# On a Local Energy Decay of Solutions of a Dissipative Wave Equation

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## §1. Introduction.

This study is concerned with a local energy decay property of solutions to the following initial boundary value problem of the dissipative wave equation:

$$(D) \begin{cases} u_{tt} + u_t - \Delta u = 0 & \text{in } \Omega \text{ and } t > 0, \\ u = 0 & \text{on } \Gamma \text{ and } t > 0, \\ u(0, x) = u_0(x), \ u_t(0, x) = u_1(x) & \text{in } \Omega, \end{cases}$$

where  $\Omega$  is an exterior domain in an n-dimensional Euclidean space  $\mathbb{R}^n$ , whose boundary  $\Gamma$  is a  $C^{\infty}$  and compact hypersurface. Below,  $r_0 > 0$  is a fixed constant such that  $\Omega^c \subset B_{r_0} = \{x \in \mathbb{R}^n \mid |x| < R\}$ . ( $\Omega^c$  is the complement of  $\Omega$ .)

In the wave equation case, the local energy decays exponentially fast if n is odd and polynomially fast if n is even, when  $\Omega$  is at least non-trapping (cf. [9], [10], [11], [16]). In fact, from a physical point of view the energy propagates along the wave fronts, so that the motion stops after time passes unless the wave front is trapped in a bounded set.

In the dissipative wave case, the energy also propagates along the wave front. Moreover, the trapped energy also decreases in virtue of the dissipative term  $u_t$ , so that we can expect to get the local energy decay result for any domains. In fact, in 1983 Shibata [14] proved the following theorem.

**Theorem 1.1.** Assume that  $n \geq 3$ . Let  $R > r_0$  and let u(t, x) be a smooth solution of (D) such that suppu(0, x),  $suppu_t(0, x) \subset \Omega_R = \{x \in \Omega \mid |x| < R\}$ . Then, there exists

a constant C > 0 depending on n and R such that

$$\int_{\Omega_R} \left\{ |u_t(t,x)|^2 + \sum_{|\alpha| \le 1} |\partial_x u(t,x)|^2 \right\} dx$$

$$\le C(1+t)^{-n} \left\{ \sum_{|\alpha| \le 3} \int_{\Omega} |\partial_x^{\alpha} u_t(0,x)|^2 dx + \sum_{|\alpha| \le 4} \int_{\Omega} |\partial_x^{\alpha} u(0,x)|^2 dx \right\},$$

where 
$$\partial_x^{\alpha} v = \partial^{|\alpha|} v / \partial_{x_1}^{\alpha_1} \cdots \partial_{x_n}^{\alpha_n}$$
,  $\alpha = (\alpha_1, \dots, \alpha_n)$  and  $|\alpha| = \alpha_1 + \dots + \alpha_n$ .

The purpose of this study is to show the decay rate of the local energy of even the weak solutions of (D) is also n/2 when  $n \geq 2$ , that is, we shall prove the following theorem.

**Theorem 1.2.** Assume that  $n \geq 2$ . Let  $R > r_0$  and  $u_0 \in H^1_{0,R}(\Omega)$  and  $u_1 \in L^2_R(\Omega)$ , where

$$L_R^2(\Omega) = \{ f \in L^2(\Omega) \mid supp f \subset \Omega_R \},$$
 
$$H_{0,R}^1(\Omega) = \{ f \in H^1(\Omega) \mid supp f \subset \Omega_R, \ f = 0 \ on \ \Gamma \}.$$

Let u(t,x) be a weak solution of (D). Then, there exists a constant C depending on n and R such that

$$\int_{\Omega_R} \{ |u_t(t,x)|^2 + \sum_{|\alpha| \le 1} |\partial_x^{\alpha} u(t,x)|^2 \} dx 
\le C(1+t)^{-n} \{ \int_{\Omega} |u_1(x)|^2 dx + \sum_{|\alpha| \le 1} \int_{\Omega} |\partial_x^{\alpha} u_0(x)|^2 dx \}.$$

Compared with Theorem 1.1, in Theorem 1.2 we remove the smoothness assumption on solutions of (D) and we consider the case that n = 2 as well as the case that  $n \ge 3$ .

For the Cauchy problem of the dissipative wave equation (i.e.  $\Omega = \mathbb{R}^n$ ), A. Matsumura [8] studied the decay rate of solutions in 1976. His argument was based on the concrete representation of solutions by using the Fourier transform. When  $\Omega$  is bounded, it is well-known that the energy of solutions decays exponentially fast. In fact, this fact is easily proved by the multiplications of the equation with  $u_t$  and u

and by use of Poincaré's inequality. Since  $\Omega$  is unbounded in our case, we cannot use Poincaré's inequality. And also, because of the boundary, we can not use the Fourier transform. Our method is based on a spectral analysis to the corresponding stationary problem.

# §2. A construction of $C_0$ semigroup solving (D).

Putting  $u_t = v$ , let us rewrite the problem (D) in the following form:

$$\frac{d}{dt} \begin{bmatrix} u \\ v \end{bmatrix} = \begin{bmatrix} 0 & 1 \\ \Delta & -1 \end{bmatrix} \begin{bmatrix} u \\ v \end{bmatrix} = A \begin{bmatrix} u \\ v \end{bmatrix}.$$

To consider A to be dissipative, we introduce a space  $H_D(\Omega)$ . For any open set  $\mathcal{O} \subset \mathbb{R}^n$ ,  $C_0^{\infty}(\Omega)$  denotes the space of all  $C^{\infty}$  functions on  $\mathbb{R}^n$  whose support is compact and lies in  $\mathcal{O}$  ( in particular, such functions vanish near the boundary of  $\mathcal{O}$ ),  $L^2(\mathcal{O})$  a usual  $L^2$  space on  $\mathcal{O}$  with norm  $\|\cdot\|_{\mathcal{O}}$  innerproduct  $(\cdot,\cdot)_{\mathcal{O}}$  and  $H^s(\mathcal{O})$  a usual Sobolev space of order s on  $\mathcal{O}$  with norm  $\|\cdot\|_{s,\mathcal{O}}$ .  $\|\cdot\|_{k,\Omega}$  will be denoted simply by  $\|\cdot\|_k$ . Likewise for  $\|\cdot\|_{\Omega}$  and  $(\cdot,\cdot)_{\Omega}$ . Then, we put

$$H_D(\Omega) = \{ u \in H^1_{loc}(\Omega) \mid \nabla u = (\frac{\partial u}{\partial x_1}, \dots, \frac{\partial u}{\partial x_n}) \in L^2(\Omega), \quad u = 0 \quad \text{on } \Gamma,$$
$$\exists \ \{ u_n \} \subset C_0^{\infty}(\Omega) \text{ s.t. } \|\nabla (u_n - u)\| \to 0 \text{ as } n \to \infty \ \},$$

where  $H^1_{loc}(\Omega) = \{u \in \mathcal{D}'(\Omega) \mid u \in H^1(\Omega_R) \ \forall R > r_0\}$ .  $H_D(\Omega)$  has the following properties.

**Theorem 2.1.** If  $u \in H_D(\Omega)$ , then u satisfies the following inequalities:

$$||u||_{0,\Omega_R} \le C(R) ||\nabla u||_{0,\Omega_R},$$
$$\int_{\Omega} \frac{|u(x)|^2}{d(x)^2} dx \le C ||\nabla u||^2.$$

Moreover,  $H_D(\Omega)$  is a Hilbert space equipped with an inner product  $(u, v)_D = (\nabla u, \nabla v)$ .

Then, an underlying space for A is

$$\mathcal{H} = \left\{ \begin{bmatrix} u \\ v \end{bmatrix} \mid u \in H_D(\Omega) \ v \in L^2(\Omega) \right\}.$$

From Theorem 2.1 we know that  $\mathcal{H}$  is a Hilbert space equipped with the innerproduct

$$\left( \begin{bmatrix} u \\ v \end{bmatrix}, \begin{bmatrix} w \\ z \end{bmatrix} \right)_{\mathcal{H}} = (u, w)_D + (v, z).$$

The domain of A is

$$D(A) = \left\{ \begin{bmatrix} u \\ v \end{bmatrix} \in \mathcal{H} \mid A \begin{bmatrix} u \\ v \end{bmatrix} \in \mathcal{H} \right\}$$
$$= \left\{ \begin{bmatrix} u \\ v \end{bmatrix} \in \mathcal{H} \mid v \in H_D(\Omega), \ \Delta u \in L^2(\Omega) \right\}.$$

Then, A has the following properties.

**Proposition 2.2.** (1) A is a closed operator. (2) A is a dissipative operator.

(3) 
$$\mathcal{R}(I-A) = \mathcal{H}$$
. (4)  $D(A)$  is dense in  $\mathcal{H}$ .

Lumer and Phillips theorem [13, Chapter 1, Theorem 4.3] implies that A generates a  $C^0$  semigroup  $\{T(t)\}$  on  $\mathcal{H}$ .

### §3. A proof of Theorem 1.2.

Our purpose in this section is to prove the following result, which implies our main theorem.

#### Theorem 3.1.

$$\|\varphi_R T(t)\mathbf{x}\|_{\mathcal{H}} \le C(1+t)^{-n/2} \|\mathbf{x}\|_{\mathcal{H}},$$

for  $\mathbf{x} \in \mathcal{H}_{1,R}$ , where C = C(R).

Sketch of proof.

Since A is dissipative, T(t) is a  $C_0$  semigroup of contractions, so that

$$||T(t)|| \le 1 \quad \forall t \ge 0.$$

Let  $\alpha$  be a positive number. In view of (3.1), we have the following expression:

(3.2) 
$$T(t)\mathbf{x} = \lim_{\omega \to \infty} \frac{1}{2\pi i} \int_{\alpha = i\omega}^{\alpha + i\omega} e^{\lambda t} (\lambda I - A)^{-1} \mathbf{x} d\lambda \quad \text{for } \mathbf{x} \in D(A^2).$$

(cf. [12, p.295] or [13, Chapter 1, Corollary 7.5]). By a lemma due to F. Huang in [4, §1, Lemma 1] (also see [7]), we have the following lemma.

**Lemma 3.2.** For any  $\alpha > 0$  and  $\mathbf{x} \in \mathcal{H}$ , put

$$g(\omega) = \|((\alpha + i\omega)I - A)^{-1}\mathbf{x}\|_{\mathcal{H}}.$$

Then  $g(\omega) \in L^2(\mathbb{R})$  and

$$\lim_{|\omega| \to \infty} g(\omega) = 0,$$

$$\int_{-\infty}^{\infty} g(\omega)^2 d\omega \le \frac{\pi}{\alpha} ||\mathbf{x}||_{\mathcal{H}}^2.$$

In view of Lemma 3.2, the high frequency part decays sufficiently fast, so that we have to investigate the low frequency part. Now we shall introduce some functional spaces. Let E be a Banach space with norm  $|\cdot|_E$ ,  $N \geq 0$  an integer and  $k = N + \sigma$  with  $0 < \sigma \leq 1$ . Put

$$C^{k}(\mathbb{R}^{1}; E) = \{ u \in C^{N-1}(\mathbb{R}^{1}; E) \cap C^{\infty}(\mathbb{R}^{1} - \{0\}; E); \ll u \gg_{k, E} < \infty \},$$

where

$$\ll u \gg_{k,E} = \sum_{j=0}^{N} \int_{\mathbb{R}} |\left(\frac{d}{d\tau}\right)^{j} u(\tau)|_{E} d\tau$$

$$+ \sup_{h \neq 0} |h|^{-\sigma} \int_{\mathbb{R}} |\left(\frac{d}{d\tau}\right)^{N} u(\tau+h) - \left(\frac{d}{d\tau}\right)^{N} u(\tau)|_{E} d\tau \quad \text{if } 0 < \sigma < 1,$$

$$\ll u \gg_{k,E} = \sum_{j=0}^{N} \int_{\mathbb{R}} |\left(\frac{d}{d\tau}\right)^{j} u(\tau)|_{E} d\tau$$

$$+ \sup_{h \neq 0} |h|^{-1} \int_{\mathbb{R}} |\left(\frac{d}{d\tau}\right)^N u(\tau + 2h) - 2\left(\frac{d}{d\tau}\right)^N u(\tau + h) + \left(\frac{d}{d\tau}\right)^N u(\tau)|_E d\tau,$$

if  $\sigma = 1$ . Here,  $\left(\frac{d}{d\tau}\right)^0 = 1$ . The following lemma is concerned with the properties of the Fourier transformation of functions belonging to  $\mathcal{C}^k(\mathbb{R}^1, E)$ , which was proved in [14, Part 1, Theorem 3.7].

**Lemma 3.3.** Let E be a Banach space with norm  $|\cdot|_E$ . Let  $N \geq 0$  be an integer and  $\sigma$  a positive number  $\leq 1$ . Assume that  $f \in \mathcal{C}^{N+\sigma}(\mathbb{R}^1; E)$ . Put

$$F(t) = \frac{1}{2\pi} \int_{-\infty}^{\infty} f(\tau) \exp(\sqrt{-1}\tau t) d\tau.$$

Then,

$$|F(t)|_E \le C(1+|t|)^{-(N+\sigma)} \ll f \gg_{N+\sigma,E}$$

Here and hereafter, we put  $\mathcal{H}_R = \{ \begin{bmatrix} u \\ v \end{bmatrix} \in \mathcal{H} \mid \text{ supp } u, \text{ supp } v \subset \Omega_R \}$ .  $\varphi_R(x)$  always refers to a function in  $C_0^{\infty}(\mathbb{R}^n)$  such that  $\varphi_R(x) = 1$  if  $|x| \leq R$  and  $|x| \geq R + 1$ . Moreover, we put

$$\mathcal{H}_{loc} = \left\{ \begin{bmatrix} u \\ v \end{bmatrix} \mid u \in H^1(\Omega_R), v \in L^2(\Omega_R) \ \forall R \geq r_0 \right\},$$
 $\mathcal{H}_{comp} = \bigcup_{R \geq r_0} \mathcal{H}_R,$ 

and  $\mathcal{L}(B_1, B_2)$  denotes the set of all bounded linear operators from  $B_1$  into  $B_2$  and Anal(I, B) the set of all B-valued analytic functions in I. In view of Lemma 3.3, if we prove the following fact, the proof of Theorem 3.1 is complete.

(F) Put  $Q_d = \{\lambda \in \mathbb{C} \mid 0 < \Re \lambda < d, |\Im \lambda| < d\}$ . Then, there exists a d > 0 and  $R(\lambda) \in Anal(Q_d; \mathcal{L}(\mathcal{H}_{comp}, \mathcal{H}_{loc}))$  such that :

(a) 
$$R(\lambda)\mathbf{x} = (\lambda I - A)^{-1}\mathbf{x}$$
 for  $\mathbf{x} \in \mathcal{H}_{comp}$  and  $\lambda \in Q_d$ ;

(b) For any  $R \geq r_0$  and  $\rho(s) \in C_0^{\infty}(\mathbb{R})$  such that  $\rho(s) = 1$  if |s| < d/2 and |s| > d, there exist  $M_1 > 0$  depending on R,  $\rho$  and d such that

$$\ll \rho(\cdot)(\varphi_R R(\alpha+i\cdot)\mathbf{x},\mathbf{y})_{\mathcal{H}} \gg_{n/2,\mathbb{R}} \leq M_1 \|\mathbf{x}\|_{\mathcal{H}} \|\mathbf{y}\|_{\mathcal{H}},$$

for any  $\mathbf{x} \in \mathcal{H}_R$ ,  $\mathbf{y} \in \mathcal{H}$  and  $0 < \alpha < d$ .

We shall conclude this report by giving a brief proof of (F).

Proof of (F).

When  $n \geq 3$ , (F) was proved by Shibata [14, Part 1], so that we shall consider the case that n = 2. Corresponding stationary problem is

(3.3) 
$$(\lambda^2 + \lambda - \Delta)u = f \text{ and } u = 0 \text{ on } \Gamma.$$

If  $|\lambda|$  is small, then in stead of (3.3), it is sufficient to consider the following problem:

$$(A_{\lambda})$$
  $(\lambda - \Delta)u = f \text{ in } \Omega \subset \mathbb{R}^2 \text{ and } u = 0 \text{ on } \Gamma,$ 

where  $\lambda \in S_{r,\varepsilon} = \{\lambda \in \mathbb{C} \setminus \{0\} \mid |\lambda| < r, |\arg \lambda| < \pi - \varepsilon\}, 0 < r < 1 \text{ and } 0 < \varepsilon < \pi/2,$  because  $\lambda^2 + \lambda$  is equivalent to  $\lambda$  for small  $|\lambda|$ . In view of Lemma 3.4 of [14], in order to prove (F), it is sufficient to prove the following propositions.

**Proposition 3.4.** For  $\lambda \in S_{r,\varepsilon}$  and  $r_0 \leq R < \infty$ , there exists  $A(\lambda) : L_R^2(\Omega) \to H^1_{loc}(\Omega)$  satisfying that

$$(\lambda-\Delta)A(\lambda)f=f\quad \text{in }\Omega\quad \text{ and }\ A(\lambda)f=0\quad \text{on }\Gamma,$$

for  $f \in L^2_R(\Omega)$ . Moreover, it satisfies that

$$\|\varphi_R A(\lambda)f\|_1 \le C\|f\|$$
 as  $\lambda \in S_{r,\varepsilon}$ .

**Proposition 3.5.** For  $\lambda$  and R as mentioned above, following estimates hold;

$$\|\varphi_R \frac{d}{d\lambda} A(\lambda) f\|_1 \le \frac{C(R)}{|\lambda| |\log \lambda|^2} \|f\| \le \frac{C(R)}{|\Im \lambda|} \|f\|,$$
  
$$\|\varphi_R \frac{d^2}{d\lambda^2} A(\lambda) f\|_1 \le \frac{C(R)}{|\lambda|^2 |\log \lambda|^2} \|f\| \le \frac{C(R)}{|\Im \lambda|^2} \|f\|,$$

for  $f \in L^2_R(\Omega)$ .

Our main idea to prove Propositions 3.4 and 3.5 is to use the single layer potential and the double layer potential, and to reduce  $(A_{\lambda})$  to a boundary integral equation. Put  $v = (\lambda - \Delta)^{-1} f$ . Then v is represented by the modified Bessel function:

$$(\lambda - \Delta)^{-1} f = \int_{\mathbb{R}^2} E_{\lambda}(x - y) f(y) dy,$$

where  $E_{\lambda}(x) = (2\pi)^{-1} K_0(|x|\sqrt{\lambda})$ ,  $K_m$  ( $m \in \mathbb{N} \cup \{0\}$ ) denotes the modified Bessel function of order m. So we want to solve the equation

$$(A'_{\lambda})$$
  $(\lambda - \Delta)w = 0$  in  $\Omega$  and  $w = f_{\lambda}$  on  $\Gamma$ ,

where  $f_{\lambda} = (\lambda - \Delta)^{-1} f|_{\Gamma}$ . To do this, let us introduce the integral operator  $B_{\lambda}$ :

$$B_{\lambda}\Phi = D_{\lambda}\Phi - \eta E_{\lambda}M\Phi + rac{2\pi\alpha}{\log\sqrt{\lambda}}E_{\lambda}\Phi \quad ext{ for } \Phi \in C^0(\Gamma).$$

Here  $\alpha$ ,  $\eta > 0$ ,  $E_{\lambda}$  is a single layer potential defined by

$$E_{\lambda}\Psi(x) = \int_{\Gamma} E_{\lambda}(x-y)\Psi(y)do_{y}$$

and  $D_{\lambda}$  is a double layer potential defined by

$$D_{\lambda}\Psi(x) = \int_{\Gamma} D_{\lambda}(x,y)\Psi(y)do_{y},$$

where

$$\begin{split} D_{\lambda}(x,y) &= \nabla_{x} E_{\lambda}(x-y) \cdot N(y) \\ &= -\frac{1}{2\pi} K_{1}(|x-y|\sqrt{\lambda}) \frac{\sqrt{\lambda}}{|x-y|} (x-y) \cdot N(y). \end{split}$$

The projection  $M:\,C^0(\Gamma) o C^0(\Gamma)$  is defined by

$$\Phi \to M\Phi = \Phi - \Phi_M \quad \text{ with } \Phi_M = \frac{1}{|\Gamma|} \int_{\Gamma} \Phi do \text{ and } |\Gamma| = \operatorname{meas}(\Gamma).$$

Obviously  $B_{\lambda}\Phi$  satisfies that  $(\lambda - \Delta)B_{\lambda}\Phi = 0$  in  $\Omega$ , so that we obtain the following boundary integral equation:

(3.4) 
$$B_{\lambda}\Phi|_{\Gamma} = K_{\lambda}\Phi = \left(-\frac{1}{2} + D_{\lambda} - \eta E_{\lambda}M + \frac{2\pi\alpha}{\log\sqrt{\lambda}}E_{\lambda}\right)\Phi = f_{\lambda}.$$

If  $\Phi$  is a solution of (3.4),  $B_{\lambda}\Phi$  satisfies  $(A'_{\lambda})$ , and  $A(\lambda)f$  is expressed by

(3.5) 
$$A(\lambda)f = (\lambda - \Delta)^{-1}f - B_{\lambda}\Phi.$$

Therefore,  $(A_{\lambda})$  was reduced to a boundary integral equation (3.4).  $K_{\lambda}$  is a Fredholm operator, so that by using the Fredholm alternative theorem, we can solve the boundary equation (3.4). If we consider that  $A(\lambda)$  is an operator from  $L_R^2(\Omega)$  to  $L_{loc}^2(\Omega)$ , by the properties of Bessel function, we know that the expansion of  $A(\lambda)$  at  $\lambda \to 0$  is

$$A(\lambda) = C_0 + C_1 \frac{1}{\log \lambda} + \cdots.$$

Therefore, we have Propositions 3.4 and 3.5.

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