# THE EXPONENTIAL STABILITY OF THE PROBLEM OF TRANSMISSION OF THE WAVE EQUATION 

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#### Abstract

The problem of exponential stability of the problem of transmission of the wave equation with lower-order terms is considered. Making use of the classical energy method and multiplier technique, we prove that this problem of transmission is exponentially stable.


## 1. Introduction

Throughout this paper, let $\Omega$ be a bounded domain (open, nonempty, and connected) in $\mathbb{R}^{n}(n \geqslant 1)$ with a boundary $\Gamma=\partial \Omega$ of class $C^{2}$ which consists of two parts, $S_{1}$ and $S_{2}$ (see Figure 1 below). $S_{1}$ is assumed to be either empty or to have a nonempty interior and $S_{2} \neq \emptyset$ and relatively open in $\Gamma$. Assume $\bar{S}_{1} \cap \bar{S}_{2}=\emptyset$. Let $S_{0}$ with $\bar{S}_{0} \cap \bar{S}_{1}=\bar{S}_{0} \cap \bar{S}_{2}=\emptyset$ be a regular hypersurface of class $C^{2}$, which separates $\Omega$ into two domains, $\Omega_{1}$ and $\Omega_{2}$, such that $S_{1} \subset \Gamma_{1}=\partial \Omega_{1}$ and $S_{2} \subset \Gamma_{2}=\partial \Omega_{2}$. For $T>0$, set $Q=\Omega \times(0, T), Q_{1}=\Omega_{1} \times(0, T), Q_{2}=\Omega_{2} \times(0, T), \Sigma_{i}=S_{i} \times(0, T)(i=0,1,2)$. The following figure is a typical domain of this kind.


Figure 1

[^0]In this paper we shall be concerned with the problem of rate of exponential decay of energy for the problem of transmission of the wave equation with lower-order terms and with dissipative boundary condition of Robin type:

$$
\begin{cases}u_{i}^{\prime \prime}-a_{i} \Delta u_{i}+q u_{i}=0 & \text { in } \Omega_{i} \times(0, \infty),  \tag{1.1}\\ u_{i}(x, 0)=u_{i}^{0}(x), \quad u_{i}^{\prime}(x, 0)=u_{i}^{1}(x) & \text { in } \Omega_{i}, \quad i=1,2, \\ u_{1}=0 & \text { on } S_{1} \times(0, \infty), \\ \frac{\partial u_{2}}{\partial \nu}+\alpha(x) u_{2}+\sigma(x) u_{2}^{\prime}=0, & \text { on } S_{2} \times(0, \infty), \\ u_{1}=u_{2}, \quad a_{1} \frac{\partial u_{1}}{\partial \nu}=a_{2} \frac{\partial u_{2}}{\partial \nu} & \text { on } S_{0} \times(0, \infty) .\end{cases}
$$

In (1.1), $\nu$ denotes the unit normal on $\Gamma$ and $S_{0}$ directed towards the exterior of $\Omega$ and $\Omega_{1}, a_{1}$ and $a_{2}$ are positive constants, the functions $q: \Omega \rightarrow \mathbb{R}, \alpha, \sigma: S_{2} \rightarrow \mathbb{R}$ are nonnegative and satisfy

$$
\begin{equation*}
q \in L^{\infty}(\Omega), \quad \alpha, \sigma \in C^{1}\left(S_{2}\right) \tag{1.2}
\end{equation*}
$$

There has been extensive work on energy decay for the wave equation. The poineering work (see [14, 18]) was first performed in the mid-seventies in studies aimed at achieving energy decay rates for the wave equation exterior to a bounded obstacle (the so-called "exterior" problem). In contrast, the "interior" problem is more difficult than the "exterior" problem, since the latter enjoys the advantage that the energy distributes itself over an infinite region as $t \rightarrow \infty$. Russell [17] made a conjecture in 1974 concerning uniform energy decay rates for the interior problem. This conjecture was verified by Chen (see $[\mathbf{2}, \mathbf{3}, \mathbf{4}, \mathbf{5}, 6]$ ) under some natural geometrical conditions on $\Omega$. Lagnese [ 9 ] further relaxed the geometrical conditions on $\Omega$ by obtaining a key inequality which is of independent interest. More recently, Bardos, Lebeau, and Rauch [1] considered general second order hyperbolic equations but with smooth coefficients.

We note that in the previous work the coefficients of the equation are required to be sufficiently smooth and it seems that the problem of transmission has not been considered yet. Therefore, by applying the classical energy method and multiplier technique, we here discuss this problem and generalise some known results to the case of transmission.

Set

$$
u=\left\{\begin{array}{ll}
u_{1}, & x \in \Omega_{1}, \\
u_{2}, & x \in \Omega_{2},
\end{array} \quad u^{0}=\left\{\begin{array}{ll}
u_{1}^{0}, & x \in \Omega_{1}, \\
u_{2}^{0}, & x \in \Omega_{2},
\end{array} \quad u^{1}= \begin{cases}u_{1}^{1}, & x \in \Omega_{1} \\
u_{2}^{1}, & x \in \Omega_{2}\end{cases}\right.\right.
$$

We define the energy of system (1.1) by

$$
\begin{aligned}
& E(u, t)=\frac{1}{2} \int_{\Omega}\left[\left|u^{\prime}(x, t)\right|^{2}+a(x)|\nabla u(x, t)|^{2}+q(x)|u(x, t)|^{2}\right] d x \\
&+\frac{1}{2} \int_{S_{2}} a(x) \alpha(x)|u|^{2} d \Gamma
\end{aligned}
$$

Let $H^{s}(\Omega)$ always denote the usual Sobolev space and $\|\cdot\|_{s, \Omega}$ its norm for any $s \in \mathbb{R}$. Let

$$
L^{2}\left(\Omega, S_{1}\right)= \begin{cases}\left\{u \in L^{2}(\Omega): \int_{\Omega} u(x) d x=0\right\}, & \text { if } S_{1}=\emptyset, q \equiv 0, \text { and } \alpha \equiv 0  \tag{1.3}\\ L^{2}(\Omega), & \text { otherwise } ;\end{cases}
$$

(1.4) $H^{1}\left(\Omega, S_{1}\right)= \begin{cases}\left\{u \in H^{1}(\Omega): \int_{\Omega} u(x) d x=0\right\}, & \text { if } S_{1}=\emptyset, q \equiv 0, \text { and } \alpha \equiv 0, \\ \left\{u \in H^{1}(\Omega): u=0 \text { on } S_{1}\right\}, & \text { otherwise; }\end{cases}$

$$
H_{S_{1}}^{1}(\Omega)= \begin{cases}H^{1}(\Omega), & \text { if } S_{1}=\emptyset  \tag{1.5}\\ \left\{u \in H^{1}(\Omega): u=0 \text { on } S_{1}\right\}, & \text { otherwise }\end{cases}
$$

The main result of this paper is as follows.
Theorem 1.1. Let $\nu$ denote the unit normal on $\Gamma$ and $S_{0}$ directed towards the exterior of $\Omega$ and $\Omega_{1}$. Assume there is a vector field $l(x)=\left(l_{1}(x), \cdots, l_{n}(x)\right)$ of class $C^{2}(\bar{\Omega})$ such that
(i) l. $\nu \leqslant 0$ almost everywhere on $S_{1}$ with respect to the (n-1)-dimensional surface measure;
(ii) $l \cdot \nu \geqslant \eta>0$ almost everywhere on $S_{2}$ with respect to the (n-1)dimensional surface measure;
(iii) $\left(a_{1}-a_{2}\right) l \cdot \nu \geqslant 0$ almost everywhere on $S_{0}$ with respect to the (n-1)dimensional surface measure;
(iv) the matrix $\left(\frac{\partial l_{i}}{\partial x_{j}}+\frac{\partial l_{j}}{\partial x_{i}}\right)$ is uniformly positive definite on $\Omega$;
(v) there exists a constant $\sigma_{0}>0$ such that

$$
\sigma \geqslant \sigma_{0} \quad \text { on } S_{2}
$$

Then there are positive constants $M, \tau$ such that

$$
\begin{equation*}
E(u, t) \leqslant M e^{-\tau t} E(u, 0), \text { for all } t \geqslant 0 \tag{1.6}
\end{equation*}
$$

for all solutious $u$ of (1.1) with $\left(u^{0}, u^{1}\right) \in H_{S_{1}}^{1}(\Omega) \times L^{2}\left(\Omega, S_{1}\right)$.
In the proof of Theorem 1.1 below, condition (iii) is crucial. Whether Theorem 1.1 still holds if condition (iii) fails is an open problem. The vector field $l(x)$ was first introduced in [4] and further improved in [9]. We here give an example of $l(x)$ which satisfies conditions (i)-(iv) of Theorem 1.1. Let $\Omega=\left\{x \in \mathbb{R}^{2}: 1<|x|<3\right\}$ and $S_{0}=\left\{x \in \mathbb{R}^{2}:|x|=2\right\}$. Then $S_{1}=\left\{x \in \mathbb{R}^{2}:|x|=1\right\}$ and $S_{2}=\left\{x \in \mathbb{R}^{2}:|x|=3\right\}$. It is easy to see that $l(x)=x$ is a vector field as required.

In comparison with existing results, Theorem 1.1 generalises the result of Lagnese [9] to the case of transmission with Robin boundary conditions. Also, it generalises [8, Theorem 8.15, p.117] in three aspects: firstly, the vector field $m(x)=x-x^{0}=$ ( $x_{1}-x_{1}^{0}, \cdots, x_{n}-x_{n}^{0}$ ) in the previous theorems is replaced by the more general vector field $l(x)$; secondly the condition $\min _{S_{2}} \alpha(x)>0$ of [ 8 , Theorem 8.15] has been moved off; thirdly, we have considered the problem of transmission. In addition, the most interesting part of this paper may be the strategy for handling the case where $\alpha$ is not necessarily small.

The rest of this paper is divided into two sections. In Section 2, we briefly discuss the well-posedness of problem (1.1) via the theory of semigroups of linear bounded operators. In Section 3, we prove Theorem 1.1.

## 2. Well-posedness

The well-posedness of problem (1.1) is by now well known in the case where $a_{1}=a_{2}$ (see [2, 10], [11, p.137-139]), and can be similarly treated without any difficulty in the case where $a_{1} \neq a_{2}$. For completeness, we give an outline.

Set

$$
u=\left\{\begin{array}{ll}
u_{1}, & x \in \Omega_{1},  \tag{2.1}\\
u_{2}, & x \in \Omega_{2},
\end{array} \quad u^{0}=\left\{\begin{array}{ll}
u_{1}^{0}, & x \in \Omega_{1}, \\
u_{2}^{0}, & x \in \Omega_{2},
\end{array} \quad u^{1}= \begin{cases}u_{1}^{1}, & x \in \Omega_{1}, \\
u_{2}^{1}, & x \in \Omega_{2},\end{cases}\right.\right.
$$

$$
a(x)= \begin{cases}a_{1}, & x \in \Omega_{1}  \tag{2.2}\\ a_{2}, & x \in \Omega_{2}\end{cases}
$$

In the sequel, $u, u^{0}, u^{1}$ are always as in (2.1); an integral of $u$ on a domain $\Omega$ means the sum of two integrals of $u_{1}$ and $u_{2}$ on the subdomains $\Omega_{1}$ and $\Omega_{2}$; that an equation related to $u$ holds on a domain $\Omega$ means that the equation holds on the subdomains $\Omega_{1}$ and $\Omega_{2}$, respectively.

Problem (1.1) can be formulated as an abstract Cauchy problem:

$$
\left\{\begin{array}{l}
\binom{u^{\prime}}{v^{\prime}}=\left(\begin{array}{cc}
0 & I \\
a(x) \Delta-q & 0
\end{array}\right)\binom{u}{v}=A\binom{u}{v}  \tag{2.3}\\
\binom{u(0)}{v(0)}=\binom{u^{0}}{u^{1}}
\end{array}\right.
$$

in the Hilbert space

$$
\mathcal{H}_{1}=H^{1}\left(\Omega, S_{1}\right) \times L^{2}\left(\Omega, S_{1}\right)
$$

for an initial condition $\left(u^{0}, u^{1}\right)$ with

$$
D(A)=\left\{(u, v):(u, v) \in H^{2}\left(\Omega_{1}, \Omega_{2}, S_{1}\right) \times H^{1}\left(\Omega, S_{1}\right), \frac{\partial u}{\partial \nu}+\alpha u+\sigma v=0 \text { on } S_{2}\right\} .
$$

The spaces used for these definition are given by (1.3)-(1.5). In addition,

$$
\begin{aligned}
& H_{S_{1}}^{2}\left(\Omega_{1}, \Omega_{2}\right)=\left\{u \in H^{1}(\Omega): u_{i} \in H^{2}\left(\Omega_{i}\right), i=1,2 ; u=0 \text { on } S_{1}\right. \\
& \left.\quad \text { and } a_{1} \frac{\partial u_{1}}{\partial \nu}=a_{2} \frac{\partial u_{2}}{\partial \nu} \text { on } S_{0}\right\}
\end{aligned}
$$

$$
\begin{aligned}
& H^{2}\left(\Omega_{1}, \Omega_{2}, S_{1}\right) \\
& \quad= \begin{cases}\left\{u \in H_{S_{1}}^{2}\left(\Omega_{1}, \Omega_{2}\right): \sum_{i=1}^{2} \int_{\Omega_{i}} u_{i} d x=0, \sum_{i=1}^{2} \int_{\Omega_{i}}\left(a_{i} \Delta u_{i}-q u_{i}\right) d x=0\right\} \\
& \text { if } S_{1}=\emptyset, q \equiv 0, \text { and } \alpha \equiv 0 \\
H_{S_{1}}^{2}\left(\Omega_{1}, \Omega_{2}\right), & \text { otherwise. }\end{cases}
\end{aligned}
$$

Note that $H^{2}\left(\Omega_{1}, \Omega_{2}, S_{1}\right) \subset H^{1}\left(\Omega, S_{1}\right)$ because $u_{1}=u_{2}$ on $S_{0}$.
If $S_{1}=\emptyset, q \equiv 0$, and $\alpha \equiv 0, L^{2}\left(\Omega, S_{1}\right)$ in $\mathcal{H}_{1}$ cannot be replaced by $L^{2}(\Omega)$ since $H^{1}\left(\Omega, S_{1}\right)$ is not dense in $L^{2}(\Omega)$.

In the sequel, we always use the energy scalar product on $\mathcal{H}_{1}$ :

$$
\begin{equation*}
\left\langle\binom{ u}{v},\binom{\phi}{\psi}\right\rangle=\int_{\Omega}[a(x) \nabla u \cdot \nabla \bar{\phi}+q u \bar{\phi}+v \bar{\psi}] d x+\int_{S_{2}} a(x) \alpha(x) u \bar{\phi} d \Gamma \tag{2.5}
\end{equation*}
$$

which is equivalent to the scalar product on $\mathcal{H}_{1}$ induced by $H^{1}(\Omega) \times L^{2}(\Omega)$.
As done in [2] or [11, p.137-139], it is easy to verify that the operator $A$ is the infinitesimal generator of a strongly continuous semigroup of contractions on $\mathcal{H}_{1}$.

We define the energy of system (1.1) by

$$
\begin{align*}
& E(u, t)=\frac{1}{2} \int_{\Omega}\left[\left|u^{\prime}(x, t)\right|^{2}+a(x)|\nabla u(x, t)|^{2}+q(x)|u(x, t)|^{2}\right] d x  \tag{2.6}\\
&+\frac{1}{2} \int_{S_{2}} a(x) \alpha(x)|u|^{2} d \Gamma
\end{align*}
$$

Let $X$ be a Banach space. We denote by $C^{k}([0, T] ; X)$ the space of all $k$ times continuously differentiable functions defined on $[0, T]$ with values in $X$, and write $C([0, T] ; X)$ for $C^{0}([0, T] ; X)$.

Now an application of the theory of semigroups [16, Chapter 1] gives
Theorem 2.1. (i) For any initial condition $\left(u^{0}, u^{1}\right) \in H_{S_{1}}^{1}(\Omega) \times L^{2}\left(\Omega, S_{1}\right)$, problem (1.1) has a unique weak solution with

$$
\begin{equation*}
u \in C\left([0, \infty) ; H_{S_{1}}^{1}(\Omega)\right) \cap C^{1}\left([0, \infty) ; L^{2}\left(\Omega, S_{1}\right)\right) \tag{2.7}
\end{equation*}
$$

and

$$
\begin{equation*}
\frac{\partial u}{\partial \nu}, \quad u^{\prime} \in L^{2}\left(\Sigma_{2}\right) \tag{2.8}
\end{equation*}
$$

Moreover,

$$
\begin{equation*}
E(u, t) \leqslant E(u, 0), \quad \forall t \geqslant 0, \tag{2.9}
\end{equation*}
$$

and there exists a constant $c=c(T)>0$ such that

$$
\begin{equation*}
\left\|\frac{\partial u}{\partial \nu}\right\|_{L^{2}\left(\Sigma_{2}\right)}^{2}+\left\|u^{\prime}\right\|_{L^{2}\left(\Sigma_{2}\right)}^{2} \leqslant c E(u, 0) \tag{2.10}
\end{equation*}
$$

(ii) For any initial condition $\left(u^{0}, u^{1}\right) \in D(A)$, problem (1.1) has a unique strong solution with

$$
\begin{equation*}
u \in C\left([0, \infty) ; H^{2}\left(\Omega_{1}, \Omega_{2}, S_{1}\right)\right) \cap C^{1}\left([0, \infty) ; H^{1}\left(\Omega, S_{1}\right)\right) \tag{2.11}
\end{equation*}
$$

Moreover, there exists a constant $c=c(T)>0$ such that for all $t \in[0, T]$

$$
\begin{equation*}
\left\|u^{\prime}(t)\right\|_{1, \Omega}+\|u(t)\|_{1, \Omega}+\sum_{i=1}^{2}\left\|\Delta u_{i}(t)\right\|_{0, \Omega_{i}} \leqslant c\left[\left\|u^{1}\right\|_{1, \Omega}+\left\|u^{0}\right\|_{1, \Omega}+\sum_{i=1}^{2}\left\|\Delta u_{i}^{0}\right\|_{0, \Omega_{i}}\right] \tag{2.12}
\end{equation*}
$$

Proof: If $S_{1} \neq \emptyset$, or $q \not \equiv 0$, or $\alpha \not \equiv 0$, then $H_{S_{1}}^{1}(\Omega)=H^{1}\left(\Omega, S_{1}\right)$. It therefore follows from semigroup theory that problem (1.1) has a unique weak solution $u$ with (2.7) for $\left(u^{0}, u^{1}\right) \in H_{S_{1}}^{1}(\Omega) \times L^{2}\left(\Omega, S_{1}\right)$, and with (2.11) for $\left(u^{0}, u^{1}\right) \in D(A)$. On the other hand, by multiplying the first equation of (1.1) by $u_{i}^{\prime}$ and integrating over $\Omega \times(0, T)$, we obtain

$$
\begin{equation*}
0 \leqslant \int_{\Sigma_{2}} a(x) \sigma(x)\left|u^{\prime}\right|^{2} d \Sigma=E(u, 0)-E(u, T) \leqslant E(u, 0) \tag{2.13}
\end{equation*}
$$

This gives (2.9). In addition, by (1.1) we have

$$
\begin{equation*}
\int_{\Sigma_{2}}\left|\frac{\partial u}{\partial \nu}\right|^{2} d \Sigma=\int_{\Sigma_{2}}\left(\alpha u+\sigma u^{\prime}\right)^{2} d \Sigma \tag{2.14}
\end{equation*}
$$

and by the trace theorem [12, p.39], we have

$$
\begin{equation*}
\int_{\Sigma_{2}}|u|^{2} d \Sigma \leqslant c E(u, 0) \tag{2.15}
\end{equation*}
$$

Thus, (2.8) and (2.10) follow from (2.13)-(2.15). To prove (2.12), let $T(t)$ be the semigroup generated by $A$. Then the solution $u$ of (1.1) can be expressed as

$$
\binom{u}{u^{\prime}}=T(t)\binom{u^{0}}{u^{1}} .
$$

If $\left(u^{0}, u^{1}\right) \in D(A)$, by the semigroup property [16, p.4] we have

$$
A\binom{u}{u^{\prime}}=T(t) A\binom{u^{0}}{u^{1}}
$$

that is,

$$
\binom{u^{\prime}}{a(x) \Delta u-q u}=T(t)\binom{u^{1}}{a(x) \Delta u^{0}-q u^{0}} .
$$

This yields
(2.16) $\quad\left\|u^{\prime}(t)\right\|_{1, \Omega}+\sum_{i=1}^{2}\left\|\Delta u_{i}(t)\right\|_{0, \Omega_{i}}$

$$
\leqslant c\left[\left\|u^{1}\right\|_{1, \Omega}+\left\|u^{0}\right\|_{0, \Omega}+\sum_{i=1}^{2}\left\|\Delta u_{i}^{0}\right\|_{0, \Omega_{i}}+\|u(t)\|_{0, \Omega}\right]
$$

Hence, (2.12) follows from (2.9) and (2.16).
Suppose $S_{1}=\emptyset, q \equiv 0$, and $\alpha \equiv 0$. Let $\left(u^{0}, u^{1}\right) \in H_{S_{1}}^{1}(\Omega) \times L^{2}\left(\Omega, S_{1}\right)=$ $H^{1}(\Omega) \times L_{S_{1}}^{2}(\Omega)$, and set

$$
\begin{equation*}
w^{0}=u^{0}-\frac{1}{m(\Omega)} \int_{\Omega} u^{0}(x) d x, \quad w^{1}=u^{1} \tag{2.17}
\end{equation*}
$$

Then $\left(w^{0}, w^{1}\right) \in \mathcal{H}_{1}$, where $m(\Omega)$ denotes the Lebesgue measure of $\Omega$. Thus, problem (1.1) has a weak solution $w$ for the initial condition ( $w^{0}, w^{1}$ ). Moreover it satisfies

$$
E(w, t) \leqslant E(w, 0), \quad \forall t \geqslant 0
$$

and

$$
\left\|\frac{\partial w}{\partial \nu}\right\|_{L^{2}\left(\Sigma_{2}\right)}^{2}+\left\|w^{\prime}\right\|_{L^{2}\left(\Sigma_{2}\right)}^{2} \leqslant c E(w, 0)
$$

Set

$$
\begin{equation*}
u=w+\frac{1}{m(\Omega)} \int_{\Omega} u^{0}(x) d x \tag{2.18}
\end{equation*}
$$

It is easy to verify that $u$ is a solution of (1.1) with the initial condition $\left(u^{0}, u^{1}\right)$ (note that $q \equiv 0$ and $\alpha \equiv 0$ ). Moreover,

$$
\begin{gathered}
E(u, t)=E(w, t) \leqslant E(w, 0)=E(u, 0), \quad \forall t \geqslant 0 \\
\frac{\partial u}{\partial \nu}=\frac{\partial w}{\partial \nu}, \quad u^{\prime}=w^{\prime} \in L^{2}\left(\Sigma_{2}\right) \\
\left\|\frac{\partial u}{\partial \nu}\right\|_{L^{2}\left(\Sigma_{2}\right)}^{2}+\left\|u^{\prime}\right\|_{L^{2}\left(\Sigma_{2}\right)}^{2} \leqslant c E(u, 0)
\end{gathered}
$$

REMARK 2.2. If $u^{1} \in L^{2}(\Omega)$ rather than $L_{S_{1}}^{2}(\Omega)$, we don't know if problem (1.1) has a solution. Although

$$
\begin{equation*}
u=w+\frac{1}{m(\Omega)} \int_{\Omega} u^{0}(x) d x+\frac{t}{m(\Omega)} \int_{\Omega} u^{1}(x) d x \tag{2.19}
\end{equation*}
$$

satisfies equation (1.1) and the initial condition ( $u^{0}, u^{1}$ ), it doesn't satisfy the boundary condition on $S_{2} \times(0, \infty)$. In (2.19), $w$ is the solution of (1.1) with the initial condition ( $w^{0}, w^{1}$ ) given by

$$
w^{0}=u^{0}-\frac{1}{m(\Omega)} \int_{\Omega} u^{0}(x) d x, \quad w^{1}=u^{1}-\frac{1}{m(\Omega)} \int_{\Omega} u^{1}(x) d x
$$

## 3. Proof of Theorem 1.1

In this section we prove Theorem 1.1. We first generalise an inequality of Lagnese [9, Theorem 2] to the case of transmission. This inequality is the key to proving Theorem 1.1.

In the sequel, all functions are assumed to be real-valued.
THEOREM 3.1. For every $\varepsilon>0$ there exists a constant $c(\varepsilon)$ such that for every $\delta>0$,

$$
\begin{equation*}
\int_{0}^{\infty} \int_{\Omega} e^{-2 \delta t}\left(u-I\left(u^{0}\right)\right)^{2} d x d t \leqslant c(\varepsilon) E(u, 0)+\varepsilon \int_{0}^{\infty} \int_{\Omega} e^{-2 \delta t}\left|u^{\prime}\right|^{2} d x d t \tag{3.1}
\end{equation*}
$$

for every solution of (1.1) with $\left(u^{0}, u^{1}\right) \in H_{S_{1}}^{1}(\Omega) \times L^{2}\left(\Omega, S_{1}\right)$, where

$$
I\left(u_{0}\right)= \begin{cases}\frac{1}{m(\Omega)} \int_{\Omega} u^{0} d x, & \text { if } S_{1}=\emptyset, q \equiv 0, \text { and } \alpha \equiv 0  \tag{3.2}\\ 0, & \text { otherwise }\end{cases}
$$

Theorem 3.1 will be proven below. We show how it can be used to prove Theorem 1.1.

For convenience, we adopt the following notation. For a vector field $l(x)=$ $\left(l_{1}(x), \cdots, l_{n}(x)\right)$ of class $C^{2}(\bar{\Omega})$ the additional subscripts in $l_{i j}$ and $l_{i i j}$ denote derivatives of the vector field $l$, for example, $l_{i j}=\frac{\partial l_{i}}{\partial x_{j}}$.

Proof of Theorem 1.1: Case I: $\alpha_{0}=\max _{x \in S_{2}} \alpha(x)$ is small enough.
We begin with the case where $\left(u^{0}, u^{1}\right) \in \mathcal{H}_{1}$. We may as well assume that $\left(u^{0}, u^{1}\right) \in D(A)$ since the general case $\left(u^{0}, u^{1}\right) \in \mathcal{H}_{1}$ can be handled by a simple limiting process. Then $u$ is a classical solution of (1.1). After a straightforward and tedious calculation, we have

$$
\begin{align*}
& \frac{\partial}{\partial t}\left[t\left(\left|u^{\prime}\right|^{2}+a(x)|\nabla u|^{2}+q|u|^{2}\right)+2 u^{\prime}(l \cdot \nabla u)+\left(l_{i i}-1\right) u u^{\prime}\right] \\
& \quad=\quad \operatorname{div}\left[2 t a(x) u^{\prime} \nabla u+2 a(x)(l \cdot \nabla u) \nabla u+\left|u^{\prime}\right|^{2} l\right. \\
& \left.\quad+a(x)\left(l_{i i}-1\right) u \nabla u-a(x)|\nabla u|^{2} l\right]  \tag{3.3}\\
& \quad+2 a(x)\left(\delta_{i j}-l_{i j}\right) \frac{\partial u}{\partial x_{i}} \frac{\partial u}{\partial x_{j}}-a(x) l_{i i j} u \frac{\partial u}{\partial x_{j}}+\left(2-l_{i i}\right) q|u|^{2}-2 q(l \cdot \nabla u) u
\end{align*}
$$

where $\delta_{i j}$ denote the Kronecker symbol, that is,

$$
\delta_{i j}= \begin{cases}1, & i=j \\ 0, & i \neq j\end{cases}
$$

and summation convention is assumed. Set

$$
\begin{align*}
P(t)=\int_{\Omega}\left[\frac{t}{2}\left(\left|u^{\prime}\right|^{2}+a(x)|\nabla u|^{2}+q|u|^{2}\right)+2 u^{\prime}(l \cdot \nabla u)\right. & \left.+\left(l_{i i}-1\right) u u^{\prime}\right] d x  \tag{3.4}\\
& +\int_{S_{2}} \frac{t}{2} a(x) \alpha(x)|u|^{2} d \Gamma
\end{align*}
$$

Since

$$
\begin{aligned}
& \frac{\partial}{\partial t} \int_{\Omega} t\left(\left|u^{\prime}\right|^{2}+a(x)|\nabla u|^{2}+q|u|^{2}\right) d x \\
& \quad=\int_{\Omega}\left(\left|u^{\prime}\right|^{2}+a(x)|\nabla u|^{2}+q|u|^{2}\right) d x \\
& \quad+2 t \int_{\Omega}\left(u^{\prime} u^{\prime \prime}+a(x) \nabla u \cdot \nabla u^{\prime}+q u u^{\prime}\right) d x \\
& \quad=\int_{\Omega}\left(\left|u^{\prime}\right|^{2}+a(x)|\nabla u|^{2}+q|u|^{2}\right) d x+2 t \int_{S_{2}} a(x) u^{\prime} \frac{\partial u}{\partial \nu} d \Gamma
\end{aligned}
$$

and

$$
\frac{\partial}{\partial t} \int_{S_{2}} \frac{t}{2} a(x) \alpha(x)|u|^{2} d \Gamma=\int_{S_{2}} \frac{1}{2} a(x) \alpha(x)|u|^{2} d \Gamma+\int_{S_{2}} t a(x) \alpha(x) u u^{\prime} d \Gamma
$$

it follows from (3.3) and the divergence theorem that

$$
\begin{aligned}
\frac{\partial P}{\partial t}=- & \frac{1}{2} \frac{\partial}{\partial t} \int_{\Omega} t\left(\left|u^{\prime}\right|^{2}+a(x)|\nabla u|^{2}+q|u|^{2}\right) d x \\
& +\int_{\Omega} 2 a(x)\left(\delta_{i j}-l_{i j}\right) \frac{\partial u}{\partial x_{i}} \frac{\partial u}{\partial x_{j}} d x \quad\left(=I_{2}\right) \\
& +\int_{\Omega}\left[-a(x) l_{i i} u \frac{\partial u}{\partial x_{j}}+\left(2-l_{i i}\right) q|u|^{2}-2 q(l \cdot \nabla u) u\right] d x \quad\left(=I_{3}\right) \\
& +\int_{S_{1}}\left[2 a_{1}\left(l \cdot \nabla u_{1}\right) \frac{\partial u_{1}}{\partial \nu}-a_{1}\left|\nabla u_{1}\right|^{2} l \cdot \nu\right] d \Gamma \quad\left(=I_{4}\right) \\
& +\int_{S_{2}}\left[a_{2} t u_{2}^{\prime}\left(2 \frac{\partial u_{2}}{\partial \nu}+\alpha(x) u_{2}\right)+2 a_{2}\left(l \cdot \nabla u_{2}\right) \frac{\partial u_{2}}{\partial \nu}+\left|u_{2}^{\prime}\right|^{2} l \cdot \nu\right. \\
& \left.+a_{2}\left(l_{i i}-1\right) u_{2} \frac{\partial u_{2}}{\partial \nu}-a_{2}\left|\nabla u_{2}\right|^{2} l \cdot \nu+\frac{1}{2} a_{2} \alpha(x)\left|u_{2}\right|^{2}\right] d \Gamma \\
& +\int_{S_{0}}\left[2 a_{1}\left(l \cdot \nabla u_{1}\right) \frac{\partial u_{1}}{\partial \nu}-2 a_{2}\left(l \cdot \nabla u_{2}\right) \frac{\partial u_{2}}{\partial \nu}\right. \\
& \left.+\left(a_{2}\left|\nabla u_{2}\right|^{2}-a_{1}\left|\nabla u_{1}\right|^{2}\right) l \cdot \nu\right] d \Gamma \quad\left(=I_{6}\right) \\
=- & \frac{1}{2} \int_{\Omega}\left(\left|u^{\prime}\right|^{2}+a(x)|\nabla u|^{2}+q|u|^{2}\right) d x-\frac{1}{2} \int_{S_{2}} a \alpha|u|^{2} d \Gamma \quad\left(=I_{1}\right) \\
& +I_{2}+I_{3}+I_{4} \\
& +\int_{S_{2}}\left[a_{2} t u_{2}^{\prime}\left(\frac{\partial u_{2}}{\partial \nu}+\alpha(x) u_{2}\right)+2 a_{2}\left(l \cdot \nabla u_{2}\right) \frac{\partial u_{2}}{\partial \nu}+\left|u_{2}^{\prime}\right|^{2} l \cdot \nu\right. \\
& \left.+a_{2}\left(l_{i i}-1\right) u_{2} \frac{\partial u_{2}}{\partial \nu}-a_{2}\left|\nabla u_{2}\right|^{2} l \cdot \nu+a_{2} \alpha(x)\left|u_{2}\right|^{2}\right] d \Gamma \quad\left(=I_{5}\right)
\end{aligned}
$$

$$
\begin{aligned}
& \quad+I_{6} \\
& =I_{1}+I_{2}+I_{3}+I_{4}+I_{5}+I_{6} .
\end{aligned}
$$

Since, for any positive constant $c, c l$ still satisfies conditions (i)-(iv) of Theorem 1.1, we may assume that ( $2 \delta_{i j}-l_{i j}-l_{j i}$ ) is negative definite in $\Omega$ by multiplying $l$ by an enough large positive constant. Thus we have $I_{2} \leqslant 0$. From condition (i) of Theorem 1.1 and the fact $\nabla u=\frac{\partial u}{\partial \nu} \nu$ on $S_{1}$ it follows that

$$
\begin{equation*}
I_{4}=\int_{S_{1}} a_{1}\left|\frac{\partial u_{1}}{\partial \nu}\right|^{2} l \cdot \nu d \Gamma \leqslant 0 . \tag{3.5}
\end{equation*}
$$

Concerning $I_{5}$, it follows from the fourth equation of (1.1) and condition (ii) of Theorem 1.1 that

$$
\begin{aligned}
I_{5}= & \int_{S_{2}}\left[-a_{2} \sigma t\left|u_{2}^{\prime}\right|^{2}-2 a_{2}\left(l \cdot \nabla u_{2}\right)\left(\alpha u_{2}+\sigma u_{2}^{\prime}\right)+\left|u_{2}^{\prime}\right|^{2} l \cdot \nu\right. \\
& \left.\quad-a_{2}\left(l_{i i}-1\right) u_{2}\left(\alpha u_{2}+\sigma u_{2}^{\prime}\right)-a_{2}\left|\nabla u_{2}\right|^{2} l \cdot \nu+a_{2} \alpha(x)\left|u_{2}\right|^{2}\right] d \Gamma \\
\leqslant & \int_{S_{2}}\left[-a_{2} \sigma t\left|u_{2}^{\prime}\right|^{2}+\frac{a_{2} \eta}{4}\left|\nabla u_{2}\right|^{2}+c(a, l, \eta) \alpha^{2}\left|u_{2}\right|^{2}\right. \\
& \quad+\frac{a_{2} \eta}{4}\left|\nabla u_{2}\right|^{2}+c(a, l, \eta, \sigma)\left|u_{2}^{\prime}\right|^{2}+\left|u_{2}^{\prime}\right|^{2} l \cdot \nu \\
& \left.\quad-a_{2}\left(l_{i i}-1\right) \alpha\left|u_{2}\right|^{2}+\varepsilon\left|u_{2}\right|^{2}+c(\varepsilon, a, l, \sigma)\left|u_{2}^{\prime}\right|^{2}-a_{2} \eta\left|\nabla u_{2}\right|^{2}+a_{2} \alpha(x)\left|u_{2}\right|^{2}\right] d \Gamma \\
\leqslant & \int_{S_{2}}\left[\left(-a_{2} t \sigma_{0}+l \cdot \nu+c(\varepsilon, a, l, \sigma, \eta)\right)\left|u_{2}^{\prime}\right|^{2}-\frac{a_{2} \eta}{2}\left|\nabla u_{2}\right|^{2}\right] d \Gamma \\
& \quad+\int_{S_{2}}\left(\varepsilon+c(a, l, \eta)\left(\alpha_{0}+\alpha_{0}^{2}\right)\right)\left|u_{2}\right|^{2} d \Gamma \\
= & I_{51}+I_{52} .
\end{aligned}
$$

By the trace theorem, we have

$$
\begin{equation*}
I_{52} \leqslant c_{1}\left(\varepsilon+c(a, l, \eta)\left(\alpha_{0}+\alpha_{0}^{2}\right)\right)\|u\|_{H^{1}\left(\Omega, S_{1}\right)}^{2} \tag{3.6}
\end{equation*}
$$

Concerning $I_{3}$, we have

$$
\begin{equation*}
I_{3} \leqslant \varepsilon \int_{\Omega}|\nabla u|^{2} d x+c(\varepsilon, a, l, q) \int_{\Omega}|u|^{2} d x \tag{3.7}
\end{equation*}
$$

If $\varepsilon$ and $\alpha_{0}$ are small enough, then by (3.6)-(3.7) we have

$$
\begin{aligned}
& I_{3}+I_{52} \leqslant \frac{1}{4} \int_{\Omega}\left(\left|u^{\prime}\right|^{2}+a(x)|\nabla u|^{2}+q|u|^{2}\right) d x+\frac{1}{4} \int_{S_{2}} a \alpha|u|^{2} d \Gamma \\
&+c(\varepsilon, a, l, q) \int_{\Omega}|u|^{2} d x
\end{aligned}
$$

It therefore follows that

$$
\begin{aligned}
I_{1}+I_{3}+I_{52} \leqslant-\frac{1}{4} \int_{\Omega}\left(\left|u^{\prime}\right|^{2}+a(x)|\nabla u|^{2}\right. & \left.+q|u|^{2}\right) d x \\
& -\frac{1}{4} \int_{S_{2}} a \alpha|u|^{2} d \Gamma+c(\varepsilon, a, l, q) \int_{\Omega}|u|^{2} d x
\end{aligned}
$$

if $\varepsilon$ and $\alpha_{0}$ are small enough. Fix $\varepsilon$ and $\alpha_{0}$, then $I_{51}<0$ if $t$ is large enough. We also prove that $I_{6} \leqslant 0$. Since $u_{1}=u_{2}$ on $S_{0}$, we have

$$
\nabla\left(u_{2}-u_{1}\right)=\frac{\partial\left(u_{2}-u_{1}\right)}{\partial \nu} \nu, \quad \text { on } S_{0}
$$

Then,

$$
\begin{aligned}
\left|\nabla u_{2}\right|^{2} & =\left|\nabla u_{1}\right|^{2}+2\left(\frac{\partial u_{2}}{\partial \nu}-\frac{\partial u_{1}}{\partial \nu}\right) \frac{\partial u_{1}}{\partial \nu}+\left(\frac{\partial u_{2}}{\partial \nu}-\frac{\partial u_{1}}{\partial \nu}\right)^{2} \\
& =\left|\nabla u_{1}\right|^{2}+\left(\frac{\partial u_{2}}{\partial \nu}\right)^{2}-\left(\frac{\partial u_{1}}{\partial \nu}\right)^{2}
\end{aligned}
$$

So,

$$
\begin{aligned}
& 2 a_{1}\left(l \cdot \nabla u_{1}\right) \frac{\partial u_{1}}{\partial \nu}-2 a_{2}\left(l \cdot \nabla u_{2}\right) \frac{\partial u_{2}}{\partial \nu}+\left(a_{2}\left|\nabla u_{2}\right|^{2}-a_{1}\left|\nabla u_{1}\right|^{2}\right) l \cdot \nu \\
&= 2 a_{1}\left(l \cdot \nabla u_{1}\right) \frac{\partial u_{1}}{\partial \nu}-2 a_{2}\left[l \cdot \nabla u_{1}+\left(\frac{\partial u_{2}}{\partial \nu}-\frac{\partial u_{1}}{\partial \nu}\right) l \cdot \nu\right] \frac{\partial u_{2}}{\partial \nu} \\
&+\left[a_{2}\left(\left|\nabla u_{1}\right|^{2}+\left(\frac{\partial u_{2}}{\partial \nu}\right)^{2}-\left(\frac{\partial u_{1}}{\partial \nu}\right)^{2}\right)-a_{1}\left|\nabla u_{1}\right|^{2}\right] l \cdot \nu \\
&= 2 a_{1}\left(l \cdot \nabla u_{1}\right) \frac{\partial u_{1}}{\partial \nu}-2 a_{1}\left[l \cdot \nabla u_{1}+\left(\frac{a_{1}}{a_{2}} \frac{\partial u_{1}}{\partial \nu}-\frac{\partial u_{1}}{\partial \nu}\right) l \cdot \nu\right] \frac{\partial u_{1}}{\partial \nu} \\
&+\left[a_{2}\left(\left|\nabla u_{1}\right|^{2}+\frac{a_{1}^{2}}{a_{2}^{2}}\left(\frac{\partial u_{1}}{\partial \nu}\right)^{2}-\left(\frac{\partial u_{1}}{\partial \nu}\right)^{2}\right)-a_{1}\left|\nabla u_{1}\right|^{2}\right] l \cdot \nu \\
&=\left(a_{2}-a_{1}\right)\left|\nabla u_{1}\right|^{2} l \cdot \nu-\frac{\left(a_{2}-a_{1}\right)^{2}}{a_{2}}\left(\frac{\partial u_{1}}{\partial \nu}\right)^{2} l \cdot \nu
\end{aligned}
$$

This show that $I_{6} \leqslant 0$ because of (iii) of Theorem 1.1. It therefore follows that

$$
\begin{align*}
& \frac{d P}{d t} \leqslant-\frac{1}{4} \int_{\Omega}\left(\left|u^{\prime}\right|^{2}+a(x)|\nabla u|^{2}+q|u|^{2}\right) d x  \tag{3.8}\\
&-\frac{1}{4} \int_{S_{2}} a \alpha|u|^{2} d \Gamma+c(\varepsilon) \int_{\Omega}|u|^{2} d x, \quad t \geqslant T_{1}
\end{align*}
$$

if $T_{1}$ is large enough.

On the other hand, there exists $T_{2}$ sufficiently large such that

$$
\begin{equation*}
0 \leqslant P(t) \leqslant c(t+1) E(u, t), \quad t \geqslant T_{2} \tag{3.9}
\end{equation*}
$$

where $c$ is a constant independent of $t, u$.
Let $\delta>0$ be fixed. Set $T=\max \left\{T_{1}, T_{2}\right\}$. By multiplying (3.8) by $e^{-2 \delta t}$ and integrating from $T$ to $+\infty$, we get

$$
\begin{equation*}
2 \delta \int_{T}^{\infty} e^{-2 \delta t} P(t) d t+\frac{1}{8} \int_{T}^{\infty} e^{-2 \delta t} E(u, t) d t \leqslant c_{1} E(u, 0)+c_{2} \int_{T}^{\infty} \int_{\Omega} e^{-2 \delta t} u^{2} d x d t \tag{3.10}
\end{equation*}
$$

where $c_{1}, c_{2}$ are independent of $\delta$. It therefore follows from (3.9) and (3.10) that

$$
\int_{0}^{\infty} e^{-2 \delta t} E(u, t) d t \leqslant c_{1} E(u, 0)+c_{2} \int_{0}^{\infty} \int_{\Omega} e^{-2 \delta t} u^{2} d x d t
$$

Applying Theorem 3.1, we conclude that

$$
\int_{0}^{\infty} e^{-2 \delta t} E(u, t) d t \leqslant c E(u, 0)
$$

Letting $\delta \rightarrow 0$, we obtain

$$
\int_{0}^{\infty} E(u, t) d t \leqslant c E(u, 0) .
$$

By [16, Theorem 4.1, p.116], there are positive constants $M, \tau$ such that

$$
\begin{equation*}
E(u, t) \leqslant M e^{-\tau t} E(u, 0), \quad t \geqslant 0 . \tag{3.11}
\end{equation*}
$$

If $S_{1}=\emptyset, q \equiv 0, \alpha \equiv 0$, and $\left(u^{0}, u^{1}\right) \in H_{S_{1}}^{1}(\Omega) \times L_{S_{1}}^{2}(\Omega)$ rather than $\mathcal{H}_{1}$, then we take $\left(w^{0}, w^{1}\right)$ as in (2.17). Let $w$ be the solution of (1.1) with the initial condition ( $w^{0}, w^{1}$ ), then

$$
u=w+\frac{1}{m(\Omega)} \int_{\Omega} u^{0}(x) d x
$$

is the solution of (1.1) with the initial condition $\left(u^{0}, u^{1}\right)$ and (3.11) holds for $w$. Therefore,

$$
E(u, t)=E(w, t) \leqslant M e^{-\tau t} E(w, 0)=M e^{-\tau t} E(u, 0)
$$

We shall use control-theoretic method given in [2] and [13] to prove Theorem 1.1 in the case that $\alpha_{0}$ is arbitrary. Therefore we now employ Russell's "controllability via stabilisability " principle (see [17]) to solve the following exact controllability problem:

For ( $y_{i}^{0}, y_{i}^{1}$ ) in a suitable Hilbert space and $T$ large enough, find a control function $\phi(x, t)$ such that the solution of

$$
\begin{cases}y_{i}^{\prime \prime}-a_{i} \Delta y_{i}+q y_{i}=0 & \text { in } Q_{i}  \tag{3.12}\\ y_{i}(0)=y_{i}^{0}, \quad y_{i}^{\prime}(0)=y_{i}^{1} & \text { in } \Omega_{i}, i=1,2 \\ y_{1}=0 & \text { on } \Sigma_{1} \\ \frac{\partial y_{2}}{\partial \nu}+\alpha y_{2}=\phi & \text { on } \Sigma_{2} \\ y_{1}=y_{2}, \quad a_{1} \frac{\partial y_{1}}{\partial \nu}=a_{2} \frac{\partial y_{2}}{\partial \nu} & \text { on } \Sigma_{0}\end{cases}
$$

satisfies

$$
\begin{equation*}
y_{i}(x, T ; \phi)=y_{i}^{\prime}(x, T ; \phi)=0 \quad \text { in } \quad \Omega_{i}, i=1,2 \tag{3.13}
\end{equation*}
$$

Because the problem is linear, this is equivalent to steering any initial state to any terminal state. This controllability problem was discussed in [1] in the general case of second order hyperbolic equations, but the coefficients of the equations are required to be smooth enough. Thus, problem (3.12) (3.13) here is not covered by [1]. With the help of the Hilbert Uniqueness Method, problem (3.12) (3.13) was also considered in [15] but with $\alpha=0$.

Theorem 3.2. Suppose all assumptions of Theorem 1.1 are satisfied. Suppose $\alpha_{0}$ is small enough and $T$ is given large enough. Then for any $\left(y^{0}, y^{1}\right) \in \mathcal{H}_{1}$, there exists a boundary control function

$$
\phi(x, t) \in L^{2}\left(\Sigma_{2}\right)
$$

such that the solution of (3.12) satisfies (3.13). Moreover, there exist positive constants $c_{1}(T), c_{2}(T)$ such that

$$
\begin{equation*}
c_{1} E(y, 0) \leqslant\|\phi\|_{L^{2}\left(\Sigma_{2}\right)}^{2} \leqslant c_{2} E(y, 0) \tag{3.14}
\end{equation*}
$$

Proof: We first consider the problem:

$$
\begin{cases}u_{i}^{\prime \prime}-a_{i} \Delta u_{i}+q u_{i}=0 & \text { in } \Omega_{i} \times(0, \infty)  \tag{3.15}\\ u_{i}(0)=u_{i}^{0}, \quad u_{i}^{\prime}(0)=u_{i}^{1} & \text { in } \Omega_{i}, i=1,2 \\ u_{1}=0 & \text { on } S_{1} \times(0, \infty) \\ \frac{\partial u_{2}}{\partial \nu}+\alpha u_{2}+\sigma u_{2}^{\prime}=0 & \text { on } S_{2} \times(0, \infty) \\ u_{1}=u_{2}, \quad a_{1} \frac{\partial u_{1}}{\partial \nu}=a_{2} \frac{\partial u_{2}}{\partial \nu} & \text { on } S_{0} \times(0, \infty)\end{cases}
$$

which has a unique weak solution with

$$
u(t) \in C\left([0, \infty) ; H^{1}\left(\Omega, S_{1}\right)\right) \cap C^{1}\left([0, \infty) ; L^{2}\left(\Omega, S_{1}\right)\right)
$$

and

$$
u^{\prime} \in L^{2}\left(\Sigma_{2}\right)
$$

for any $\left(u^{0}, u^{1}\right) \in \mathcal{H}_{1}$, thanks to Theorem 2.1.
Using the solution $u$ of (3.15), we then consider the backwards problem:

$$
\begin{cases}w_{i}^{\prime \prime}-a_{i} \Delta w_{i}+q w_{i}=0 & \text { in } \Omega_{i} \times(0, \infty) \\ w_{i}(T)=u_{i}(T), \quad w_{i}^{\prime}(T)=u_{i}^{\prime}(T) & \text { in } \Omega_{i}, i=1,2 \\ w_{1}=0 & \text { on } S_{1} \times(0, \infty) \\ \frac{\partial w_{2}}{\partial \nu}+\alpha w_{2}-\sigma w_{2}^{\prime}=0 & \text { on } S_{2} \times(0, \infty) \\ w_{1}=w_{2}, \quad a_{1} \frac{\partial w_{1}}{\partial \nu}=a_{2} \frac{\partial w_{2}}{\partial \nu} & \text { on } S_{0} \times(0, \infty)\end{cases}
$$

which has a unique weak solution with

$$
w \in C\left([0, T] ; H^{1}\left(\Omega, S_{1}\right)\right) \cap C^{1}\left([0, T] ; L^{2}\left(\Omega, S_{1}\right)\right)
$$

and

$$
w^{\prime} \in L^{2}\left(\Sigma_{2}\right)
$$

since $\left(u(x, T), u^{\prime}(x, T)\right) \in \mathcal{H}_{1}$.
Set

$$
y=u-w
$$

and

$$
\phi=-\sigma(x)\left(w^{\prime}+u^{\prime}\right) \in L^{2}\left(\Sigma_{2}\right) .
$$

Then $y$ satisfies

$$
\begin{cases}y_{i}^{\prime \prime}-a_{i} \Delta y_{i}+q y_{i}=0 & \text { in } \Omega_{i} \times(0, \infty) \\ y_{i}(0)=u_{i}^{0}-w_{i}(0), \quad y_{i}^{\prime}(0)=u_{i}^{1}-w_{i}^{\prime}(0) & \text { in } \Omega_{i} \\ y_{i}(T)=0, \quad y_{i}^{\prime}(T)=0 & \text { in } \Omega_{i}, i=1,2 \\ y_{1}=0 & \text { on } S_{1} \times(0, \infty) \\ \frac{\partial y_{2}}{\partial \nu}+\alpha y_{2}=\phi(x) & \text { on } S_{2} \times(0, \infty) \\ y_{1}=y_{2}, \quad a_{1} \frac{\partial y_{1}}{\partial \nu}=a_{2} \frac{\partial y_{2}}{\partial \nu} & \text { on } S_{0} \times(0, \infty)\end{cases}
$$

We define an operator $\Lambda$ by

$$
\Lambda\left(u^{0}, u^{1}\right)=\left(w(x, 0), w^{\prime}(x, 0)\right)
$$

Then it is clear that $\Lambda$ is a linear operator from $\mathcal{H}_{1}$ into $\mathcal{H}_{1}$. Moreover, by Theorem 1.1 (in the case where $\alpha_{0}$ is small enough) we have

$$
\begin{aligned}
\left\|\Lambda\left(u^{0}, u^{1}\right)\right\|_{\mathcal{H}_{1}}^{2} & =E(w, 0) \\
& \leqslant M e^{-\tau T} E(w, T) \\
& \leqslant M^{2} e^{-2 \tau T}\left\|\left(u^{0}, u^{1}\right)\right\|_{\mathcal{H}_{1}}
\end{aligned}
$$

Therefore,

$$
\|\Lambda\| \leqslant M e^{-\tau T}
$$

Taking $T$ large enough so that $M e^{-\tau T}<1$, then $I-\Lambda$ is an isomorphism from $\mathcal{H}_{1}$ onto $\mathcal{H}_{1}$. Thus, for any $\left(y^{0}, y^{1}\right) \in \mathcal{H}_{1}$, there exists a unique $\left(u^{0}, u^{1}\right) \in \mathcal{H}_{1}$ such that

$$
\begin{equation*}
\left(y^{0}, y^{1}\right)=\left(u^{0}, u^{1}\right)-\Lambda\left(u^{0}, u^{1}\right)=\left(u^{0}, u^{1}\right)-\left(w(x, 0), w^{\prime}(x, 0)\right) \tag{3.17}
\end{equation*}
$$

Consequently, we have constructed a control function $\phi=-\sigma(x)\left(w^{\prime}+u^{\prime}\right)$ solving the exact controllability problem (3.12)-(3.13).

On the other hand, by multiplying the first equation of (3.15) by $u^{\prime}$ and integrating over $Q$, we obtain

$$
\int_{\Sigma_{2}} a(x) \sigma(x)\left|u^{\prime}\right|^{2} d \Sigma=E(u, 0)-E(u, T)
$$

Likewise, we have

$$
\int_{\Sigma_{2}} a(x) \sigma(x)\left|w^{\prime}\right|^{2} d \Sigma=E(w, T)-E(w, 0)
$$

It therefore follows from Theorem 1.1 that there exist positive constants $c_{1}, c_{2}$ such that

$$
c_{1}\left(1-M e^{-\tau T}\right)^{1 / 2} E^{1 / 2}(u, 0) \leqslant\left\|u^{\prime}\right\|_{L^{2}\left(\Sigma_{2}\right)} \leq_{2} E^{1 / 2}(u, 0)
$$

and

$$
c_{1}\left(1-M e^{-\tau T}\right)^{1 / 2} E^{1 / 2}(w, T) \leqslant\left\|w^{\prime}\right\|_{L^{2}\left(\Sigma_{2}\right)} \leq_{2} E^{1 / 2}(w, T)
$$

Noting $E(u, T)=E(w, T)$, we deduce from the triangle inequality and Theorem 1.1 that

$$
c_{1}\left[1-M e^{-\tau T}\right]^{1 / 2}\left[1-\left(M e^{-\tau T}\right)^{1 / 2}\right] E^{1 / 2}(u, 0) \leqslant\left\|u^{\prime}+w^{\prime}\right\|_{L^{2}\left(\Sigma_{2}\right)}
$$

$$
\begin{equation*}
\leqslant c_{2}\left[1+\left(M e^{-\tau T}\right)^{1 / 2}\right] E^{1 / 2}(u, 0) \tag{3.18}
\end{equation*}
$$

Since $I-\Lambda$ is an isomorphism, (3.14) follows from (3.17) and (3.18).
Proof of Theorem 1.1: Case II: $\alpha_{0}$ is arbitrary.
Let $\varepsilon>0$ be small enough and $T$ large enough. It then follows from Theorem 3.2 that there exists a control $\phi$ such that

$$
\begin{cases}y_{i}^{\prime \prime}-a_{i} \Delta y_{i}+q y_{i}=0 & \text { in } Q_{i}  \tag{3.19}\\ y_{i}(0)=0, \quad y_{i}^{\prime}(0)=0 & \text { in } \Omega_{i} \\ y_{i}(T)=u_{i}(T), \quad y_{i}^{\prime}(T)=u_{i}^{\prime}(T) & \text { in } \Omega_{i}, i=1,2 \\ y_{1}=0 & \text { on } \Sigma_{1}, \\ \frac{\partial y_{2}}{\partial \nu}+\varepsilon y_{2}=\phi(x) & \text { on } \Sigma_{2}, \\ y_{1}=y_{2}, \quad a_{1} \frac{\partial y_{1}}{\partial \nu}=a_{2} \frac{\partial y_{2}}{\partial \nu} & \text { on } \Sigma_{0}\end{cases}
$$

According to the proof of Theorem 3.2, $y$ and $\phi$ can be written as

$$
\begin{equation*}
y=v-w, \quad \phi=\sigma\left(v^{\prime}+w^{\prime}\right) \tag{3.20}
\end{equation*}
$$

where $v$ and $w$ are respectively the solutions of

$$
\begin{align*}
& \begin{cases}v_{i}^{\prime \prime}-a_{i} \Delta v_{i}+q v_{i}=0 & \text { in } Q_{i}, \\
v_{i}(T)=v_{i}^{0}, \quad v_{i}^{\prime}(T)=v_{i}^{1} & \text { in } \Omega_{i}, i=1,2 \\
v_{1}=0 & \text { on } \Sigma_{1}, \\
\frac{\partial v_{2}}{\partial \nu}+\varepsilon v_{2}-\sigma v_{2}^{\prime}=0 & \text { on } \Sigma_{2}, \\
v_{1}=v_{2}, \quad a_{1} \frac{\partial v_{1}}{\partial \nu}=a_{2} \frac{\partial v_{2}}{\partial \nu} & \text { on } \Sigma_{0},\end{cases}  \tag{3.21}\\
& \begin{cases}w_{i}^{\prime \prime}-a_{i} \Delta w_{i}+q w_{i}=0 & \text { in } Q_{i}, \\
w_{i}(0)=v_{i}(0), \quad w_{i}^{\prime}(0)=v_{i}^{\prime}(0) & \text { in } \Omega_{i}, i=1,2, \\
w_{1}=0 & \text { on } \Sigma_{1}, \\
\frac{\partial w_{2}}{\partial \nu}+\varepsilon w_{2}+\sigma w_{2}^{\prime}=0 & \text { on } \Sigma_{2}, \\
w_{1}=w_{2}, \quad a_{1} \frac{\partial w_{1}}{\partial \nu}=a_{2} \frac{\partial w_{2}}{\partial \nu} & \text { on } \Sigma_{0} .\end{cases} \tag{3.22}
\end{align*}
$$

In (3.21), $\left(v^{0}, v^{1}\right)$ are chosen to be such that

$$
\begin{equation*}
\left(v^{0}, v^{1}\right)-\left(w(T), w^{\prime}(T)\right)=\left(v^{0}, v^{1}\right)-\Lambda\left(v^{0}, v^{1}\right)=\left(u(T), u^{\prime}(T)\right) \tag{3.23}
\end{equation*}
$$

Integrating by parts, we obtain

$$
\begin{align*}
0 & =\int_{Q}\left[y^{\prime}\left(u^{\prime \prime}-a \Delta u+q u\right)+u^{\prime}\left(y^{\prime \prime}-a \Delta y+q y\right)\right] d x d t  \tag{3.24}\\
& =\int_{Q} \frac{\partial}{\partial t}\left(u^{\prime} y^{\prime}+a \nabla u \cdot \nabla y+q u y\right) d x d t-\int_{\Sigma_{2}} a\left(y^{\prime} \frac{\partial u}{\partial \nu}+u^{\prime} \frac{\partial y}{\partial \nu}\right) d \Sigma \\
& =\int_{\Omega}\left(\left|u^{\prime}(T)\right|^{2}+a|\nabla u(T)|^{2}+q|u(T)|^{2}\right) d x+\int_{\Sigma_{2}} a\left[y^{\prime}\left(\alpha u+\sigma u^{\prime}\right)+u^{\prime}(\varepsilon y-\phi)\right] d \Sigma \\
& =2 E(u, T)+\int_{\Sigma_{2}} a u^{\prime}\left(\sigma y^{\prime}-\alpha y+\varepsilon y-\phi\right) d \Sigma
\end{align*}
$$

The following constants $c=c(T)$ denote various positive constans depending on $T$. By the trace theorem and (3.20) and (3.23), we have

$$
\begin{align*}
\|y\|_{L^{2}\left(\Sigma_{2}\right)}^{2} & \leqslant c \int_{0}^{T}\|y(t)\|_{H^{1}\left(\Omega, S_{1}\right)}^{2} d t \\
& \leqslant c \int_{0}^{T}\left(\|v(t)\|_{H^{1}\left(\Omega, S_{1}\right)}^{2}+\|w(t)\|_{H^{1}\left(\Omega, S_{1}\right)}^{2}\right) d t  \tag{3.25}\\
& \leqslant c(E(v, T)+E(w, 0)) \\
& \leqslant c E(v, T) \\
& \leqslant c E(u, T)
\end{align*}
$$

Moreover, since

$$
\begin{equation*}
\int_{\Sigma_{2}} a \sigma\left|v^{\prime}\right|^{2} d \Sigma=E(v, T)-E(v, 0) \tag{3.26}
\end{equation*}
$$

and

$$
\begin{equation*}
\int_{\Sigma_{2}} a \sigma\left|w^{\prime}\right|^{2} d \Sigma=E(w, 0)-E(w, T) \tag{3.27}
\end{equation*}
$$

it follows from (3.20) and (3.23) that

$$
\begin{equation*}
\left\|y^{\prime}\right\|_{L^{2}\left(\Sigma_{2}\right)}^{2} \leqslant c E(u, T) . \tag{3.28}
\end{equation*}
$$

By (3.14) we have

$$
\begin{equation*}
\|\phi\|_{L^{2}\left(\Sigma_{2}\right)}^{2} \leqslant c E(u, T) \tag{3.29}
\end{equation*}
$$

By the Cauchy-Schwartz inequality we deduce from (3.24) that

$$
\begin{equation*}
(E(u, T))^{2} \leqslant c \int_{\Sigma_{2}}\left|u^{\prime}\right|^{2} d \Sigma \int_{\Sigma_{2}}\left|\sigma y^{\prime}-\alpha y+\varepsilon y-\phi\right|^{2} d \Sigma \tag{3.30}
\end{equation*}
$$

which, combining (3.25), (3.28), and (3.29), yields

$$
\begin{equation*}
\int_{\Sigma_{2}}\left|u^{\prime}\right|^{2} d \Sigma \geqslant c E(u, T) \tag{3.31}
\end{equation*}
$$

On the other hand, we have

$$
\begin{equation*}
\int_{\Sigma_{2}} a \sigma\left|u^{\prime}\right|^{2} d \Sigma=E(u, 0)-E(u, T) . \tag{3.32}
\end{equation*}
$$

We then conclude

$$
E(u, 0)-E(u, T) \geqslant c E(u, T)
$$

so that

$$
E(u, T) \leqslant \frac{1}{1+c} E(u, 0) .
$$

By repeating the above reasoning, we get

$$
\begin{aligned}
E(u,(k+1) T) & \leqslant \frac{1}{1+c} E(u, k T) \\
& \leqslant \frac{1}{(1+c)^{k+1}} E(u, 0), \quad k=0,1,2, \cdots
\end{aligned}
$$

This implies (1.6) with

$$
M=1+c, \quad \tau=\frac{1}{T} \ln (1+c)
$$

The proof of Theorem 1.1 is complete.
Remark 3.3. From Theorem 1.1 and the proof of Theorem 3.2 we can conclude that Theorem 3.2 still holds true for arbitrary $\alpha_{0}$.

If $S_{1}=\emptyset, q \equiv 0$, and $\alpha \equiv 0$, then $\mathcal{H}_{1}$ in Theorem 3.2 cannot be replaced by $H_{S_{1}}^{1}(\Omega) \times L^{2}\left(\Omega, S_{1}\right)$ because $(E(u, 0))^{1 / 2}$ is no longer a norm on $H^{1}(\Omega) \times L^{2}\left(\Omega, S_{1}\right)$. Nevertheless, (3.16) still holds on the quotient space $\left(H^{1}(\Omega) \times L^{2}\left(\Omega, S_{1}\right)\right) / N$, where $N=\{(c, 0): c \in \mathbb{R}\}$. Therefore, Theorem 3.2 still holds when $\mathcal{H}_{1}$ is replaced by $\left(H^{1}(\Omega) \times L^{2}\left(\Omega, S_{1}\right)\right) / N$. Because the zero element in $\left(H^{1}(\Omega) \times L^{2}\left(\Omega, S_{1}\right)\right) / N$ is $N$, we can only drive any initial state $\left(y^{0}, y^{1}\right) \in H^{1}(\Omega) \times L^{2}\left(\Omega, S_{1}\right)$ to a constant function $(c, 0)$. In fact, we can explain this in the following way. For any $\left(y^{0}, y^{1}\right) \in$ $H^{1}(\Omega) \times L^{2}(\Omega)$, set

$$
w^{0}=y^{0}-\frac{1}{m(\Omega)} \int_{\Omega} y^{0}(x) d x, \quad w^{1}=y^{1}-\frac{1}{m(\Omega)} \int_{\Omega} y^{1}(x) d x
$$

Then $\left(w^{0}, w^{1}\right) \in \mathcal{H}_{1}$. By Theorem 3.2, there exists a control function $\phi(x, t)$ such that the solution $w$ of (3.12) with the initial state $\left(w^{0}, w^{1}\right)$ satisfies (3.13). It is easy to check that

$$
y=w+\frac{1}{m(\Omega)} \int_{\Omega} y^{0}(x) d x+\frac{t}{m(\Omega)} \int_{\Omega} y^{1}(x) d x
$$

is the solution of (3.12) with the initial state $\left(y^{0}, y^{1}\right)$, but

$$
\begin{aligned}
y(x, T ; \phi) & =\frac{1}{m(\Omega)} \int_{\Omega} y^{0}(x) d x+\frac{T}{m(\Omega)} \int_{\Omega} y^{1}(x) d x, \quad \text { (a constant) } \\
y^{\prime}(x, T ; \phi) & =\frac{1}{m(\Omega)} \int_{\Omega} y^{1}(x) d x, \quad \text { (a constant). }
\end{aligned}
$$

At last, we want to prove Theorem 3.1. For this, we need the following lemma.
Lemma 3.4. Let $\phi \in H^{1 / 2}(\Gamma)$. Then there exists $u \in H^{2}\left(\Omega_{1}, \Omega_{2}\right)$ such that

$$
\begin{equation*}
u=0, \quad \frac{\partial u}{\partial \nu}=\phi, \quad \text { on } \Gamma \tag{3.33}
\end{equation*}
$$

and

$$
\begin{equation*}
\left\|u_{1}\right\|_{2, \Omega_{1}}+\left\|u_{2}\right\|_{2, \Omega_{2}} \leqslant c\|\phi\|_{H^{1 / 2}(\Gamma)} \tag{3.34}
\end{equation*}
$$

where $c$ is a positive constant and

$$
H^{2}\left(\Omega_{1}, \Omega_{2}\right)=\left\{u: u_{i}=\left.u\right|_{\Omega_{i}} \in H^{2}\left(\Omega_{i}\right), i=1,2 ; u_{1}=u_{2}, a_{1} \frac{\partial u_{1}}{\partial \nu}=a_{2} \frac{\partial u_{2}}{\partial \nu} \text { on } S_{0}\right\}
$$

Proof: By the trace theorem, it follows that there exists $w \in H^{2}(\Omega)$ such that

$$
w=0, \quad \frac{\partial w}{\partial \nu}=\phi, \quad \text { on } \Gamma
$$

and

$$
\|w\|_{2, \Omega} \leqslant c\|\phi\|_{H^{1 / 2}(\Gamma)}
$$

Since $w \in H^{2}\left(\Omega_{2}\right)$, again by the trace theorem, we have

$$
w \in H^{3 / 2}\left(S_{0}\right), \quad \frac{\partial w}{\partial \nu} \in H^{1 / 2}\left(S_{0}\right)
$$

Also by the trace theorem, it follows that there exists $v \in H^{2}\left(\Omega_{1}\right)$ such that

$$
\begin{array}{lll}
v=0, & \frac{\partial v}{\partial \nu}=\phi & \text { on } \Gamma \cap \Gamma_{1} \\
v=w, & \frac{\partial v}{\partial \nu}=\frac{a_{2}}{a_{1}} \frac{\partial w}{\partial \nu} & \text { on } S_{0}
\end{array}
$$

and

$$
\begin{aligned}
\|v\|_{2, \Omega_{1}} & \leqslant c\left[\|\phi\|_{H^{1 / 2}\left(\Gamma \cap \Gamma_{1}\right)}+\|w\|_{H^{3 / 2}\left(S_{0}\right)}+\left\|\frac{\partial w}{\partial \nu}\right\|_{H^{1 / 2}\left(S_{0}\right)}\right] \\
& \leqslant c\left[\|\phi\|_{H^{1 / 2}\left(\Gamma \cap \Gamma_{1}\right)}+\|w\|_{2, \Omega_{2}}\right] \\
& \leqslant c\left[\|\phi\|_{H^{1 / 2}\left(\Gamma \cap \Gamma_{1}\right)}+\|w\|_{2, \Omega}\right] \\
& \leqslant c\|\phi\|_{H^{1 / 2}(\Gamma)} .
\end{aligned}
$$

Then, $u$ defined by

$$
u= \begin{cases}v, & x \in \Omega_{1} \\ w, & x \in \Omega_{2}\end{cases}
$$

belongs to $H^{2}\left(\Omega_{1}, \Omega_{2}\right)$ and satisfies (3.33) and (3.34).
We also need the following unique continuation theorem for elliptic operators given in [7].

Let $A(x)=\left(a_{i j}(x)\right)$ be a real symmetric matrix-valued function on $\Omega$ satisfying the assumptions:
(i) there exists a $\rho \in(0,1)$ such that, for every $x \in \Omega$ and $\xi \in \mathbb{R}^{n}$,

$$
\begin{equation*}
\rho|\xi|^{2} \leqslant \sum_{i, j=1}^{n} a_{i j}(x) \xi_{i} \xi_{j} \leqslant \rho^{-1}|\xi|^{2} \tag{3.35}
\end{equation*}
$$

(ii) there exists a $K>0$ such that, for every $x, y \in \Omega$,

$$
\begin{equation*}
\left|a_{i j}(x)-a_{i j}(y)\right| \leqslant K|x-y|, \quad i, j=1,2, \cdots, n \tag{3.36}
\end{equation*}
$$

Let the potential $V$ satisfy the assumption: for every $x_{0} \in \Omega$ there exist $r_{0}>0$ and two constants $C_{1}, C_{2}>0$ such that if $V(x)=V^{+}(x)-V^{-}(x)$, then

$$
\begin{align*}
& 0 \leqslant V^{+}(x) \leqslant \frac{C_{1}}{\left|x-x_{0}\right|^{2}}  \tag{3.37}\\
& 0 \leqslant V^{-}(x) \leqslant \frac{C_{2}}{\left|x-x_{0}\right|^{2}} \tag{3.38}
\end{align*}
$$

for any $x \in B_{r_{0}}\left(x_{0}\right) \cap \Omega$, where $B_{r_{0}}\left(x_{0}\right)=\left\{x \in \mathbb{R}^{n}:\left|x-x_{0}\right|<r_{0}\right\}$.
Theorem 3.5. [7, Corollary 1.1] Assume that $A(x)$ and $V(x)$ satisfy the above assumptions (3.35)-(3.38). Then the operator $L=-\operatorname{div}(A(x) \nabla)+V(x)$ has the unique
continuation property in $\Omega$, that is, the only $H_{l o c}^{1}(\Omega)$ solution of $L u=0$ which can vanish in an open subset of $\Omega$ is $u \equiv 0$.

Proof of Theorem 3.1: Case I: $S_{1} \neq \emptyset$, or $q \neq 0$, or $\alpha \neq 0$. The proof is the same as the one of Theorem 2 of [ 9$]$ except for the following two points:
(i) In the proof of Theorem 2 of [9], Lagnese used the analyticity of the solution of

$$
\Delta W-\omega^{2} W=0 \quad \text { in } \quad \Omega, \quad \omega \text { real }
$$

to conclude that $W=0$ in $\Omega$ if $W=0$ in an open subset of $\Omega$. For the present case of transmission, we use Theorem 3.5 since we now can no longer appeal to the analyticity of solutions.
(ii) The usual trace theorem used in the proof of Theorem 2 of [9] is replaced by Lemma 3.4.
Case II: $\quad S_{1}=\emptyset, q \equiv 0, \alpha \equiv 0$. If $\left(u^{0}, u^{1}\right) \in H_{S_{1}}^{1}(\Omega) \times L^{2}\left(\Omega, S_{1}\right)$ rather than $\mathcal{H}_{1}$, then we take $\left(w^{0}, w^{1}\right)$ as in (2.17). Let $w$ be the solution of (1.1) with the initial condition ( $w^{0}, w^{1}$ ). Then

$$
u=w+\frac{1}{m(\Omega)} \int_{\Omega} u^{0}(x) d x
$$

is the solution of (1.1) with the initial condition $\left(u^{0}, u^{1}\right)$ and (3.1) holds for $w$. Therefore,

$$
\begin{aligned}
\int_{0}^{\infty} \int_{\Omega} e^{-2 \delta t}\left(u-I\left(u^{0}\right)\right)^{2} d x d t & =\int_{0}^{\infty} \int_{\Omega} e^{-2 \delta t} w^{2} d x d t \\
& \leqslant c(\varepsilon) E(w, 0)+\varepsilon \int_{0}^{\infty} \int_{\Omega} e^{-2 \delta t}\left|w^{\prime}\right|^{2} d x d t \\
& =c(\varepsilon) E(u, 0)+\varepsilon \int_{0}^{\infty} \int_{\Omega} e^{-2 \delta t}\left|u^{\prime}\right|^{2} d x d t
\end{aligned}
$$

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