x}^{2} v\right)(t) \mathrm{dx}\right| \leq\left\|v_{x}(t)\right\|_{2}^{2}\|v(t)\|_{\infty} \leq\left\|v_{x}(t)\right\|_{2}^{3} \leq\left\|v_{x x}(t)\right\|_{2}^{2} C^{\prime} e^{-k^{\prime} t}
$$

It follows that for any fixed $\varepsilon>0$ (to be chosen later) there exists $T^{\prime \prime}>0$ such that

$$
\begin{equation*}
\int\left|\left(v_{x}^{2} v\right)(t)\right|+v(t)^{4} \mathrm{dx} \leq \varepsilon \int v_{x x}(t)^{2} \mathrm{dx}, \quad \forall t>T^{\prime \prime} \tag{32}
\end{equation*}
$$

Choosing $\varepsilon>0$ sufficiently small we conclude from (32) that

$$
\begin{equation*}
\int\left(v_{x x}^{2}-(5 / 3) v_{x}^{2} v+(5 / 36) v^{4}\right)(t) \mathrm{dx} \geq(1 / 3) \int v_{x x}(t)^{2} \mathrm{dx} \geq 0 \tag{33}
\end{equation*}
$$

and

$$
\begin{align*}
& -2 k \int\left(v_{x x}^{2}-(5 / 2) v_{x}^{2} v+(5 / 18) v^{4}\right)(t) \mathrm{dx}  \tag{34}\\
\leq & -2 k^{\prime \prime} \int\left(v_{x x}^{2}-(5 / 3) v_{x}^{2} v+(5 / 36) v^{4}\right)(t) \mathrm{dx}
\end{align*}
$$

for all $t>T^{\prime \prime}$. We deduce from (23), (33) and (34) that

$$
\begin{aligned}
& (1 / 3) \int v_{x x}(t)^{2} \mathrm{dx} \leq \int\left(v_{x x}^{2}-(5 / 3) v_{x}^{2} v+(5 / 36) v^{4}\right)(t) \mathrm{dx} \\
\leq & \int\left(v_{x x}^{2}-(5 / 3) v_{x}^{2} v+(5 / 36) v^{4}\right)(0) \mathrm{dxx} \\
-2 k^{\prime}\left(t-T^{\prime}\right) & =:\left(C^{\prime}\right)^{2} e^{-2 k^{\prime} t}
\end{aligned}
$$

proving (32) for all $t>T^{\prime}$. The left-hand side of (32) being continuous, the estimate (32) remains valid for every $t \geq 0$ if we choose some larger constant $C^{\prime \prime}$.

Now we may easily complete the proof of the theorem. By (27), (28) and (32) for every fixed $k^{\prime}<k$ there exists a positive constant $C_{1}>0$ such that

$$
\begin{equation*}
\|v(t)\|_{H_{p}^{2}} \leq C_{1} e^{-k^{\prime} t}, \quad \forall t \geq 0 \tag{35}
\end{equation*}
$$

Using the equation (16) hence we conclude easily that

$$
\begin{equation*}
\left\|u^{\prime}(t)\right\|_{H_{p}^{-1}} \leq C_{2} e^{-k^{\prime} t}, \quad \forall t \geq 0 \tag{36}
\end{equation*}
$$

with some constant $C_{2}>0$. The estimate (18) follows from (35) and (36).

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