

Scattering Theory for the Elastic Wave Equation

By

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§ 0. Introduction

Let us set

$$Lu(x) = \sum_{i,j=1}^n \partial_i(a_{ij}(x)\partial_j u(x)), \quad \partial_i = \partial/\partial x_i, \quad u = {}^t(u_1, \dots, u_n),$$

where $a_{ij}(x)$ are $n \times n$ matrices of real-valued functions $\in \mathcal{B}^\infty(\mathbf{R}^n)$ ($= \{f \in C^\infty(\mathbf{R}^n) \mid \sup_{x \in \mathbf{R}^n} |\partial_x^\alpha f(x)| < \infty \text{ for any multi-index } \alpha\}$). Let \mathcal{Q} be an exterior domain in \mathbf{R}^n whose boundary $\partial\mathcal{Q}$ is a compact and C^∞ hypersurface contained in $\{x \in \mathbf{R}^n \mid |x| < r_0\}$. We consider the elastic wave equation:

$$(0.1) \quad \begin{cases} (\partial_t^2 - L)u(t, x) = 0 \text{ in } \mathbf{R} \times \mathcal{Q} \quad (\partial_t = \partial/\partial t), \quad Bu(t, x') = 0 \text{ on } \mathbf{R} \times \partial\mathcal{Q}, \\ u(0, x) = f_1(x) \text{ in } \mathcal{Q}, \quad \partial_t u(0, x) = f_2(x) \text{ in } \mathcal{Q}. \end{cases}$$

Here the boundary operator B is of the form

$$(0.2) \quad Bu = u|_{\partial\mathcal{Q}}, \quad \text{or}$$

$$(0.3) \quad Bu = \sum_{i,j=1}^n \nu_i(x) a_{ji}(x) \partial_j u|_{\partial\mathcal{Q}},$$

where $\nu(x) = (\nu_1(x), \dots, \nu_n(x))$ denotes the unit outer normal to $\partial\mathcal{Q}$ at $x \in \partial\mathcal{Q}$. Lax and Phillips [6, 7] formulated the scattering theory for the scalar-valued wave equation. The purpose of this paper is to make an analogous formulation for the elastic wave equation (0.1).

Throughout this paper, we assume that the space dimension $n \geq 3$. Furthermore, it is assumed that

$$(A.1) \quad a_{ipjq}(x) = a_{pijq}(x) = a_{jqip}(x), \quad i, j, p, q = 1, \dots, n \text{ (hyperelasticity)},$$

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$$(A.2) \quad \sum_{i,j,p,q=1}^n a_{ipjq}(x) \varepsilon_{jq} \overline{\varepsilon_{ip}} \geq \delta_1 \sum_{i,p=1}^n |\varepsilon_{ip}|^2 \quad (\text{stability}),$$

where a_{ipjq} denote the (p, q) components of the $n \times n$ matrix a_{ij} , (ε_{ip}) is any $n \times n$ Hermitian matrix and δ_1 is some positive constant independent of ε_{ip} .

Lax and Phillips [6, 7] employed the following energy for the data $f = (f_1, f_2)$:

$$\frac{1}{2} \int_{\Omega} (|\nabla_x f_1|^2 + |f_2|^2) dx \quad (\nabla_x f_1 = (\partial_1 f_1, \dots, \partial_n f_1)).$$

This played an essential role on their formulation. For the elastic equation, the energy is defined (in physics) by

$$\frac{1}{2} \int_{\Omega} \left\{ \sum_{i,j,p,q=1}^n a_{ipjq} \varepsilon_{jq}(f_1) \overline{\varepsilon_{ip}(f_1)} + |f_2|^2 \right\} dx,$$

where $f_i = (f_{i1}, \dots, f_{in})$, $i=1, 2$, and $\varepsilon_{ip}(f_1) = \frac{1}{2}(\partial_i f_{1p} + \partial_p f_{1i})$ called the stress tensor. By (A.1) we see that this energy is equal to

$$\|f\|_{E,\Omega}^2 = \frac{1}{2} \int_{\Omega} \left(\sum_{i,p,j,q=1}^n a_{ipjq} \partial_j f_{1q} \overline{\partial_i f_{1p}} + |f_2|^2 \right) dx.$$

From (A.2) it follows that

$$\|f\|_{E,\Omega}^2 \geq \delta_1 \sum_{i,p=1}^n \|\varepsilon_{ip}(f_1)\|_{L^2(\Omega)}^2 + \|f_2\|_{L^2(\Omega)}^2,$$

and then the energy is non-negative for any $f = (f_1, f_2)$. However, it does not mean that $\|f\|_{E,\Omega}$ is a norm equivalent to $\|\partial_x^1 f_1\|_{L^2(\Omega)} + \|f_2\|_{L^2(\Omega)}$ ($\partial_x^1 f_1 = (\nabla_x f_{11}, \dots, \nabla_x f_{1n})$). To prove this equivalence, we need to verify the inequality:

$$(0.4) \quad \sum_{i,p=1}^n \|\varepsilon_{ip}(f_1)\|_{L^2(\Omega)}^2 \geq \delta_2 \|\partial_x^1 f_1\|_{L^2(\Omega)}^2,$$

where δ_2 is some positive constant independent of f_1 . The estimate (0.4) is the key to the construction of the scattering theory for the elastic equation. When the domain Ω is bounded and the boundary condition is the displacement one (the case (0.2)), (0.4) is well-known as Korn's first inequality (cf. § 3 of Duvaut and Lions [1]). We shall give the proof of (0.4) in § 1 below. K.O. Friedrichs [2] also proved Korn's first inequality in the exterior domain. But the function spaces are quite different, and our result cannot follow from his.

We introduce the Hilbert space H of the data $f = (f_1, f_2)$ defined as the completion of $\{f \in C^\infty(\bar{\Omega}) \mid Bf_1 = 0 \text{ on } \partial\Omega, \text{ supp } [f] \text{ is bounded}\}$ in the norm $\|f\|_{E,\Omega}$. This space H coincides with $H_B^1(\Omega) \times L^2(\Omega)$, where $H_B^1(\Omega)$ is defined by

$$H_B^1(\mathcal{Q}) = \{u \in H_{loc}^1(\mathcal{Q}) \mid \partial_x^1 u \in L^2(\mathcal{Q}), \lim_{R \rightarrow \infty} R^{-2} \int_{R \leq |x| \leq 2R} |u(x)|^2 dx = 0, \\ u|_{\partial\mathcal{Q}} = 0 \text{ on } \partial\mathcal{Q} \text{ only when } Bu = u|_{\partial\mathcal{Q}}\} .$$

$H_B^1(\mathcal{Q})$ is complete in the norm $\|\partial_x^1 u\|_{L^2(\mathcal{Q})}$ (cf. Corollary 1.6 and Theorem 1.9 in § 1 below). We define

$$H_B^2(\mathcal{Q}) = \{u \in H_B^1(\mathcal{Q}) \mid \partial_x^1 u \in H^1(\mathcal{Q}), Bu = 0 \text{ on } \partial\mathcal{Q}\} .$$

Set

$$(0.5) \quad A = \begin{bmatrix} 0 & I \\ L & 0 \end{bmatrix}, \quad v(t) = \begin{bmatrix} u(t, \cdot) \\ \partial_t u(t, \cdot) \end{bmatrix} .$$

Then the equation (0.1) is transformed into

$$(0.6) \quad \frac{dv(t)}{dt} = Av(t), \quad v(0) = f .$$

The domain of A is defined by $D(A) = H_B^2(\mathcal{Q}) \times (H_B^1(\mathcal{Q}) \cap L^2(\mathcal{Q}))$. Then, A is a closed operator in H ; furthermore, A is skew self-adjoint in H (cf. Theorem 1.11 in § 1 below). Therefore, by Stone’s theorem (cf. Appendix I in Lax and Phillips [6]), we see that A generates a group $\{U(t)\}_{t \in \mathbb{R}}$ of unitary operators on H , and that for any $f \in D(A)$ $v(t) = U(t)f$ becomes a H -valued C^1 function and the solution of (0.6).

Using the Radon transformation: $u(x) \rightarrow \tilde{u}(s, \omega)$, Lax and Phillips [6, 7] constructed concretely the translation representations for the scalar-valued wave equation. In § 2, we shall construct the analogous representations for (0.1) in the unperturbed case (i.e., $\mathcal{Q} = \mathbb{R}^n$, and $a_{ij} = a_{ij}^0$ are constant) and study their properties under an additional assumption:

$$(A.3) \quad L^0(\xi) = \sum_{i,j=1}^n a_{ij}^0 \xi_i \xi_j \text{ has eigenvalues of constant} \\ \text{multiplicity for } \xi \in \mathbb{R}^n - \{0\} .$$

In the elastic wave equation case, there exist waves of the different modes (i.e., $L^0(\xi)$ may have different eigenvalues). We need to notice this phenomenon when defining the translation representations and studying their properties; however, in this process we do not encounter a serious difficulty caused by that phenomenon. The idea of our methods is essentially the same as in Chapter VI of Lax and Phillips [6], which deals with the scattering for symmetric hyperbolic systems of first order.

In § 3, assuming that the coefficients of L are constant, we shall consider the scattering for (0.1). Set

$$D_{\pm} = \{f \in H \mid \text{supp } [U(t)f(x)] \subset \{(t, x) \mid \pm \eta t + r_0 \leq |x|\} \text{ if } \pm t > 0\},$$

where η^2 ($\eta > 0$) is the minimum of the eigenvalues of $L^0(\omega)$ ($\omega \in S^{n-1}$). Then it is seen from the discussion in the unperturbed case (in § 2) that D_+ (resp. D_-) has the properties of the outgoing (resp. incoming) subspace for $U(t)$, except

$$(0.7) \quad \overline{\bigcup_t U(t)D_{\pm}} = H.$$

By this equality we can construct the translation representations for $U(t)$, and also can derive the completeness of the wave operators in the same way as in Lax and Phillips [6, 7]. Thus, the main task in § 3 is to prove (0.7). The proof of (0.7) can be reduced to verifying non-existence of the point spectrums of A (cf. § 2 of Chapter V in Lax and Phillips [6]), and therefore we have only to show that there exist no eigen-functions of A in the domain $D(A)$ (cf. Theorem 3.5 in § 3 below). The key to this show is to prove that if $f \in D(A)$ satisfies $(A - i\sigma)f(x) = 0$ for every large $|x|$, then $f(x) = 0$ for every large $|x|$. The methods in Lax and Phillips [6] are not applicable to this proof in our case. Multiplying $A - i\sigma I$ by its cofactor, we transform it into a diagonal operator and carry out the proof by means of a uniqueness theorem for single equations of higher order obtained in Littman [8], Hörmander [4] and Murata [10].

Yamamoto [14] makes a related study. He considers the isotropic equation in three dimensional space (i.e., $a_{ipjq}(x) = \mu(\delta_{pq}\delta_{ij} + \delta_{iq}\delta_{jp}) + \lambda\delta_{ip}\delta_{jq}$, δ_{ip} being Kronecker's delta; λ and μ being the Lamé constants, and $n = 3$) with the displacement boundary condition (the case (0.2)), which is contained in our case. And then, he obtains the same results as ours. But it seems difficult to apply his methods to the traction boundary condition case (the case (0.3)). For, he does not derive the estimate (0.4), which is essential in the traction boundary condition case.

We note that Iwashita and Shibata [5] investigate the analyticity of spectral functions of $L + \sigma^2 I$ and the rate of the local energy decay of the solutions to (0.1).

§ 1. Spaces of the Data and Properties of the Generator A

In what follows, the Roman letters u, v, w and the Greek letters ϕ, ψ are used to denote n -dimensional row vector and scalar-valued functions, respectively.

For any domain G in \mathbf{R}^n , $C_0^\infty(G)$ denotes the space of all C^∞ functions on \mathbf{R}^n whose supports are compact and lie in G . $C_0^\infty(\bar{G})$ denotes the space of all functions in $C^\infty(\bar{G})$ whose supports are compact. In particular, functions in $C_0^\infty(G)$ vanish near the boundary of G . By $H^m(G)$ we denote the Sobolev space of order m on G ; put $H_{loc}^m(G) = \{\phi \mid \psi\phi \in H^m(G) \text{ for any } \psi \in C_0^\infty(\mathbf{R}^n)\}$. Set $\Omega_R = \{x \in \Omega \mid |x| < R\}$ ($R \geq r_0$). For any space Γ of scalar valued functions we abbreviate the product space $\Gamma \times \cdots \times \Gamma$ by also Γ . $\chi(x)$ will always refer to a real and scalar valued function in $C_0^\infty(\mathbf{R}^n)$ such that $0 \leq \chi \leq 1$, $\chi(x) = 1$ for $|x| \leq 1$ and $= 0$ for $|x| \geq 2$.

As is stated in § 0, the purpose of this section is to discuss the skew self-adjointness of the operator A given in (0.5). Set

$$(1.1) \quad \dot{H}^m(\Omega) = \{u \in H_{loc}^m(\Omega) \mid \partial_x^\alpha u \in L^2(\Omega) \text{ for any multi-index } \alpha \text{ with } 1 \leq |\alpha| \leq m, \\ \lim_{R \rightarrow \infty} R^{-2} \int_{R \leq |x| \leq 2R} |u(x)|^2 dx = 0\}.$$

Put

$$(u, v)_{m,\Omega} = \sum_{1 \leq |\alpha| \leq m} \sum_{j=1}^n \int_{\Omega} \partial_x^\alpha u_j(x) \overline{\partial_x^\alpha v_j(x)} dx, \quad |u|_{m,\Omega}^2 = (u, u)_{m,\Omega}$$

for any $u = {}^t(u_1, \dots, u_n)$ and $v = {}^t(v_1, \dots, v_n) \in \dot{H}^m(\Omega)$ (m being integers ≥ 1).

Theorem 1.1. *Let m be an integer ≥ 1 . Then, $\dot{H}^m(\Omega)$ is a Hilbert space equipped with the innerproduct $(\ , \)_{m,\Omega}$. Furthermore, $C_0^\infty(\bar{\Omega})$ is dense in $\dot{H}^m(\Omega)$.*

To prove Theorem 1.1, we need the following three technical lemmas.

Lemma 1.2. *For any $u \in \dot{H}^m(\Omega)$ ($m \geq 1$) and $R \geq r_0$, there exists a sequence $\{u_{k,R}\}_{k=1,2,\dots} \subset C_0^\infty(\bar{\Omega})$ such that*

$$|u_{k,R} - u|_{m,\Omega} \rightarrow 0 \quad \text{and} \quad \|u_{k,R} - u\|_{L^2(\Omega_R)} \rightarrow 0 \quad \text{as } k \rightarrow \infty.$$

Lemma 1.3. *For any $u \in \dot{H}^1(\Omega)$ and $R \geq r_0 + 1$,*

$$(1.2) \quad \int_{|x| \geq R} |u(x)|^2 |x|^{-2} dx \leq 4 \int_{|x| \geq R} |\partial_x^1 u(x)|^2 dx.$$

Lemma 1.4. *There exists a constant $C > 0$ independent of $R \geq r_0$ such that*

$$(1.3) \quad \|u\|_{L^2(\Omega_R)} \leq CR \|\partial_x^1 u\|_{L^2(\Omega)} \quad \text{for any } u \in \dot{H}^1(\Omega) \text{ and } R \geq r_0.$$

Deferring the proofs of Lemmas 1.2–1.4, we give

Proof of Theorem 1.1. From Lemma 1.4 we see that if $(u, u)_{m,\Omega} = 0$, then

$u=0$. Thus, we can easily check that $(u, v)_{m,\Omega}$ is an innerproduct. In view of Lemma 1.2, for any $u \in \dot{H}^m(\Omega)$ we can choose a sequence $\{u_{k,r_0}\}_{k=1,2,\dots} \subset C_0^\infty(\bar{\Omega})$ so that $\|u_{k,r_0} - u\|_{m,\Omega} \rightarrow 0$ as $k \rightarrow \infty$, from which it follows that $C_0^\infty(\bar{\Omega})$ is dense in $\dot{H}^m(\Omega)$. The rest of our task is to prove the completeness.

Let $\{u_k\}_{k=1,2,\dots}$ be any Cauchy sequence in $\dot{H}^m(\Omega)$. Since

$$\|u_k - u_l\|_{L^2(\Omega_R)} \leq CR \|\partial_x^1(u_k - u_l)\|_{L^2(\Omega)} \leq CR \|u_k - u_l\|_{m,\Omega}$$

as follows from Lemma 1.4, $\{u_k\}$ becomes a Cauchy sequence in $H^m(\Omega_R)$ for any $R \geq r_0 + 1$. Let u_R denote the limits of $\{u_k\}$ in $H^m(\Omega_R)$ and set $u(x) = u_R(x)$ for $x \in \Omega_R$. Then, $u(x)$ is well-defined as a function on Ω and $\in H_{loc}^m(\Omega)$. It is easily seen that $\partial_x^\alpha u(x) \in L^2(\Omega)$ for $1 \leq |\alpha| \leq m$ and that

(1.4.a) $\|u_k - u\|_{m,\Omega} \rightarrow 0,$

(1.4.b) $\|u_k - u\|_{L^2(\Omega_R)} \rightarrow 0$ for any $R \geq r_0 + 1$

as $k \rightarrow \infty$. We must prove that $u(x)$ satisfies (1.1). Noting that $u_k \in \dot{H}^1(\Omega)$, by Lemma 1.3 we have

$$\begin{aligned} (1.5) \quad & \{R^{-2} \int_{R \leq |x| \leq 2R} |u(x)|^2 dx\}^{1/2} \\ & \leq \{R^{-2} \int_{R \leq |x| \leq 2R} |u(x) - u_k(x)|^2 dx\}^{1/2} + 4 \{ \int_{|x| \geq R} |\partial_x^1 u_k(x)|^2 dx \}^{1/2} \\ & \leq R^{-1} \|u - u_k\|_{L^2(\Omega_{2R})} + 4 \|u_k - u\|_{m,\Omega} + 4 \{ \int_{|x| \geq R} |\partial_x^1 u(x)|^2 dx \}^{1/2}. \end{aligned}$$

Letting $k \rightarrow \infty$ in (1.5) and using (1.4), we have

$$(1.6) \quad R^{-2} \int_{R \leq |x| \leq 2R} |u(x)|^2 dx \leq 16 \int_{|x| \geq R} |\partial_x^1 u(x)|^2 dx.$$

Since we already know that $|\partial_x^1 u(x)| \in L^2(\Omega)$, it follows from (1.6) that u satisfies (1.1), which completes the proof of the theorem.

Now, we shall prove Lemmas 1.2–1.4.

Proof of Lemma 1.2. Set $\chi_r(x) = \chi(r^{-1}x)$. Then, it follows that for some constant C $\|\chi_r u - u\|_{m,\Omega}^2 \leq C(I_1(r) + I_2(r))$ where

$$\begin{aligned} I_1(r) &= \sum_{\substack{1 \leq |\alpha| + |\beta| \leq m \\ |\beta| \geq 1}} \int_{\Omega} |\partial_x^\alpha (1 - \chi_r(x)) \partial_x^\beta u(x)|^2 dx, \\ I_2(r) &= \sum_{1 \leq |\alpha| \leq m} \int_{\Omega} |\partial_x^\alpha \chi_r(x) u(x)|^2 dx. \end{aligned}$$

Since $\partial_x^\beta u(x) \in L^2(\Omega)$ for $|\beta| \geq 1$, it is seen that $I_1(r) \rightarrow 0$ as $r \rightarrow \infty$. Noting that

$$I_2(r) \leq \sum_{1 \leq |\alpha| \leq m} \sup |\partial_x^\alpha \chi|^{-2|\alpha|} \int_{r \leq |x| \leq 2r} |u(x)|^2 dx,$$

by (1.1) we have also that $I_2(r) \rightarrow 0$ as $r \rightarrow \infty$. These imply that for any integer $k \geq 1$ there exists a $v_k \in H^m(\Omega)$ satisfying

$$(1.7) \quad |u - v_k|_{m, \Omega} \leq 1/k \quad \text{and} \quad v_k(x) = u(x) \quad \text{for} \quad x \in \Omega_R.$$

As is well-known, $C_0^\infty(\bar{\Omega})$ is dense in $H^m(\Omega)$. Thus, we can find $u_{k,R} \in C_0^\infty(\bar{\Omega})$ such that

$$(1.8) \quad \sum_{|\alpha| \leq m} \|\partial_\alpha^x (u_{k,R} - v_k)\|_{L^2(\Omega)} \leq 1/k.$$

Combining (1.7) and (1.8), we see easily that the sequence $\{u_{k,R}\}$ has the desired properties.

Proof of Lemma 1.3. First, we shall prove that

$$(1.9) \quad \int_{|x| \geq R} |\phi(x)|^2 |x|^{-2l} dx \leq \left\{ \frac{2}{n(1-(2l/n))} \right\}^2 \int_{|x| \geq R} |\nabla_x \phi(x)|^2 |x|^{-2(l-1)} dx$$

for any $\phi \in C_0^\infty(\bar{\Omega})$, $R \geq R_0$ and $l < n/2$. We use the polar coordinates: $r = |x|$ and $\omega = x/r$. Since $|\phi(r\omega)|^2 r^{-2l} = - \int_r^\infty \frac{\partial}{\partial s} [|\phi(s\omega)|^2 s^{-2l}] ds$, we have

$$(1.10) \quad \int_R^\infty |\phi(r\omega)|^2 r^{n-2l-1} dr \leq 2 \int_R^\infty r^{n-1} dr \int_r^\infty |\partial_s \phi(s\omega)| |\phi(s\omega)| s^{-2l} ds + 2l \int_R^\infty r^{n-1} dr \int_r^\infty |\phi(s\omega)|^2 s^{-2l-1} ds.$$

By integration by parts we have

$$(1.11) \quad 2 \int_R^\infty r^{n-1} dr \int_r^\infty |\partial_s \phi(s\omega)| |\phi(s\omega)| s^{-2l} ds \leq \frac{2}{n} \int_R^\infty |(\nabla_x \phi)(r\omega)| |\phi(r\omega)| r^{n-2l} dr,$$

$$(1.12) \quad 2l \int_R^\infty r^{n-1} dr \int_r^\infty |\phi(s\omega)|^2 s^{-2l-1} ds \leq \frac{2l}{n} \int_R^\infty |\phi(r\omega)|^2 r^{n-2l-1} dr.$$

By Schwarz's inequality we have

$$(1.13) \quad \begin{aligned} & \text{the right-hand side of (1.11)} \\ & \leq \frac{2}{n} \left\{ \int_R^\infty |(\nabla_x \phi)(r\omega)|^2 r^{n-2l+1} dr \right\}^{1/2} \left\{ \int_R^\infty |\phi(r\omega)|^2 r^{n-2l-1} dr \right\}^{1/2}. \end{aligned}$$

Combining (1.10)–(1.13) implies that

$$(1.14) \quad (1-(2l/n)) \left\{ \int_R^\infty |\phi(r\omega)|^2 r^{2n-2l-1} dr \right\}^{1/2} \leq \frac{2}{n} \left\{ \int_R^\infty |(\nabla_x \phi)(r\omega)|^2 r^{n-2l+1} dr \right\}^{1/2}.$$

Since $1-(2l/n) > 0$, (1.9) follows immediately from (1.14). In particular, if we take $l=1$ (note that $n \geq 3$), from (1.9) we have that (1.2) is valid for any $u \in C_0^\infty(\bar{\mathcal{D}})$.

Now, using Lemma 1.2, we shall prove that (1.2) is also valid for any $u \in \dot{H}^1(\mathcal{Q})$. Let R' be any number $> R$. By Lemma 1.2 we know that there exists a sequence $\{u_{k,R'}\}_{k=1,2,\dots} \subset C_0^\infty(\bar{\mathcal{D}})$ such that

$$(1.15) \quad |u_{k,R'} - u|_{1,\mathcal{Q}} \rightarrow 0 \quad \text{and} \quad \|u_{k,R'} - u\|_{L^2(\mathcal{Q}_{R'})} \rightarrow 0 \quad \text{as} \quad k \rightarrow \infty.$$

On the other hand, we have

$$(1.16) \quad \left\{ \int_{R \leq |x| \leq R'} |u(x)|^2 |x|^{-2} dx \right\}^{1/2} \\ \leq \left\{ \int_{R \leq |x| \leq R'} |u(x) - u_{k,R'}(x)|^2 |x|^{-2} dx \right\}^{1/2} \\ + \left\{ \int_{R \leq |x| \leq R'} |u_{k,R'}(x)|^2 |x|^{-2} dx \right\}^{1/2} \\ \leq R^{-1} \|u_{k,R'} - u\|_{L^2(\mathcal{Q}_{R'})} + 2|u_{k,R'} - u|_{1,\mathcal{Q}} + \left\{ 4 \int_{|x| \geq R} |\partial_x^1 u(x)|^2 dx \right\}^{1/2}.$$

Letting $k \rightarrow \infty$ in (1.16) and using (1.15), we have

$$\int_{R \leq |x| \leq R'} |u(x)|^2 |x|^{-2} dx \leq 4 \int_{|x| \geq R} |\partial_x^1 u(x)|^2 dx.$$

The arbitrariness of the choice of $R' > R$ implies the lemma.

Proof of Lemma 1.4. In view of Lemma 1.3, it suffices to prove that there exists a constant $C > 0$ such that

$$(1.17) \quad \|u\|_{L^2(\mathcal{Q}_{r_0+1})} \leq C \{ \|\partial_x^1 u\|_{L^2(\mathcal{Q}_{r_0+1})} + \|u\|_{L^2(\mathcal{G})} \}$$

for any $u \in H^1(\mathcal{Q}_{r_0+1})$ where $\mathcal{G} = \{x \in \mathbb{R}^n \mid r_0 \leq |x| \leq r_0 + 1\}$. In fact, if (1.17) is valid, then we have

$$\|u\|_{L^2(\mathcal{Q}_R)}^2 = \|u\|_{L^2(\mathcal{Q}_{r_0+1})}^2 + \|u\|_{L^2(\{x \in \mathbb{R}^n \mid r_0+1 \leq |x| \leq R\})}^2 \\ \leq 2C^2 \{ \|\partial_x^1 u\|_{L^2(\mathcal{Q}_{r_0+1})}^2 + (r_0+1)^2 \int_{\mathcal{G}} |u(x)|^2 |x|^{-2} dx \} + R^2 \int_{r_0+1 \leq |x| \leq R} |u(x)|^2 |x|^{-2} dx \\ \leq (2C^2 + 8C^2(r_0+1)^2 + 4R^2) \|\partial_x^1 u\|_{L^2(\mathcal{Q})}^2$$

for any $u \in \dot{H}^1(\Omega)$ and $R \geq r_0 + 1$, from which the lemma follows.

Now, we prove (1.17). Suppose to the contrary that for any integer $k \geq 1$ there exists a $u_k \in H^1(\Omega_{r_0+1})$ such that

$$(1.18) \quad \|u_k\|_{L^2(\Omega_{r_0+1})} = 1,$$

$$(1.19) \quad \|\partial_x^1 u_k\|_{L^2(\Omega_{r_0+1})} + \|u_k\|_{L^2(G)} < 1/k.$$

Since $\{u_k\}$ is a bounded sequence in $H^1(\Omega_{r_0+1})$ as follows from (1.18) and (1.19), $\{u_k\}$ is weakly compact in $H^1(\Omega_{r_0+1})$. Furthermore, since Ω_{r_0+1} is bounded, by Rellich's compactness theorem we see that $\{u_k\}$ is strongly compact in $L^2(\Omega_{r_0+1})$. By passing to a subsequence if necessary, we can conclude that there exists a $v \in H^1(\Omega_{r_0+1})$ such that

$$(1.20) \quad u_k \rightarrow v \quad \text{weakly in } H^1(\Omega_{r_0+1}),$$

$$(1.21) \quad u_k \rightarrow v \quad \text{strongly in } L^2(\Omega_{r_0+1}),$$

as $k \rightarrow \infty$. It follows that from (1.19) and (1.20) that

$$\|\partial_x^1 v\|_{L^2(\Omega_{r_0+1})} + \|v\|_{L^2(G)} = 0,$$

which implies that $v=0$. In fact, the fact that $\partial_x^1 v=0$ implies that v is a constant vector. Thus, $|v|^2 \times (\text{the Lebesgue measure of } G) = 0$, from which we see that $v=0$.

On the other hand, (1.18) and (1.21) imply that the L^2 norm of v on Ω_{r_0+1} is equal to 1, which contradicts the fact that $v=0$. Thus, we have seen that (1.17) is valid, which completes the proof of the lemma.

Now, let us introduce the spaces and bilinear forms connected closely with the operator L and the boundary condition. We define

$$\begin{aligned} C_B^\infty(\bar{\Omega}) &= \{u \in C_0^\infty(\bar{\Omega}) \mid BL^k u = 0 \text{ on } \partial\Omega \text{ for any } k \geq 0\}, \\ H_B^\infty(\Omega) &= \{u \in \bigcap_{m=1}^\infty \dot{H}^m(\Omega) \mid BL^k u = 0 \text{ on } \partial\Omega \text{ for any } k \geq 0\}, \\ H_B^1(\Omega) &= \begin{cases} \{u \in \dot{H}^1(\Omega) \mid u = 0 \text{ on } \partial\Omega\} & \text{when } Bu = u|_{\partial\Omega}, \\ \dot{H}^1(\Omega) & \text{when } Bu = \sum_{i,j=1}^n \nu_i a_{ij} \partial_j u|_{\partial\Omega}, \end{cases} \\ H_B^2(\Omega) &= \{u \in \dot{H}^2(\Omega) \mid Bu = 0 \text{ on } \partial\Omega\}, \\ \langle u, v \rangle_{1,\Omega} &= \sum_{i,j,p,q=1}^n \int_{\Omega} a_{ipjq}(x) \partial_j u_q(x) \overline{\partial_i v_p(x)} dx, \\ \langle u, v \rangle_{2,\Omega} &= \langle u, v \rangle_{1,\Omega} + (Lu, Lv)_{L^2(\Omega)}, \end{aligned}$$

where $u = {}^t(u_1, \dots, u_n)$ and $v = {}^t(v_1, \dots, v_n)$ and $(\cdot, \cdot)_{L^2(\mathcal{Q})}$ is the usual innerproduct of $L^2(\mathcal{Q})$.

First, we would like to prove that $\langle u, u \rangle_{1, \mathcal{Q}}$ is the equivalent norm to $(u, u)_{1, \mathcal{Q}}$. To do this, we need

Theorem 1.5 (*Korn's first inequality*). *There exists a constant $C > 0$ such that*

$$C^{-1} \int_{\mathcal{Q}} |\partial_x^1 u|^2 dx \leq \sum_{i,j=1}^n \int_{\mathcal{Q}} |\varepsilon_{ij}(u)|^2 dx \leq C \int_{\mathcal{Q}} |\partial_x^1 u|^2 dx$$

for any $u = {}^t(u_1, \dots, u_n) \in \dot{H}^1(\mathcal{Q})$, where $\varepsilon_{ij}(u) = \frac{1}{2}(\partial_i u_j + \partial_j u_i)$.

When G is a bounded domain in \mathbb{R}^n , the following estimate is well-known as Korn's second inequality (cf. § 3 of Duvaut and Lions [1] or Nitsche [11]):

$$(1.22) \quad \int_G |\partial_x^1 u(x)|^2 dx \leq C \left\{ \sum_{i,j=1}^n \int_G |\varepsilon_{ij}(u)|^2 dx + \int_G |u(x)|^2 dx \right\}$$

for any $u \in H^1(G)$. Friedrichs [2] also derived Korn's first inequality in the unbounded domain for the functions u in the different classes. But, Theorem 1.5 does not follow directly from [2]. To prove Theorem 1.5 we need other ideas.

Proof of Theorem 1.5. It is trivial that the second part of the inequalities holds. Thus, we prove only the first part. Since $C_0^\infty(\bar{\mathcal{Q}})$ is dense in $\dot{H}^1(\mathcal{Q})$ as follows from Theorem 1.1, we may assume that $u \in C_0^\infty(\bar{\mathcal{Q}})$. Since by integration by parts we have for any $u \in C_0^\infty(\mathbb{R}^n)$

$$\sum_{i,j=1}^n \operatorname{Re} \int_{\mathbb{R}^n} \partial_i u_j \overline{\partial_j u_i} dx = \sum_{i,j=1}^n \int_{\mathbb{R}^n} \partial_j u_j \overline{\partial_i u_i} dx \geq 0,$$

it follows immediately that

$$(1.23) \quad \sum_{i,j=1}^n \int_{\mathbb{R}^n} |\varepsilon_{ij}(u)|^2 dx \geq \frac{1}{2} \int_{\mathbb{R}^n} |\partial_x^1 u|^2 dx \quad \text{for any } u \in C_0^\infty(\mathbb{R}^n).$$

Using (1.22) and (1.23), we can prove that there exists a constant $C_1 > 0$ such that

$$(1.24) \quad \|\partial_x^1 u\|_{L^2(\mathcal{Q})}^2 \leq C_1 \left\{ \sum_{i,j=1}^n \int_{\mathcal{Q}} |\varepsilon_{ij}(u)|^2 dx + \|u\|_{L^2(\Omega_{r_0+1})}^2 \right\}$$

for any $u \in C_0^\infty(\bar{\mathcal{Q}})$. In fact, if we choose real-valued C^∞ functions ϕ, ψ on \mathbb{R}^n so that $\phi(x) = 0$ for $|x| \geq r_0 + 1$, $\psi(x) = 0$ for $|x| \leq r_0$ and $\phi^2(x) + \psi^2(x) = 1$, then

$$\|\partial_x^1 u\|_{L^2(\mathcal{Q})}^2 \leq \int_{\Omega_{r_0+1}} |\partial_x^1(\phi u)|^2 dx + \int_{\mathbb{R}^n} |\partial_x^1(\psi u)|^2 dx + C \int_{\Omega_{r_0+1}} |u|^2 dx,$$

where $C = \sup |\nabla_x \phi|^2 + \sup |\nabla_x \psi|^2$. Applying (1.22) and (1.23) and noting that $|\varepsilon_{ij}(\phi u)|^2 + |\varepsilon_{ij}(\psi u)|^2 \leq 2(\phi^2 + \psi^2)|\varepsilon_{ij}(u)|^2 + 4(|\nabla_x \phi|^2 + |\nabla_x \psi|^2)|u|^2$ we have (1.24).

In view of (1.24), to complete the proof we have only to prove that there exists a constant $C_2 > 0$ such that

$$(1.25) \quad \|u\|_{L^2(\Omega_{r_0+1})} \leq C_2 \sum_{i,j=1}^n \int_{\Omega} |\varepsilon_{ij}(u)|^2 dx \quad \text{for any } u \in C_0^\infty(\bar{\Omega}).$$

Suppose to the contrary that for any integer $k \geq 1$ there exists a $u_k \in C_0^\infty(\bar{\Omega})$ such that

$$(1.26) \quad \|u_k\|_{L^2(\Omega_{r_0+1})} = 1,$$

$$(1.27) \quad \sum_{i,j=1}^n \int_{\Omega} |\varepsilon_{ij}(u_k)|^2 dx < 1/k.$$

Combining (1.24), (1.26) and (1.27), we see that $\{u_k\}_{k=1,2,\dots}$ is a bounded set in $H^1(\Omega_{r_0+1})$. Thus, by Rellich's compactness theorem we may assume that $\{u_k\}$ is a Cauchy sequence in $L^2(\Omega_{r_0+1})$. Applying (1.24) to $u_k - u_l$ and using (1.27), we have

$$\|u_k - u_l\|_{L^2(\Omega_{r_0+1})}^2 \leq C_2[(1/k) + (1/l) + \|u_k - u_l\|_{L^2(\Omega_{r_0+1})}^2].$$

This implies that $\{u_k\}$ is a Cauchy sequence in $\dot{H}^1(\Omega)$. By Theorem 1.1 we see that there exists a $w = {}^t(w_1, \dots, w_n) \in \dot{H}^1(\Omega)$ such that $\|u_k - w\|_{L^2(\Omega)} \rightarrow 0$ as $k \rightarrow \infty$. (1.27) implies that

$$(1.28) \quad \partial_i w_j(x) + \partial_j w_i(x) = 0 \quad \text{in } \Omega \text{ for any } i, j = 1, \dots, n.$$

Since $\partial_i \partial_j w_k = -\partial_i \partial_k w_j = \partial_k \partial_j w_i = -\partial_i \partial_j w_k$ as follows from (1.28), we see that $\partial_i \partial_j w_k = 0$. These mean that w_k are polynomials of order at most 1. Noting that $\nabla_x w_k \in L^2(\Omega)$, we have that w_k are constant. Then, we have

$$R^{-2} \int_{R \leq |x| \leq 2R} |w(x)|^2 dx = |w|^2 R^{n-2} c_n (2^n - 1),$$

where c_n is the Lebesgue measure of the unit ball. Since $n \geq 3$, by (1.1) we have that $w = 0$. Namely, we can conclude that $\|u_k\|_{L^2(\Omega)} \rightarrow 0$ as $k \rightarrow \infty$.

On the other hand, by (1.26) and Lemma 1.4 we have

$$1 = \|u_k\|_{L^2(\Omega_{r_0+1})} \leq C(r_0 + 1) \|u_k\|_{L^2(\Omega)},$$

which leads the contradiction. Thus, we have proved (1.25), which completes the proof.

Corollary 1.6. *Assume that (A.1) and (A.2) are valid. Then, there exists a constant $C > 0$ such that*

$$C^{-1}|u|_{1,\Omega}^2 \leq \langle u, u \rangle_{1,\Omega} \leq C|u|_{1,\Omega}^2 \quad \text{for any } u \in \dot{H}^1(\Omega).$$

Proof. By (A.1) and (A.2) we have

$$\langle u, u \rangle_{1,\Omega}^2 = \sum_{i,j,p,q=1}^n \int_{\Omega} a_{ipjq} \varepsilon_{jq}(u) \overline{\varepsilon_{ip}(u)} dx \geq \delta \sum_{i,j=1}^n \int_{\Omega} |\varepsilon_{ij}(u)|^2 dx.$$

Combining this and Theorem 1.5 implies that the first part of the inequalities is valid. The second part is trivial. The proof is complete.

When getting fundamental properties of the spaces $H_B^1(\Omega)$ and $H_B^2(\Omega)$ and proving the skew self-adjointness of A , the existence theorem of the solutions to the following problem plays an essential role:

$$(1.29) \quad Lu = v \text{ in } \Omega, \quad Bu = 0 \text{ on } \partial\Omega.$$

To solve (1.29), first we consider the variational equation:

$$(1.30) \quad \langle u, w \rangle_{1,\Omega} = -(v, w)_{L^2(\Omega)} \quad \text{for any } w \in H_B^1(\Omega).$$

By Theorem 1.1 and Corollary 1.6 we see that $H_B^1(\Omega)$ is a Hilbert space equipped with the innerproduct $\langle \cdot, \cdot \rangle_{1,\Omega}$. It follows from Riesz's representation theorem that for any $v \in L_R^2(\Omega) = \{v \in L^2(\Omega) \mid v(x) = 0 \text{ for } |x| > R\}$ ($R > r_0$) there exists a unique $u \in H_B^1(\Omega)$ satisfying (1.30). Since

$$|(v, w)_{L^2(\Omega)}| \leq \|v\|_{L^2(\Omega)} \|w\|_{L^2(\Omega_R)} \leq CR \|v\|_{L^2(\Omega)} |w|_{1,\Omega}$$

as follows from Lemma 1.4, putting $u = w$ in (1.30) and using Corollary 1.6, we have

$$(1.31) \quad |u|_{1,\Omega} \leq CR \|v\|_{L^2(\Omega)}.$$

Theorem 1.7. *Assume that (A.1) and (A.2) are valid. Let R be any number $> r_0$. If $v \in L_R^2(\Omega)$, then (1.29) admits a unique solution $u \in H_B^2(\Omega)$. In addition, we assume that $v \in H^m(\Omega)$ ($m \geq 1$). Then, $u \in \dot{H}^{m+2}(\Omega)$.*

Furthermore, there exists a constant $C > 0$ independent of R , u and v such that

$$(1.32) \quad \sum_{|\alpha|=2} \|\partial_x^\alpha u\|_{L^2(\Omega)} \leq CR \|v\|_{L^2(\Omega)}.$$

Gilberg and Trudinger [3] obtain the same results as in this theorem in the case where $Bu = u|_{\partial\Omega}$ (cf. Theorems 8.8, 8.12 and 8.13 of [3]), and Shibata

[12] in the case where $Bu = \sum \nu_i a_{ij} \partial_j u|_{\partial\Omega}$ (cf. Theorem 3.4 of [12]). Since, in the above discussion about the variational equation (1.30), we have got a unique u in $H_B^1(\Omega)$ satisfying (1.30), the main part of the proof is to show that $\partial_x^\alpha u \in L^2(\Omega)$ for any α with $|\alpha| = 2$. For detailed proof, see [3] and [12].

Using Corollary 1.6 and Theorem 1.7, we can prove

Lemma 1.8. *Assume that (A.1) and (A.2) are valid. Then, there exists a constant $C > 0$ such that*

$$C^{-1} |u|_{2,\Omega}^2 \leq \langle u, u \rangle_{2,\Omega} \leq C |u|_{2,\Omega}^2 \quad \text{for any } u \in H_B^2(\Omega).$$

Proof. The second part of the inequalities is trivial. Our task is to verify the first part. Choose $\phi \in C_0^\infty(\mathbb{R}^n)$ so that $\phi(x) = 1$ for $|x| \leq r_0 + 2$ and $= 0$ for $|x| \geq r_0 + 3$. For $u \in H_B^2(\Omega)$, put $v^1 = -\phi Lu$, $v^2 = \phi Lu - L(\phi u)$. Then, $v^1, v^2 \in L_{r_0+3}^2(\Omega)$, and by integration by parts we have

$$\langle \phi u, w \rangle_{1,\Omega} = (v^1 + v^2, w)_{L^2(\Omega)} \quad \text{for any } w \in H_B^1(\Omega).$$

The uniqueness of the solution of (1.30) and the estimates (1.31), (1.32) yield that

$$(1.33) \quad \begin{aligned} |\phi u|_{2,\Omega} &\leq C_1(r_0 + 3)(\|v^1\|_{L^2(\Omega)} + \|v^2\|_{L^2(\Omega)}) \\ &\leq C_2(\|Lu\|_{L^2(\Omega)} + |u|_{1,\Omega} + \|u\|_{L^2(\Omega_{r_0+3})}). \end{aligned}$$

Next let us estimate $v = (1 - \phi)u$. Noting that $\text{supp } [v] \subset \{x \mid |x| \geq r_0 + 2\}$, by Corollary 1.6 and integration by parts we have for $i = 1, 2, \dots, n$,

$$\begin{aligned} \sum_{j=1}^n \|\partial_j \partial_i v\|_{L^2(\Omega)}^2 &\leq C_3 \langle \partial_i v, \partial_i v \rangle_{1,\Omega} \\ &= C_3 \{ (Lv, \partial_i^2 v)_{L^2(\Omega)} + \sum_{k,l,p,q=1}^n \int_{\Omega} (\partial_i a_{kplq}) \partial_l v_q \overline{\partial_k \partial_p v} dx \} \\ &\leq C_4 \{ \|Lu\|_{L^2(\Omega)} + |u|_{1,\Omega} + \|u\|_{L^2(\Omega_{r_0+3})} \} \sum_{|\alpha|=2} \|\partial_x^\alpha v\|_{L^2(\Omega)}. \end{aligned}$$

From this it follows that

$$(1.34) \quad |(1 - \phi)u|_{2,\Omega} \leq C_5 \{ \|Lu\|_{L^2(\Omega)} + |u|_{1,\Omega} + \|u\|_{L^2(\Omega_{r_0+3})} \}.$$

Combining (1.33), (1.34) and Lemma 1.4, we obtain

$$|u|_{2,\Omega} \leq C_6 \{ \|Lu\|_{L^2(\Omega)} + |u|_{1,\Omega} \}.$$

From this and Corollary 1.6 the lemma follows.

Now, we summarize fundamental properties of the spaces $H_B^1(\Omega)$ and $H_B^2(\Omega)$.

Theorem 1.9. *Assume that (A.1) and (A.2) are valid. Then, $H_B^1(\mathcal{Q})$ and $H_B^2(\mathcal{Q})$ are Hilbert spaces equipped with the innerproducts $\langle \cdot, \cdot \rangle_{1,\mathcal{Q}}$ and $\langle \cdot, \cdot \rangle_{2,\mathcal{Q}}$, respectively.*

Furthermore, $H_B^\infty(\mathcal{Q})$ and $C_B^\infty(\bar{\mathcal{Q}})$ are dense in $H_B^\infty(\mathcal{Q})$ for $k=1, 2$.

Proof. The first assertion immediately follows from Theorem 1.1, Corollary 1.6 and Lemma 1.8. Our task is to prove the second assertion.

At first, we shall prove that $H_B^\infty(\mathcal{Q})$ is dense in $H_B^1(\mathcal{Q})$. Suppose to the contrary that there exists a non-trivial $u \in H_B^1(\mathcal{Q})$ such that

$$(1.35) \quad \langle u, v \rangle_{1,\mathcal{Q}} = 0 \quad \text{for any } v \in H_B^\infty(\mathcal{Q}).$$

By Theorem 1.7, for any $w \in C_0^\infty(\mathcal{Q})$ we have a solution $v \in H_B^\infty(\mathcal{Q})$ of the equation: $Lv = w$ in \mathcal{Q} . By integration by parts we obtain

$$(1.36) \quad \sum_{i,j,p,q=1}^n \int_{\mathcal{Q}} \chi(x/R) a_{ipjq}(x) \partial_j u_q(x) \overline{\partial_i v_p(x)} dx \\ = -(\chi(\cdot/R)u, Lv)_{L^2(\mathcal{Q})} - \sum_{i,j,p,q=1}^n R^{-1} \int_{\mathcal{Q}} (\partial_j \chi)(x/R) a_{ipjq}(x) u_q(x) \overline{\partial_i v_p(x)} dx$$

for any $R > r_0$. Noting that u satisfies (1.1) and letting $R \rightarrow \infty$ in (1.36), we get

$$(1.37) \quad \langle u, v \rangle_{1,\mathcal{Q}} = -(u, w)_{L^2(\mathcal{Q})}.$$

Combining (1.35) and (1.37) yields that $(u, w)_{L^2(\mathcal{Q})} = 0$ for any $w \in C_0^\infty(\mathcal{Q})$, which leads the contradiction. Thus, $H_B^\infty(\mathcal{Q})$ is dense in $H_B^1(\mathcal{Q})$.

As is easily seen, proof of the fact that $C_B^\infty(\bar{\mathcal{Q}})$ is dense in $H_B^1(\mathcal{Q})$ can be reduced to verifying that for any $v \in H_B^\infty(\mathcal{Q})$ $|(1 - \chi(\cdot/R))v|_{1,\mathcal{Q}} \rightarrow 0$ as $R \rightarrow \infty$. This follows from the inequality:

$$|(1 - \chi(\cdot/R))v|_{1,\mathcal{Q}}^2 \leq \int_{R \leq |x|} |\partial_x^1 v|^2 dx + (\sup_{R \leq |x| \leq 2R} |\nabla_x \chi|^2) R^{-2} \int_{R \leq |x| \leq 2R} |v|^2 dx,$$

since v satisfies the condition (1.1).

Next, let us prove that $H_B^\infty(\mathcal{Q})$ is dense in $H_B^2(\mathcal{Q})$. Suppose to the contrary that there exists a non-trivial $u \in H_B^2(\mathcal{Q})$ such that

$$(1.38) \quad \langle u, v \rangle_{2,\mathcal{Q}} = 0 \quad \text{for any } v \in H_B^\infty(\mathcal{Q}).$$

Introducing the innerproduct $\langle w, v \rangle_{1,\mathcal{Q}} + (w, v)_{L^2(\mathcal{Q})}$ instead of $\langle w, v \rangle_{1,\mathcal{Q}}$, in the similar manner to the proof of Theorem 1.7 we can prove that for any $w \in C_0^\infty(\mathcal{Q})$ there exists a $v \in H_B^\infty(\mathcal{Q})$ satisfying

$$(1.39) \quad -Lv + v = w \quad \text{in } \Omega, \quad Bv = 0 \quad \text{on } \partial\Omega.$$

In the same way as (1.37), (1.38) and (1.39) yield that $0 = -\langle u, v \rangle_{2, \Omega} = (Lu, w)_{L^2(\Omega)}$ for any $w \in C_0^\infty(\Omega)$. This implies that $Lu = 0$ in Ω . Hence, by the uniqueness of the solutions stated in Theorem 1.7 we have $u(x) = 0$ in Ω , which leads the contradiction. Hence $H_B^\infty(\Omega)$ is dense in $H_B^2(\Omega)$.

The fact that $C_B^\infty(\bar{\Omega})$ is dense in $H_B^2(\Omega)$ can be proved in the same way as in $H_B^1(\Omega)$. Thus, we obtain the lemma.

When proving that A is a skew self-adjoint operator, we need the following lemma which follows from Theorems 1.7 and 1.9.

Lemma 1.10. *Assume that (A.1) and (A.2) are valid. Let $v \in L^2(\Omega)$. If $u \in H_B^1(\Omega)$ satisfies*

$$(1.40) \quad \langle u, w \rangle_{1, \Omega} = (v, w)_{L^2(\Omega)} \quad \text{for any } w \in C_B^\infty(\bar{\Omega}),$$

then u belongs to $H_B^2(\Omega)$.

Proof. Put $\chi_R(x) = \chi(x/R)$ for $R > r_0$. Then, it is obvious that $\chi_R u \rightarrow u$ in the distribution sense as $R \rightarrow \infty$. Therefore, if $\{\chi_R u\}_{R > r_0}$ is a bounded set in the Hilbert space $H_B^2(\Omega)$, u has to belong to $H_B^2(\Omega)$. For, there exists a subsequence of $\{\chi_R u\}$ converging weakly in $H_B^2(\Omega)$ and this limit coincides with the limit in the distribution sense.

Let us prove the boundedness of $\{\chi_R u\}$. By (1.40) and integration by parts we have

$$(1.41) \quad \langle \chi_R u, w \rangle_{1, \Omega} = (\chi_R v + v'_R, w)_{L^2(\Omega)} \quad \text{for any } w \in C_B^\infty(\bar{\Omega}),$$

where $v'_R = -\sum \partial_i(a_{ij}(\partial_j \chi_R)u) - \sum a_{ij}(\partial_i \chi_R)\partial_j u$. It is easy to see that $\chi_R u \in H_B^1(\Omega)$ and $\chi_R v + v'_R \in L^2_{2R}(\Omega)$ for any $R > r_0$. Since $C_B^\infty(\bar{\Omega})$ is dense in $H_B^1(\Omega)$ (cf. Theorem 1.9), (1.41) is valid for any $w \in H_B^1(\Omega)$. Therefore, from the uniqueness of the solution of (1.30) and Theorem 1.7 it follows that $\chi_R u \in H_B^2(\Omega)$ and $-L(\chi_R u) = \chi_R v + v'_R$ for any $R > r_0$. By Lemma 1.8 we have

$$\begin{aligned} |\chi_R u|_{2, \Omega}^2 &\leq C_1 (\|\chi_R v + v'_R\|_{L^2(\Omega)}^2 + |\chi_R u|_{1, \Omega}^2) \\ &\leq C_2 (\|v\|_{L^2(\Omega)}^2 + \|u\|_{1, \Omega}^2 + R^{-2} \int_{R \leq |x| \leq 2R} |u(x)|^2 dx), \end{aligned}$$

where the constants C_1 and C_2 are independent of $R > r_0$. Hence, by (1.1) we obtain the boundedness of $\{\chi_R u\}$ in $H_B^2(\Omega)$. The proof is complete.

From now on, we return to analysis of the operator A (defined in (0.5)) and the equation (0.6). Put

$$H = H_B^1(\mathcal{Q}) \times L^2(\mathcal{Q}) \text{ and } (f, g)_{E, \mathcal{Q}} = \frac{1}{2} \{ \langle f_1, g_1 \rangle_{1, \mathcal{Q}} + (f_2, g_2)_{L^2(\mathcal{Q})} \}$$

for any $f = (f_1, f_2)$ and $g = (g_1, g_2) \in H$. As the innerproduct of H , we adopt $(\cdot, \cdot)_{E, \mathcal{Q}}$. It is easily seen that H is a Hilbert space. We introduce the following space as the domain of A :

$$D(A) = H_B^2(\mathcal{Q}) \times \{ H_B^1(\mathcal{Q}) \cap L^2(\mathcal{Q}) \},$$

whose innerproduct is defined by

$$(f, g)_{D(A), \mathcal{Q}} = \frac{1}{2} \{ \langle f_1, g_1 \rangle_{2, \mathcal{Q}} + \langle f_2, g_2 \rangle_{1, \mathcal{Q}} + (f_2, g_2)_{L^2(\mathcal{Q})} \}$$

for $f = (f_1, f_2)$ and $g = (g_1, g_2) \in D(A)$. For the notational convenience, we define Af by $Af = (f_2, Lf_1)$ for $f = (f_1, f_2) \in D(A)$. It is obvious that $Af \in H$ if $f \in D(A)$. $D(A)$ is a Hilbert space, and then, noting that $(f, g)_{D(A), \mathcal{Q}} = (f, g)_{E, \mathcal{Q}} + (Af, Ag)_{E, \mathcal{Q}}$, we see that A is a closed operator on H . Furthermore, it is seen from Theorem 1.9 that A is densely defined in H .

Under these preparations, we shall prove the main result in this section.

Theorem 1.11. *Assume that (A.1) and (A.2) are valid. Then, A is a skew self-adjoint operator in H with domain $D(A)$.*

Proof. Let us recall the definition of the adjoint operator A^* of A and the domain $D(A^*)$:

$$(1.42) \quad D(A^*) = \{ g \in H \mid \text{there exists an } h \in H \text{ such that} \\ (g, Af)_{E, \mathcal{Q}} = (h, f)_{E, \mathcal{Q}} \text{ for any } f \in D(A) \},$$

and then for $g \in D(A^*)$ A^*g is defined by $A^*g = h$. Since integration by parts yields that $(g, Af)_{E, \mathcal{Q}} = (-Ag, f)_{E, \mathcal{Q}}$ for any $f, g \in D(A)$, we know that $D(A) \subset D(A^*)$. Our task is only to prove that $D(A^*) \subset D(A)$.

Let $g = (g_1, g_2)$ and $h = (h_1, h_2)$ be the elements in (1.42). Then, choosing $f = (f_1, f_2)$ so that $f_1 \in H_B^\infty(\mathcal{Q})$ and $f_2 = 0$, by (1.42) we have

$$(1.43) \quad \langle h_1, f_1 \rangle_{1, \mathcal{Q}} = (g_2, Lf_1)_{L^2(\mathcal{Q})} \quad \text{for any } f_1 \in H_B^\infty(\mathcal{Q}).$$

From Theorem 1.7 it is seen that for any $v \in C_0^\infty(\mathcal{Q})$ there exists an $f_1 \in H_B^\infty(\mathcal{Q})$ such that

$$(1.44) \quad Lf_1 = v \quad \text{in } \mathcal{Q}.$$

Employing the same arguments as for (1.37), we have $\langle h_1, f_1 \rangle_{1, \mathcal{Q}} = -(h_1, Lf_1)_{L^2(\mathcal{Q})}$.

Combining this, (1.43) and (1.44) implies that

$$(g_2 + h_1, v)_{L^2(\Omega)} = 0 \quad \text{for any } v \in C_0^\infty(\Omega),$$

which shows that $g_2 = -h_1 \in H_B^1(\Omega) \cap L^2(\Omega)$.

Next, choosing $f = (f_1, f_2)$ so that $f_1 = 0$ and $f_2 \in C_B^\infty(\bar{\Omega})$, we have

$$\langle g_1, f_2 \rangle_{1, \Omega} = (h_2, f_2)_{L^2(\Omega)} \quad \text{for any } f_2 \in C_B^\infty(\bar{\Omega}).$$

Since $g_1 \in H_B^1(\Omega)$ and $h_2 \in L^2(\Omega)$, by Lemma 1.10 we see that $g_1 \in H_B^2(\Omega)$. Accordingly, we have proved that $g \in D(A)$, which completes the proof.

By Theorem 1.11 and Stone's theorem (cf. Appendix I in Lax and Phillips [6]), we have

Theorem 1.12. *Assume that (A.1) and (A.2) are valid. Then, there exists a one parameter group $\{U(t)\}_{t \in \mathbf{R}}$ generated by A and having the following properties:*

- (i) $U(t)$ is a unitary operator from H to itself for any $t \in \mathbf{R}$.
- (ii) $U(t)f$ is an H -valued continuous function in $t \in \mathbf{R}$ for any $f \in H$.
- (iii) $U(t)f$ is an H -valued C^1 function in $t \in \mathbf{R}$ if and only if $f \in D(A)$.
- (iv) $U(t)$ is a unitary operator also from $D(A)$ to itself for any $t \in \mathbf{R}$.
- (v) When $f \in D(A)$, $\frac{d}{dt} U(t)f = AU(t) = U(t)Af$ for any $t \in \mathbf{R}$.

From this theorem we see that for any $f = (f_1, f_2) \in D(A)$ $u(t, x) = (U(t)f)_1(x)$ (= the first component of $U(t)f$) belongs to $\bigcap_{j=0}^2 C^j(\mathbf{R}; H_B^{2-j}(\Omega))$ ($H_B^0(\Omega) = L^2(\Omega)$) and satisfies (0.1).

§ 2. The Problem in the Free Space

In this section we consider the unperturbed problem under the assumptions (A.1), (A.2) and (A.3) stated in § 0:

$$(2.1) \quad \begin{cases} (\partial_t^2 - L^0)u(t, x) = 0 & \text{in } \mathbf{R} \times \mathbf{R}^n, \\ u(0, x) = f_1(x), \quad \partial_t u(0, x) = f_2(x) & \text{in } \mathbf{R}^n \end{cases}$$

where $L^0 = \sum_{i,j=1}^n a_{ij}^0 \partial_i \partial_j$. Obviously, for this problem we can obtain the same result as in § 1. We employ the same notations as in § 1, and particularly we denote the space H and the operators $U(t)$ for (2.1) by H_0 and $U_0(t)$, respectively.

By (A.1) $L^0(\xi)$ is a symmetric matrix. So, all eigenvalues of $L^0(\xi)$ are real. Let N denote the number of the distinct eigenvalues $\lambda_j(\xi)$ of $L^0(\xi)$. Then, by (A.3) N is independent of $\xi \in \mathbb{R}^n - \{0\}$ and it can be shown that the $\lambda_j(\xi)$ can be enumerated so as to form N distinct and analytic branches in the following way: $\lambda_1(\xi) < \lambda_2(\xi) < \dots < \lambda_N(\xi)$. Since $L^0(t\xi) = t^2 L^0(\xi)$, $t \in \mathbb{R}$, with the enumeration used above the eigenvalues $\lambda_j(\xi)$ are C^∞ functions homogeneous of order 2 in $\xi \in \mathbb{R}^n$. If we put $\varepsilon_{ip} = \frac{1}{2}(\xi_i \eta_p + \xi_p \eta_i)$ for $\xi = (\xi_1, \dots, \xi_n)$ and $\eta = (\eta_1, \dots, \eta_n) \in \mathbb{R}^n$, by (A.1) and (A.2) we see that

$$\eta L^0(\xi) \eta = \sum_{i,p,j,q=1}^n a_{ipjq}^0 \varepsilon_{iq} \varepsilon_{jp} \geq \delta_1 \sum_{i,p=1}^n |\varepsilon_{ip}|^2,$$

where a_{ipjq}^0 are the (p, q) elements of the matrices a_{ij}^0 . Since $\sum_{i,p=1}^n |\varepsilon_{ip}|^2 \geq \frac{1}{2} |\xi|^2 |\eta|^2$, we see that there exists a $\delta > 0$ such that

$$\eta L^0(\xi) \eta \geq \delta |\xi|^2 |\eta|^2 \quad \text{for any } \xi \text{ and } \eta \in \mathbb{R}^n.$$

From this we see that $\lambda_j(\xi) \geq \delta |\xi|^2$ for any $\xi \in \mathbb{R}^n$. Let $P_j(\xi)$ be the orthogonal projection into the eigenspace of each $\lambda_j(\xi)$, which an $n \times n$ matrix of C^∞ functions on $\mathbb{R}^n - \{0\}$ homogeneous of order 0.

The Radon transform $\tilde{g}(s, \omega)$ of $g(x) \in \mathcal{S}$ (\mathcal{S} is the Schwartz space of rapidly decreasing functions) is defined by

$$\tilde{g}(s, \omega) = \int_{x \cdot \omega = s} g(x) dS_x, \quad (s, \omega) \in \mathbb{R} \times S^{n-1},$$

where S^{n-1} denotes the $n-1$ dimensional unit sphere. We denote the Fourier transform of $g(x)$ ($x \in \mathbb{R}^n$) by $\mathcal{F}[g](\xi) = \hat{g}(\xi) (= \int e^{-ix \cdot \xi} g(x) dx, i = \sqrt{-1})$, and the one of $k(s)$ ($s \in \mathbb{R}$) by $F[k](\sigma) (= \int e^{-is\sigma} k(s) ds)$. Let their inverse transformations be denoted by \mathcal{F}^{-1} and F^{-1} , respectively. Then, it follows that

$$(2.2) \quad g(x) = 2^{-1}(2\pi)^{1-n} \int_{S^{n-1}} F^{-1}(|\sigma|^{n-1} F\tilde{g})(x \cdot \omega, \omega) d\omega,$$

$$(2.3) \quad \hat{g}(\sigma\omega) = F\tilde{g}(\sigma, \omega).$$

Set

$$\lambda_{\pm}(\sigma) = \begin{cases} \frac{1-i}{\sqrt{2}} \sigma^{1/2} & \text{for } \sigma \geq 0, \\ \pm \frac{1+i}{\sqrt{2}} |\sigma|^{1/2} & \text{for } \sigma < 0. \end{cases}$$

Then, we see that $\text{supp} [\lambda_+(D_s)k] \subset (-\infty, s_0]$ (resp. $\text{supp} [\lambda_-(D_s)k] \subset [s_0, \infty)$) provided that $\text{supp} [k] \subset (-\infty, s_0]$ (resp. $\subset [s_0, \infty)$) and that $(\lambda_{\pm}(D_s))^2 = -\partial_s$, where $\lambda_{\pm}(D_s)k = F^{-1}[\lambda_{\pm}(\sigma)Fk(\sigma)]$ (cf. Lemma 1.1 of Soga [13]). Setting

$$J_{\pm} = \begin{cases} (-\partial_s)^{(n-1)/2} & \text{for odd } n, \\ (-\partial_s)^{(n/2)-1}\lambda_{\pm}(D_s) & \text{for even } n, \end{cases}$$

for the data $f = (f_1, f_2) \in \mathcal{S}$ in (2.1) we define $T_{\bar{0}}^{\pm}f$ by

$$(2.4) \quad T_{\bar{0}}^{\pm}f(s, \omega) = \sum_{j=1}^N \lambda_j(\omega)^{1/4} P_j(\omega) (-\lambda_j(\omega)^{1/2} \partial_s J_{\pm} \tilde{f}_1 + J_{\pm} \tilde{f}_2) (\lambda_j(\omega)^{1/2} s, \omega).$$

Then, $f(x) \in \mathcal{S}$ is reconstructed with $T_{\bar{0}}^{\pm}f(s, \omega)$ as follows:

Theorem 2.1. *For any $f = (f_1, f_2) \in \mathcal{S}$ we have*

$$(2.5) \quad f_1(x) = -2^{-1}(2\pi)^{1-n} \int_{S^{n-1}} \sum_{j=1}^N \lambda_j(\omega)^{-n/4} P_j(\omega) (\partial_s^{-1} J_{\pm}^* T_{\bar{0}}^{\mp} f) (\lambda_j(\omega)^{-1/2} x \cdot \omega, \omega) d\omega,$$

$$(2.6) \quad f_2(x) = 2^{-1}(2\pi)^{1-n} \int_{S^{n-1}} \sum_{j=1}^N \lambda_j(\omega)^{-n/4} P_j(\omega) (J_{\pm}^* T_{\bar{0}}^{\pm} f) (\lambda_j(\omega)^{-1/2} x \cdot \omega, \omega) d\omega.$$

The above formulas are corresponding to Theorem 2.2 in Chapter IV of Lax and Phillips [6] (cf. Proposition 1.1 of Soga [13] also). Noting that $P_j(\omega)$ are the orthogonal projections, we can derive the above theorem from (2.2) in the same way as in [6] and [13].

For any $k(s, \omega) \in \mathcal{S}(\mathbb{R} \times S^{n-1})$ we define

$$(2.7) \quad (Q_{\bar{0}}^{\pm}k)(x) = 2^{-1}(2\pi)^{1-n} \int_{S^{n-1}} \sum_{j=1}^N \lambda_j(\omega)^{-n/4} P_j(\omega) (\partial_s^{-1} J_{\pm}^* k) (\lambda_j(\omega)^{-1/2} x \cdot \omega, \omega) d\omega,$$

$$(2.8) \quad Q^{\pm}k = (-Q_{\bar{0}}^{\pm}k, Q_{\bar{0}}^{\pm}(\partial_s k)).$$

Obviously $(Q^{\pm}k)(x)$ belong to $C^{\infty}(\mathbb{R}^n)$; Theorem 2.1 means that $Q^{\pm}T_{\bar{0}}^{\pm}f = f$ for any $f \in \mathcal{S}$. Furthermore, we obtain

Lemma 2.2. *Set $\mathring{S} = \{k(s, \omega) \in \mathcal{S}(\mathbb{R} \times S^{n-1}) \mid Fk(\sigma, \omega) = 0 \text{ in a neighborhood of } \sigma = 0\}$. Then, for any $k \in \mathring{S}$ we have $(Q^{\pm}k)(x) \in \mathcal{S}$ and $T_{\bar{0}}^{\pm}Q^{\pm}k = k$.*

Proof. For $k(s, \omega) \in \mathring{S}$ set

$$\mathring{k}(\sigma, \omega) = \sum_{j=1}^N \lambda_j(\omega)^{(2-n)/4} P_j(\omega) F[\partial_s^{-1} J_{\pm}^* k] (\lambda_j(\omega)^{1/2} \sigma, \omega).$$

Then $\mathring{k}(\sigma, \omega)$ belongs to $\mathcal{S}(\mathbb{R} \times S^{n-1})$ and is equal to 0 in a neighborhood of $\sigma = 0$. Furthermore, it follows that

$$(2.9) \quad (Q_{\bar{0}}^{\mp}k)(x) = 2^{-1} \mathcal{F}^{-1}[\{\mathring{k}(|\xi|, \xi/|\xi|) + \mathring{k}(-|\xi|, -\xi/|\xi|)\} |\xi|^{1-n}](x),$$

which yields that $(Q^\pm k)(x) \in \mathcal{S}$.

Combining (2.3), (2.9) and the definitions (2.7), (2.8) and noting that $P_j(\omega)$ are the orthogonal projections, we have $T_0^\pm Q^\pm k(s, \omega) = \sum_{j=1}^4 \Phi_j(s, \omega)$ where

$$\begin{aligned} \Phi_1(s, \omega) &= \frac{1}{2} \sum_{j=1}^N \lambda_j(\omega)^{3/4} P_j(\omega) F^{-1}[K_1(\cdot, \omega)](\lambda_j(\omega)^{1/2} s), \\ &K_1(\sigma, \omega) = \dot{k}(\sigma, \omega) |\sigma|^{1-n} i \sigma J_\pm(\sigma), \\ \Phi_2(s, \omega) &= \frac{1}{2} \sum_{j=1}^N \lambda_j(\omega)^{1/4} P_j(\omega) F^{-1}[K_2(\cdot, \omega)](\lambda_j(\omega)^{1/2} s), \\ &K_2(\sigma, \omega) = (\partial_s k^0)(\sigma, \omega) |\sigma|^{1-n} J_\pm(\sigma), \\ \Phi_3(s, \omega) &= \frac{1}{2} \sum_{j=1}^N \lambda_j(\omega)^{3/4} P_j(\omega) F^{-1}[K_3(\cdot, \omega)](\lambda_j(\omega)^{1/2} s), \\ &K_3(\sigma, \omega) = \dot{k}(-\sigma, -\omega) |\sigma|^{1-n} i \sigma J_\pm(\sigma), \\ \Phi_4(s, \omega) &= \frac{1}{2} \sum_{j=1}^N \lambda_j(\omega)^{1/4} P_j(\omega) F^{-1}[K_4(\cdot, \omega)](\lambda_j(\omega)^{1/2} s), \\ &K_4(\sigma, \omega) = \partial_s k^0(-\sigma, -\omega) |\sigma|^{1-n} i \sigma J_\pm(\sigma). \end{aligned}$$

Noting that $F[\partial_s^{-1} J_\pm^* k](\sigma, \omega) = (i\sigma)^{-1} \overline{J_\pm(\sigma)} [Fk](\sigma, \omega)$ and that $J_\pm(\sigma) \overline{J_\pm(\sigma)} = |\sigma|^{n-1}$, by the change of the variable: $\lambda_j(\omega)^{1/2} \sigma \rightarrow \sigma$, we have

$$F^{-1}[K_l(\cdot, \omega)](\lambda_j(\omega)^{1/2} s) = \sum_{l=1}^N \lambda_l(\omega)^{-3/4} P_l(\omega) k(\lambda_l(\omega)^{-1/2} \lambda_j(\omega)^{1/2} s, \omega).$$

Using the fact that the $P_j(\omega)$ are the orthogonal projections again, it follows that $\Phi_1(s, \omega) = 2^{-1} k(s, \omega)$. In the same way, we also obtain $\Phi_2 = \Phi_1$ and $\Phi_3 = -\Phi_4$. These prove that $(T_0^\pm Q^\pm k)(s, \omega) = k(s, \omega)$, which completes the proof of the lemma.

Lax and Phillips [6] say that a unitary operator T_0^+ (resp. T_0^-) from H_0 to $L^2(\mathbf{R} \times S^{n-1})$ is the outgoing (resp. incoming) translation representation if there exists a closed subspace D_0^+ (resp. D_0^-) in H_0 such that

$$(2.10) \quad T_0^\pm \text{ map } D_0^\pm \text{ onto } L_\pm^2(\mathbf{R} \times S^{n-1}) = \{k(s, \omega) \in L^2(\mathbf{R} \times S^{n-1}) \mid k(s, \omega) = 0 \text{ for } \pm s < 0\},$$

$$(2.11) \quad U_0^\pm(t) D_0^\pm \subset D_0^\pm \quad \text{for } \pm t > 0,$$

$$(2.12) \quad \bigcap_{t \in \mathbf{R}} U_0(t) D_0^\pm = \{0\},$$

$$(2.13) \quad \overline{\bigcup_{t \in \mathbf{R}} U_0(t) D_0^\pm} = H_0,$$

and T_0^\pm satisfy

$$(2.14) \quad T_0^\pm U_0(t) = \mathcal{Q}_t T_0^\pm,$$

where $(\mathcal{Q}_t k)(s) = k(s-t)$. And, D_0^+ (resp. D_0^-) is called the outgoing (resp. incoming) subspace. The operator T_0^+ (resp. T_0^-) defined by (2.4) becomes the outgoing (resp. incoming) translation representation:

Theorem 2.3. (i) *The operators T_0^\pm (in (2.4)) can be extended as unitary operators from H_0 to $L^2(\mathbf{R} \times S^{n-1})$.*

(ii) *Let us denote these extensions of T_0^\pm also by T_0^\pm . Set $D_0^\pm = \{f \in H_0 \mid T_0^\pm f(s, \omega) = 0 \text{ when } \pm s < 0\}$. Then, T_0^+ (resp. T_0^-) is the outgoing (resp. incoming) translation representation, and D_0^+ (resp. D_0^-) is the outgoing (resp. incoming) subspace (i.e., all the conditions (2.10)–(2.14) are satisfied).*

Proof. (i) Noting that $P_j(\omega)$ are the orthogonal projections, for any $f = (f_1, f_2) \in \mathcal{S}$ we have

$$\begin{aligned} \|T_0^\pm f\|_{L^2(\mathbf{R} \times S^{n-1})}^2 &= \iint_{\mathbf{R} \times S^{n-1}} \sum_{j=1}^N |\lambda_j^{1/2} P_j \partial_s J_\pm \tilde{f}_1|^2 ds d\omega \\ &+ \iint_{\mathbf{R} \times S^{n-1}} \sum_{j=1}^N |P_j J_\pm \tilde{f}_2|^2 ds d\omega - 2 \operatorname{Re} \iint_{\mathbf{R} \times S^{n-1}} \sum_{j=1}^N \lambda_j^{1/2} P_j \partial_s J_\pm \tilde{f}_1 \cdot \overline{P_j J_\pm \tilde{f}_2} ds d\omega. \end{aligned}$$

Therefore, using the facts that $P_j(\omega) = P_j(-\omega)$, $F\tilde{g}(\sigma, \omega) = F\tilde{g}(-\sigma, -\omega)$ and that $\sum_{j=1}^N \sigma^2 \lambda_j(\omega) P_j(\omega) = L^0(\sigma\omega)$, we have

$$\|T_0^\pm f\|_{L^2(\mathbf{R} \times S^{n-1})}^2 = 4(2\pi)^{n-1} \|f\|_{E, \mathbf{R}^n}^2 \quad \text{for any } f \in \mathcal{S}.$$

Let us take $2^{-1}(2\pi)^{(1-n)/2} \|k\|_{L^2(\mathbf{R} \times S^{n-1})}$ as the norm of $k \in L^2(\mathbf{R} \times S^{n-1})$. Then, T_0^\pm can be regarded as isometric operators from H_0 to $L^2(\mathbf{R} \times S^{n-1})$, because $\mathcal{S} (\supset C_0^\infty(\mathbf{R}^n))$ is dense in H_0 (cf. Theorem 1.9). Lemma 2.2 yields that $Q^\pm \mathring{S} \subset \mathcal{S}$ and $\mathring{S} \subset T_0^\pm \mathcal{S}$. Hence, the (extended) operators T_0^\pm map H_0 onto $L^2(\mathbf{R} \times S^{n-1})$ since \mathring{S} is dense in $L^2(\mathbf{R} \times S^{n-1})$. Thus, (i) of the theorem is proved.

(ii) From (i) of the theorem we see easily that the spaces D_0^\pm are closed subspaces in H_0 and mapped onto $L_\pm^2(\mathbf{R} \times S^{n-1})$ by T_0^\pm (i.e., (2.10) is satisfied).

Let us check (2.14). Assume that the data $f = (f_1, f_2)$ in (2.1) belong to $C_0^\infty(\mathbf{R}^n)$. Then, from finiteness of the propagation speed (cf. Theorem 3.1 in § 3 below), it is seen that the solution $u(t, x)$ is a $C_0^\infty(\mathbf{R}_x^n)$ -valued C^∞ function on \mathbf{R}_t . Therefore, for any fixed t the concrete definition form of the Radon transform $\tilde{u}(t, s, \omega)$ is valid. Noting that $\tilde{\partial}_i u(t, s, \omega) = \omega_i \partial_s \tilde{u}(t, s, \omega)$ and $L^0(\omega) = \sum_{j=1}^N \lambda_j(\omega) P_j(\omega)$ ($\omega \in S^{n-1}$), we have $(\partial_t^2 u - L^0 u)^\sim(s, \omega) = \sum_{j=1}^N (\partial_t^2 - \lambda_j(\omega)) \partial_s^2$

$P_j(\omega)\tilde{u}(t, s, \omega)$. If we put $v_j(t, s, \omega) = (\partial_t - \lambda_j(\omega)^{1/2} \partial_s)P_j(\omega)\tilde{u}(t, s, \omega)$, then $(\partial_t + \lambda_j(\omega)^{1/2} \partial_s)v_j(t, s, \omega) = 0$ for any $(t, s, \omega) \in \mathbb{R} \times \mathbb{R} \times S^{n-1}$ and $j = 1, \dots, N$, which yields that $v_j(t, s, \omega) = v_j(0, s - \lambda_j(\omega)^{1/2}t, \omega)$. Combining this and the fact that $P_j(\omega)(T_0^\pm U_0(t)f)(s, \omega) = \lambda_j(\omega)^{1/4}(J_\pm v_j)(t, \lambda_j(\omega)^{1/2}s, \omega)$, we obtain $P_j(\omega) \cdot (T_0^\pm U_0(t)f)(s, \omega) = P_j(\omega)(T_0^\pm f)(s-t, \omega)$ for any $(s, t, \omega) \in \mathbb{R} \times \mathbb{R} \times S^{n-1}$ and $j = 1, \dots, N$, which proves (2.14) (note that $C_0^\infty(\mathbb{R}^n)$ is dense in H_0).

Finally we have only to verify (2.11)–(2.13). From what we have proved, we see that (2.11)–(2.13) mean that $\mathcal{G}_t L_\pm^2(\mathbb{R} \times S^{n-1}) \subset L_\pm^2(\mathbb{R} \times S^{n-1}) (\pm t > 0)$, $\cap \mathcal{G}_t L_\pm^2(\mathbb{R} \times S^{n-1}) = 0$ and $\overline{\cup \mathcal{G}_t L_\pm^2(\mathbb{R} \times S^{n-1})} = L^2(\mathbb{R} \times S^{n-1})$, respectively. These are obvious. The proof is complete.

D_0^\pm are characterized in the following theorem. This characterization plays a fundamental role when defining the translation representations for the mixed problem (0.1) in § 3 below.

Theorem 2.4. *f belongs to D_0^\pm if and only if*

$$(2.15) \quad \text{supp } [U_0(t)f] \subset \{(t, x) \mid \pm \eta t \leq |x|\},$$

where $\eta = \min \{\lambda_j(\omega)^{1/2} \mid \omega \in S^{n-1}, j = 1, \dots, N\}$.

Proof. The idea of the proof is due to Lax and Phillips [6, 7] (cf. Theorem 1.2 in Chapter VI of [6] and Corollary 4.2 of [7]). We shall give only an outline of the proof.

Let $f \in D_0^\pm$. Then there exists a sequence $\{k_\pm^i\}_{i=1,2,\dots}$ in $C_0^\infty(\mathbb{R} \times S^{n-1})$ such that $\text{supp } [k_\pm^i] \subset \mathbb{R}_\pm \times S^{n-1}$ ($\mathbb{R}_\pm = \{s \in \mathbb{R} \mid \pm s > 0\}$) and $\|k_\pm^i - T_0^\pm f\|_{L^2(\mathbb{R} \times S^{n-1})} \rightarrow 0$ as $i \rightarrow \infty$. In the same way as in Soga [13] (see Corollary 1.1 of [13]), we see that the formulas (2.5) and (2.6) are valid for $f = f_\pm^i = (T_0^\pm)^{-1}k_\pm^i$. Therefore, combining (2.5), (2.6) and (2.14), we have

$$(U_0(t)f_\pm^i)_l(x) = \frac{1}{2} (2\pi)^{1-n} \sum_{j=1}^N \int \lambda_j(\omega)^{-n/4} P_j(\omega) (\partial_s^{l-2} J_\pm^* k_\pm^i)(\lambda_j(\omega)^{-1/2}x \cdot \omega - t, \omega) d\omega$$

for $l = 1$ and 2 . From these formulas and the fact that $\text{supp } [\partial_s^{l-2} J_\pm^* k_\pm^i] \subset \overline{\mathbb{R}_\pm} \times S^{n-1}$, it follows that $\text{supp } [U_0(t)f_\pm^i] \subset \{(t, x) \in \mathbb{R}^{n+1} \mid \pm \eta t \leq |x|\}$ for $\pm t > 0$. On the other hand, for any fixed t $U_0(t)f_\pm^i$ converge to $U_0(t)f$ in $L_{\text{loc}}^2(\mathbb{R}^n)$ as $i \rightarrow \infty$. Hence, we obtain (2.15).

Conversely, let (2.15) be satisfied. We can assume without loss of generality that $f \in D(A_0^\infty) \cap \bigcap_{m=1}^\infty D_0^\infty(A_0^m)$ (the operator A is denoted by A_0 here), which means that $\partial_s^m T_0^\pm f(s, \omega) \in L^2(\mathbb{R} \times S^{n-1})$ for any $m \geq 1$. Therefore, $J_\pm^* T_0^\pm f(s, \omega)$ is $L^2(S_\omega^{n-1})$ -valued (consequently $L^1(S_\omega^{n-1})$ -valued) C^∞ functions in

$s \in \mathbf{R}$. The proof can be reduced to verifying that

$$(2.16) \quad (\partial_x^\alpha U_0(t)f)_2(0) = \frac{1}{2}(2\pi)^{1-n}(-\partial_t)^{|\alpha|} \int_{S^{n-1}} \sum_{j=1}^N (\lambda_j(\omega)^{-1/2} \omega)^\alpha \cdot \lambda_j(\omega)^{-n/4} P_j(\omega)(J_\pm^* T_0^\pm f)(-t, \omega) d\omega$$

for any multi-index α (note that $U_0(t)f(x) \in C^\infty$, which follows from $f \in D(A_0^\infty)$). In fact, since $\partial_x^\alpha U_0(t)f(0) = 0$ for any $\pm t > 0$, by (2.16) we have

$$\int_{S^{n-1}} \sum_{j=1}^N (\lambda_j(\omega)^{-1/2} \omega)^\alpha \lambda_j(\omega)^{-n/4} P_j(\omega)(J_\pm^* T_0^\pm f)(-t, \omega) d\omega = 0$$

for any α and $\pm t > 0$. Combining this fact and Lemma 1.1 in Chapter VI of Lax and Phillips [6] yields that $P_j(\omega)(J_\pm^* T_0^\pm f)(s, \omega) = 0$ for any $\pm s < 0, \omega \in S^{n-1}$ and $j = 1, \dots, N$, which proves that $f \in D_0^\pm$.

(2.16) can be verified as follows. Combining (2.6) and (2.14), we can easily obtain (2.16) if $f \in \mathcal{S}$. For $f \in D(A_0^\infty)$, take a sequence $\{f^i\}_{i=1,2,\dots} \subset \mathcal{S}$ such that $\sum_{k=0}^m \|A_0^k(f^i - f)\|_{E, \mathbf{R}^n} \rightarrow 0$ as $i \rightarrow \infty$ ($m > |\alpha| + n/2$). Then, it is seen that for any fixed $t U_0(t)f^i \rightarrow U_0(t)f$ in $H_{loc}^m(\mathbf{R}^n)$ as $i \rightarrow \infty$, and that $T_0^\pm f^i(s, \omega) \rightarrow T_0^\pm f(s, \omega)$ in $H^m(\mathbf{R}_s; L^2(S_0^{n-1}))$ as $i \rightarrow \infty$. These facts imply that (2.16) is valid for any $f \in D(A_0^\infty)$.

§ 3. The Problem in the Exterior Domain

Throughout this section we use the notations in §§ 1 and 2. It is expected that the problem (0.1) has finite propagation speed: There exists a constant $\mu (> 0)$ such that $u(t_0, x_0) = 0$ when $u(T, x) = 0$ for any $x \in \{y \in \mathcal{Q} \mid |y - x_0| < \mu(t_0 - T)\}$ ($T < t_0$). This is derived from the following local energy estimate.

Theorem 3.1. *Let the assumptions (A.1) and (A.2) be satisfied. Then, there exists a constant $\mu (> 0)$ depending only on $n, \sup_{x \in \mathcal{Q}} |a_{ij}(x)|$ and the constant δ in (A.2) such that if $u(t, x) \in C^2(\mathbf{R} \times \bar{\mathcal{Q}})$ is a solution of (0.1), for any $T > 0, R > 0$ and $y \in \bar{\mathcal{Q}} U(t)f (= (u(t, x), (\partial_t u(t, x))))$ satisfies*

$$(3.1) \quad \|U(T)f\|_{E, \mathcal{Q}_R(y)} \leq \|f\|_{E, \mathcal{Q}_{R+\mu T}(y)},$$

where $\mathcal{Q}_R(y) = \{x \in \mathcal{Q} \mid |x - y| < R\}$ and

$$\|f\|_{E, \mathcal{G}}^2 = \frac{1}{2} \int_{\mathcal{G}} \left\{ \sum_{i,j,p,q=1}^n a_{ipjq}(x) \partial_j f_{1q}(x) \overline{\partial_i f_{1p}(x)} + |f_2(x)|^2 \right\} dx$$

($f = (f_1, f_2)$ and $f_i = {}^t(f_{i1}, \dots, f_{in}), l = 1, 2$).

Proof. The idea of the proof is the same as in the proof of Theorem 1.1 in Chapter V of Lax and Phillips [6]. Multiply the equation $(\partial_t^2 - L)u = 0$ by $\overline{\partial_t u}$ and integrate over $D = \{(t, x) | 0 < t < T, |x - y| < R + (T - t)\mu, x \in \Omega\}$. Then, by integration by parts we have

$$\begin{aligned} & \|U(T)f\|_{E, \Omega_R(y)}^2 - \|U(0)f\|_{E, \Omega_{R+\mu T}(y)}^2 + \int_{((0, T) \times \partial\Omega) \cap D} (\partial_t u \cdot \overline{Nu} + Nu \cdot \overline{\partial_t u}) dt d(\partial\Omega) \\ & + (1 + \mu^2)^{-1/2} \int_{\Gamma} \{ \mu (|\partial_t u|^2 + \sum_{i,j,p,q=1}^n a_{ipjq} \partial_j u_q \cdot \overline{\partial_t u_p}) \\ & \quad - 2 \operatorname{Re} \sum_{i,p,j,q=1}^n a_{ipjq} r_i \partial_j u_q \cdot \overline{\partial_t u_p} \} d\Gamma = 0, \end{aligned}$$

where $\Gamma = \{(t, x) | 0 < t < T, |x - y| = R + (T - t)\mu, x \in \Omega\}$, $r_i = (x_i - y_i) / |x - y|$ and $Nu = \sum_{i,j=1}^n \nu_i a_{ij} \partial_j u$. Therefore, if u satisfies (0.1), it follows from (A.1) and (A.2) that

$$\begin{aligned} (3.2) \quad & \|U(T)f\|_{E, \Omega_R(y)}^2 \leq \|f\|_{E, \Omega_{R+\mu T}(y)}^2 - (1 + \mu^2)^{-1/2} \int_{\Gamma} \{ \mu (|\partial_t u|^2 + \delta \sum_{i,p=1}^n |\varepsilon_{ip}(u)|^2) \\ & \quad - n^2 (\sup_{i,j=1, \dots, n} |a_{ij}|) (n |\partial_t u|^2 + \sum_{i,p=1}^n |\varepsilon_{ip}(u)|^2) \} d\Gamma. \end{aligned}$$

Take the μ so that $\mu \geq (\max \{n^3, n^2/\delta\}) \sup |a_{ij}|$. Then, the integral over Γ in (3.2) is non-negative, and consequently (3.1) is obtained. The proof is complete.

Hereafter, we assume that the coefficients of L in (0.1) are constant (i.e., $L = L^0$), and that (A.1), (A.2) and (A.3) in § 0 are satisfied. From now on, using the translation representations $T_{\pm}^{\#}$ in the free space, we make the translation representations for the perturbed equation (0.1) in the same way as in § 2 of Chapter V of Lax and Phillips [6]. Set

$$D_{\pm} = U_0(\pm r_0 \eta^{-1}) D_0^{\pm}.$$

Then, it follows from Theorem 2.4 that if $f \in D_{\pm}$, $U_0(t)f(x) = 0$ in a neighborhood of $\partial\Omega$ for any $\pm t > 0$. Therefore, D_{\pm} become closed subspaces in H , and we have

$$(3.3) \quad U_0(t)f = U(t)f \quad \text{for any } \pm t > 0 \text{ if } f \in D_{\pm}.$$

Combining this and Theorem 2.3, we obtain easily (i) and (ii) in the following theorem.

Theorem 3.2. D_+ (resp. D_-) is an outgoing (resp. incoming) subspace for $U(t)$:

- (i) $U(t)D_{\pm} \subset D_{\pm} \ (\pm t > 0),$
- (ii) $\bigcap_{t \in \mathbf{R}} U(t)D_{\pm} = \{0\},$
- (iii) $\overline{\bigcup_{t \in \mathbf{R}} U(t)D_{\pm}} = H.$

The (iii) of Theorem 3.2 is a key point, and is closely related to the local energy decay as was stated in § 2 of Chapter V of Lax and Phillips [6]. We shall prove (iii) later, and we use this theorem without the proof for a while. For any $f \in \bigcup_{t \in \mathbf{R}} U(t)D_{\pm}$ we set

$$T_{\pm}f = \lim_{\tau \rightarrow \infty} \mathcal{G}_{\mp\tau\mp r_0\eta}^{-1} T_0^{\pm} U(\tau)f.$$

Then, T_{\pm} are densely defined in H (from (iii) of Theorem 3.2 and (3.3)), and can be regarded as isometric operators from H to $L^2(\mathbf{R} \times S^{n-1})$ (by Theorem 2.3). Furthermore, we obtain

Corollary 3.3. T_+ (resp. T_-) is the outgoing (resp. incoming) translation representation for $U(t)$ with the outgoing (resp. incoming) subspace D_+ (resp. D^-).

In the same way as in Remark 2.2 in Chapter V of Lax and Phillips [6], we see from (iii) of Theorem 3.2 also that the wave operators:

$$W_{\pm} = \lim_{t \rightarrow \pm\infty} U(-t)U_0(t)$$

are well-defined and complete. Namely, we have

Corollary 3.4. W_{\pm} are unitary operators from H_0 to H .

Note that the scattering operator $S = (T_+)^{-1}T_-$ is well-defined and unitary from $L^2(\mathbf{R} \times S^{n-1})$ to itself.

Now, we shall give a proof of Theorem 3.2. As we mentioned earlier, it suffices to prove only (iii) of Theorem 3.2. Lax and Phillips in § 2 of Chapter V of [6] showed that (iii) could be derived from the non-existence of the point spectrum of A and the local energy estimate stated in Theorem 3.1. Thus, the proof is complete if the following theorem is verified.

Theorem 3.5. A has no point spectrum. Namely, if $f \in D(A)$ and $Af = \tau f$ for a $\tau \in \mathbf{C}$ (= the field of complex numbers), then $f = 0$.

To prove this theorem, we need the following lemma.

Lemma 3.6. Let $\sigma \in \mathbf{R} - \{0\}$. If $u \in L^2(\mathbf{R}^n)$ satisfies: $(L^0 + \sigma^2 I)u(x) = 0$ for $|x| > R$ where I is the $n \times n$ identity matrix, then $u(x) = 0$ for $|x| > R$.

The proof of this lemma will be reduced to the following theorem due to Hörmander [4], Littman [8] and Murata [10].

Theorem 3.7. *Let the polynomial $P(\xi)$ in $\xi \in \mathbb{R}^n$ be decomposed into the form: $c \prod_{j=1}^l P_j(\xi)^{\beta_j}$ ($c \in \mathbb{C}$ and β_j are positive integers), where*

- (i) $P_j(\xi)$ ($j = 1, \dots, l$) are irreducible and real polynomials,
 - (ii) there exists a $\xi^j \in \mathbb{R}^n$ for each P_j such that $P_j(\xi^j) = 0$ and $\text{grad } P_j(\xi^j) \neq 0$.
- If $\phi \in L^2_{\text{loc}}(\mathbb{R}^n) \cap S'$ (S' is the dual space of S) satisfies

$$(3.4) \quad \lim_{R \rightarrow \infty} R^{-1} \int_{R \leq |x| \leq 2R} |\phi(x)|^2 dx = 0$$

and $\text{supp } [P(D_x)\phi]$ is compact $\left(D_x = \frac{1}{i} (\partial_1, \dots, \partial_n) \right)$, then $\text{supp } [\phi]$ is compact; more precisely, $\text{supp } [\phi] \subset$ the convex hull of $\text{supp } [P(D_x)\phi]$.

Proof of Lemma 3.6. Let $u = {}^t(u_1, \dots, u_n)$ and set $P(\xi) =$ the determinant of the matrix $\sigma^2 I - L^0(\xi)$. Multiplying $L^0 + \sigma^2 I$ by its cofactor, we see that $P(D_x)u_j(x) = 0$ for $|x| > R$. Therefore, if we check that $P(\xi)$ has all the properties stated in Theorem 3.7, the lemma follows from this theorem, since $\phi = u_j \in L^2(\mathbb{R}^n)$ satisfies (3.4).

Let us check that $P(\xi)$ satisfies (i) and (ii) of Theorem 3.7. Put $\xi = (\xi', \xi_n) \in \mathbb{R}^{n-1} \times \mathbb{R}$. For $j = 1, \dots, N$ we have $\sigma^2 - \lambda_j(0, \xi_n^j) = 0$ and $(\partial \lambda_j / \partial \xi_n)(0, \xi_n^j) \neq 0$ where $\xi_n^j = |\sigma \lambda_j(0, 1)|^{-1/2}$ by the Euler identity. Hence, there exist positive-valued C^∞ functions $\eta^j(\xi')$ defined in a neighborhood ω of $\xi' = 0$ such that $\sigma^2 - \lambda_j(\xi', \eta^j(\xi')) = 0$ for any $\xi' \in \omega$ and $\eta^j(0) = \xi_n^j$. It is seen that $P(\xi)$ is of the form:

$$P(\xi', \xi_n) = (\det a_{nn}^0) \prod_{j=1}^N (\xi_n - \eta^j(\xi'))^{\alpha_j} (\xi_n + \eta^j(\xi'))^{\alpha_j},$$

where α_j are the multiplicities of $\lambda_j(\xi)$ and $\det a_{nn}^0$ denotes the determinant of the matrix a_{nn}^0 . Note that $\det a_{nn}^0 > 0$, which follows from (A.1) and (A.2). Decompose $P(\xi)$ into the product of the irreducible polynomials $P_j(\xi)$, $j = 1, \dots, l$, i.e., $P(\xi) = \prod_{j=1}^l P_j(\xi)^{\beta_j}$. Then, the set $\{\mu_k^j(\xi') \mid k = 1, \dots, m_j, j = 1, \dots, l\}$ of the roots of the equations: $P_j(\xi', \xi_n) = 0$ in ξ_n ($j = 1, \dots, l$) coincides with the set $\{\pm \eta^j(\xi') \mid j = 1, \dots, N\}$ for any ξ' in an open set $\subset \omega$. From this fact we can see that every $P_j(\xi)$ satisfies (i) and (ii) in Theorem 3.7, which completes the proof of the lemma.

Proof of Theorem 3.5. Since $0 = (Af, f)_{E, \Omega} + (f, Af)_{E, \Omega} = 2 \text{Re } \tau \|f\|_{E, \Omega}^2$ for

any $f \in D(A)$ satisfying: $Af = \tau f$, the theorem is trivial when $\operatorname{Re} \tau \neq 0$. Let $\tau = i\sigma$ ($\sigma \in \mathbb{R}$). Then, from the equation: $Af = i\sigma f$ ($f = (f_1, f_2) \in D(A)$) we have

$$(3.5) \quad (L^0 + \sigma^2 I)f_1(x) = 0 \quad \text{in } \mathcal{Q},$$

$$(3.6) \quad f_2 = i\sigma f_1.$$

If $\sigma = 0$, then $f_2 = 0$ and $L^0 f_1 = 0$ in \mathcal{Q} . The uniqueness of the solutions (cf. Theorem 1.7) implies that $f_1 = 0$, and then $f = 0$. If $\sigma \neq 0$, it follows from (3.5) that $f_1 = -\sigma^{-2} L^0 f_1 \in L^2(\mathcal{Q})$. Thus, if we put $u(x) = \chi(x/r_0)f_1(x)$ for $x \in \mathcal{Q}$ and $= 0$ for $x \notin \mathcal{Q}$, then $u \in L^2(\mathbb{R}^n)$ and $(L^0 + \sigma^2 I)u(x) = 0$ for $|x| > 2r_0$. Hence, by Lemma 3.6 we have $u(x) = 0$ for $|x| > 2r_0$, which means that

$$(3.7) \quad f_1(x) = 0 \quad \text{for } |x| > 2r_0.$$

Note that $L^0 + \sigma^2 I$ is strongly elliptic. Thus, the analytic-hypoellipticity implies that $f_1(x)$ is analytic in \mathcal{Q} . It follows from (3.7) that $f_1(x) = 0$ for all $x \in \mathcal{Q}$. Consequently, we have $f = 0$, which completes the proof.

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