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# ON THE EXPONENTIAL DECAY OF SOLUTIONS OF THE WAVE EQUATION WITH THE POTENTIAL FUNCTION 

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1. Introduction and results. We consider the Schroedinger operator $L=-\Delta+c(x)$ with the potential function $c(x)$ satisfying the following condition $\left(A_{1}\right)$ in the whole 3-dimensional Euclidean space $E$, where $\Delta$ designates the 3dimensional Laplacian and $x$ a position vector in $E$ with its length $|x|$ :

$\left(\mathrm{A}_{1}\right) \quad\left\{\right.$| $c(x)$ is a real-valued $\mathscr{B}^{1}$ function defined on $E$ and satisfies |
| :---: |
| $\|c(x)\| \leqslant C_{1} \mathrm{e}^{-\delta \mid x_{1}} \quad(x \in E)$ |
| where $C_{1}$ and $\delta$ are positive constants. |

Here $\mathscr{B}^{1}$ stands for the space of all bounded, continuous functions $f(x)$ defined over $E$ with bounded, continuous first derivatives.

It is well known that under condition $\left(A_{1}\right)$ the symmetric operator $L$ on $\mathscr{D}$ is lower semi-bounded and essentially self-adjoint in $L^{2}=L^{2}(E)$, where $\mathscr{D}$ consists of all infinitely many times differentiable functions with compact support in $E$ (see T. Kato [2], Section 6, Theorem 1). Then we denote again by $L$ the unique self-adjoint extension with domain $\mathscr{D}_{L^{2}}^{2} . \quad \mathscr{D}_{L^{2}}^{m}(m=1,2)$ is the completion of the space $\mathscr{D}$ with respect to the norm $\|f\|_{m}=\left(\sum_{|\alpha| \leqslant m} \int_{E}\left|D^{\alpha} f(x)\right|^{2} d x\right)^{1 / 2}$, or is equivalently the space with the norm $\|\cdot\|_{m}$ of all functions $f(x) \in L^{2}$ whose derivatives $D^{\alpha} f(x)(|\alpha| \leqslant m)$ in the distribution sense all belong to $L^{2}(\alpha=$ $\left(\alpha_{1}, \alpha_{2}, \alpha_{3}\right)$ with the $\alpha_{k}$ 's non-negative integers; $D_{k}=\frac{\partial}{\partial x_{k}}, D^{\alpha}=D_{1}^{\alpha_{1}} D_{2}^{\alpha} D_{3}^{\alpha}$, and $\left.|\alpha|=\alpha_{1}+\alpha_{2}+\alpha_{3}\right)$.

The spectrum of $L$ can be only on the real axis. On the whole positive axis there exists only the essential spectrum, which is, in fact, absolutely continuous, while on the negative axis we have only the discrete point spectrum, if any (see T. Ikebe [1], Chapter 2, Section 7). Here we assume
$\left(\mathrm{A}_{2}\right) \quad$ The operator $L$ has no negative eigenvalues.
Let $B_{\delta}$ be the space of all continuous functions $f(x)$ defined on $E$ with
$\|f\|_{\delta}=\sup _{x \in M}\left|\mathrm{e}^{-(\delta / 2)|x|} f(x)\right|<\infty$. Then $B_{\delta}$ becomes a Banach space with the norm $\|f\|_{\delta}$. Furthermore, we impose another condition on $c(x)$ :
$\left(\mathrm{A}_{3}\right) \quad\left\{\begin{array}{l}\text { The homogeneous integral equation } \\ f(x)=-\frac{1}{4 \pi} \int_{E} \frac{c(y)}{|x-y|} f(y) d y \\ \text { has only the trivial solution } f=0 \text { in } B_{\delta} .\end{array}\right.$
An appendix will be added for a remark on $\left(A_{3}\right)$.
Let $j_{a}(t)=\left\{\begin{array}{l}1(t \geqslant a) \\ t / a(0 \leqslant t<a)\end{array}\right.$, where $a$ is any fixed positive number, and $\omega$ be real. Now we consider in the free space $E$ the initial value problem for the wave equation

$$
\left\{\begin{array}{c}
\frac{\partial^{2}}{\partial t^{2}} v_{1}(x, t)+L v_{1}(x, t)=q(x) \mathrm{e}^{i \omega t} j_{a}(t)  \tag{1.1a}\\
v_{1}(x, 0)=0, \quad \frac{\partial}{\partial t} v_{1}(x, 0)=g(x),
\end{array}\right.
$$

and the reduced wave equation

$$
\begin{equation*}
L u(x)=\omega^{2} u(x)+q(x) \tag{1.2}
\end{equation*}
$$

under conditions $\left(A_{1}\right),\left(A_{2}\right),\left(A_{3}\right)$,

$$
\left\{\begin{array}{l}
q(x) \text { is a measurable function defined on } E \text { and there exists a pair of }  \tag{1}\\
\text { positive numbers } Q_{1} \text { and } \gamma \text { such that for any } x \text { in } E \\
\qquad|q(x)| \leqslant Q_{1} \mathrm{e}^{-\gamma_{|x|}}
\end{array}\right.
$$

and
(C) $\quad\left\{\begin{array}{c}g(x) \text { is a } C^{2} \text { function defined over } E \text { and has the estimate } \\ \left|D^{\infty} g(x)\right| \leqslant G \mathrm{e}^{-\mu\left|x_{1}\right|} \quad(x \in E,|\alpha| \leqslant 2), \\ \text { where } G \text { and } \mu \text { are positive constants. }\end{array}\right.$

Then the initial value problem (1.1a) has a unique solution $v_{1}(\cdot, t)$ in $L^{2}$, which will be set forth more in detail in Proposition 2.1.

Our first result can be stated as follows:
Theorem 1. Suppose $\left(A_{1}\right),\left(A_{2}\right),\left(A_{3}\right),\left(B_{1}\right)$ and $(C)$. Then the following assertions hold:
(i) (Limiting amplitude principle) There exists the limit function $u(x)=$ $\lim _{t \rightarrow \infty} v_{1}(x, t) \mathrm{e}^{-i \omega t}$ uniformly on any bounded set in $E$, which is a solution of (1.2).
(ii) (Exponential decay) $v_{1}(x, t)$ can be expressed as

$$
\left\{\begin{array}{l}
v_{1}(x, t)=u(x) \mathrm{e}^{i \omega t}+\delta_{1}(x, t) \\
\left|\delta_{1}(x, t)\right| \leqslant C \mathrm{e}^{\mathrm{e} / 2)|x|} \mathrm{e}^{-a_{1} t},
\end{array}\right.
$$

where $\alpha_{1}$ is some positive constant $<\min (\delta / 2, \gamma, \mu)$, and $C$ is a positive constant depending only on $c(x), q(x), g(x), \omega, a$ and $\alpha_{1}$.
(iii) (Sommerfeld radiation principle) $u(x)$ satisfies the Sommerfeld radiation conditions

$$
\left\{\begin{array}{l}
u(x)=O\left(|x|^{-1}\right) \\
\frac{d u(x)}{d|x|}+i \omega u(x)=o\left(|x|^{-1}\right)
\end{array}\right.
$$

as $|x| \rightarrow \infty$.
In the sequel the letter $C$ exclusively means a positive constant. $C$ does not always denote the same one.

In addition to $\left(A_{1}\right),\left(A_{2}\right),\left(A_{3}\right),\left(B_{1}\right)$ and $(C)$ assume
$\left(\mathrm{B}_{2}\right) \quad\left\{\begin{array}{c}q(x) \text { has continuous first derivatives satisfying } \\ |\operatorname{grad} q(x)| \leqslant Q_{2} \mathrm{e}^{-\nu|x|} \quad(x \in E), \\ \text { where } Q_{2} \text { and } \nu \text { are positive contants, }\end{array}\right.$
and consider the initial value problem

$$
\left\{\begin{array}{l}
\frac{\partial^{2}}{\partial t^{2}} v_{2}(x, t)+L v_{2}(x, t)=q(x) \mathrm{e}^{i \omega t}  \tag{1.1b}\\
\quad v_{2}(x, 0)=0, \quad \frac{\partial}{\partial t} v_{2}(x, 0)=g(x)
\end{array}\right.
$$

where $j_{a}(t)$ which appeared in (1.1a) has been deleted. Then there exists a unique solution $v_{2}(\cdot, t) \in L^{2}$ (see Proposition 2.1), and our second result is as follows:

Theorem 2. Under assumptions $\left(A_{1}\right),\left(A_{2}\right),\left(A_{3}\right),\left(B_{1}\right),\left(B_{2}\right)$ and $(C)$ the statements ( $i$ ), (ii) and (iii) in Theorem 1 are all valid for (1.1b) and (1.2), if $v_{1}(x, t)$ is replaced there by $v_{2}(x, t)=u(x) \mathrm{e}^{i \omega t}+\delta_{2}(x, t)$, and the estimate for $\delta_{1}(x, t)$ by $\left|\delta_{2}(x, t)\right| \leqslant C \mathrm{e}^{(\delta / 2)|x|} \mathrm{e}^{-\alpha_{2} t}$, where $\alpha_{2}$ is some positive constant $<\min (\delta / 2, \gamma, \nu, \mu)$, and where $C$ depends only on $c(x), q(x), g(x), \omega$ and $\alpha_{2}$.
O.A. Ladyženskaja [3] has given a proof of Theorem 1 on the basis of the Laplace transformation theory where both $c(x)$ and $q(x)$ are assumed to have a compact support, and (1.1a) to have zero initial data. Fundamentally on the same line as hers we shall prove the two above-mentioned theorems.

Meanwhile, C.S. Morawetz has studied the decay of solutions of the initialboundary value problem for the wave equation $v_{t t}(x, t)-\Delta v(x, t)=0$ with zero Dirichlet condition in the exterior of a star-shaped reflecting body in $E$. In [8]
and [9] she has obtained the rate of decay with time $t$ at least like $t^{-1 / 2}$ and $t^{-1}$ respectively by the Kirchhoff formula and certain estimates derived from socalled energy identities (Friedrichs' $a, b, c$-method). Moreover, P.D. Lax, C.S. Morawetz and R.S. Phillips [4] have proved the exponential energy decay of the solutions, following the ideas developed in [8] and [5]. Recently C.S. Morawetz [10] has deduced an exponential energy decay from a certain preassumed decay rate for not-necessarily star-shaped domains (she has considered Robin as well as Dirichlet boundary conditions), and, applying this criterion for the exponential decay to the problem in [4], has given another direct proof. [9] also contains the result that the solution $v$ of the initial-boundary value problem for the inhomogeneous equation $v_{t t}(x, t)-\Delta v(x, t)=q(x) \mathrm{e}^{i \omega t}$ with zero Dirichlet condition approaches to a solution $u$ of its reduced equation $\Delta u(x)=\omega^{2} u(x)+$ $q(x)$ as fast as $t^{-1 / 2}$.
2. Laplace transformations. Under the hypotheses $\left(A_{1}\right),\left(B_{1}\right)$ and $(C)$ the general existence theorem on the initial value problem for hyperbolic equations (refer e.g to S. Mizohata [7], Chapter 6 or [6]) guarantees

Proposition 2.1. Suppose $\left(A_{1}\right),\left(B_{1}\right)$ and $(C)$. Then there exists a unique solution $v_{1}(x, t)$ of $(1.1 a)$ such that $\left(v_{1}(x, t), \frac{\partial}{\partial t} v_{1}(x, t)\right) \in \mathcal{E}_{t \geqslant 0}^{0}\left(\mathscr{D}_{L^{2}}^{2}\right) \times \mathcal{E}_{i>0}^{0}\left(\mathscr{D}_{L^{2}}^{1}\right)$ and

$$
\begin{equation*}
\left\|v_{1}(\cdot, t)\right\|_{2}+\left\|\frac{\partial}{\partial t} v_{1}(\cdot, t)\right\|_{1} \leqslant C \mathrm{e}^{\beta t} \tag{2.1}
\end{equation*}
$$

where $\beta$ is some positive constant, and $f(t) \in \mathcal{E}_{t>0}^{0}\left(\mathscr{D}_{L^{2}}^{m_{2}}\right)(m=1,2)$ means that $f(t) \in \mathscr{D}_{L^{2}}^{m}$ and is continuous on the interval $t \geqslant 0$ in the topology of $\mathscr{D}_{L^{2}}^{m}$. The same is true for $v_{2}(x, t)$ of $(1.1 b)$.

By Proposition 2.1. the Laplace transform $w_{k}(x, \lambda)$ of $v_{k}(x, t)$

$$
\begin{equation*}
w_{k}(x, \lambda)=\int_{0}^{\infty} v_{k}(x, t) \mathrm{e}^{-\lambda t} d t \quad(\operatorname{Re} \lambda>\beta ; k=1,2) \tag{2.2}
\end{equation*}
$$

exists in $\mathscr{D}_{L^{2}}^{2}$ and is analytic in $\operatorname{Re} \lambda>\beta$, where $\operatorname{Re} \lambda$ denotes the real part of $\lambda$. Moreover, the inverse transformation of $w_{k}(x, \lambda)$

$$
\begin{equation*}
v_{k}(x, t)=\lim _{A \rightarrow \infty} \frac{1}{2 \pi i} \int_{\sigma-i A}^{\sigma+i A} w_{k}(x, \lambda) \mathrm{e}^{\lambda t} d \lambda \quad(k=1,2) \tag{2.3}
\end{equation*}
$$

can be carried out along any path $\operatorname{Re} \lambda=\sigma>\beta$.
Applying the Laplace transformation to the initial value problems (1.1a) and (1.1b) in $\operatorname{Re} \lambda>\beta$, we have

$$
\begin{equation*}
\left(L+\lambda^{2}\right) w_{k}(x, \lambda)=q(x) f_{k}(\lambda)+g(x) \quad(k=1,2) \tag{2.4}
\end{equation*}
$$

where

$$
\left\{\begin{array}{l}
f_{1}(\lambda)=\frac{1-\mathrm{e}^{a_{i(i \omega-\lambda)}}}{a(i \omega-\lambda)^{2}} \\
f_{2}(\lambda)=\frac{1}{\lambda-i \omega}
\end{array}\right.
$$

From $\left(A_{2}\right)$ and the statement on the spectrum of $L$ in Section 1 it follows that every point $-\lambda^{2}$ with $\operatorname{Re} \lambda>0$ becomes a regular point of $L$. Hence there exists the resolvent $\left(L+\lambda^{2}\right)^{-1}$ carrying $L^{2}$ onto $\mathscr{D}_{L^{2}}^{2}$ in $\operatorname{Re} \lambda>0$. Therefore (2.4) has a unique solution in $\mathscr{D}_{L^{2}}^{2}$ which is analytic in $\operatorname{Re} \lambda>0$. Then we can consider the Laplace transform $w_{k}(x, \lambda)(k=1,2)$ to be extended analytically to the whole half-plane $\operatorname{Re} \lambda>0$.

Putting

$$
\begin{equation*}
w_{k}(x, \lambda)=u_{k}(x, \lambda) f_{k}(\lambda) \quad(\operatorname{Re} \lambda>0 ; k=1,2), \tag{2.5}
\end{equation*}
$$

we can see from (2.4) and the discussion following it that $u_{k}(x, \lambda)$ is a unique solution in $\mathscr{D}_{L^{2}}^{2}$ of the equation

$$
\begin{equation*}
\left(L+\lambda^{2}\right) u_{k}(x, \lambda)=q(x)+g(x) f_{k}^{-1}(\lambda) \quad(\operatorname{Re} \lambda>0 ; k=1,2), \tag{2.6}
\end{equation*}
$$

which is analytic in $\operatorname{Re} \lambda>0$. Since in $\operatorname{Re} \lambda>0$ there exists the resolvent $\left(-\Delta+\lambda^{2}\right)^{-1}$, which is an integral operator of Carleman type with the kernel $(4 \pi|x-y|)^{-1} \mathrm{e}^{-\lambda|x-y|}, u_{k}(x, \lambda)$ in (2.6) satisfies the integral equation

$$
\begin{align*}
u_{k}(x, \lambda) & =\frac{1}{4 \pi} \int_{E} \frac{\mathrm{e}^{-\lambda \mid x-y_{\mid}}}{|x-y|}\left(q(y)+g(y) f_{k}^{-1}(\lambda)\right) d y-  \tag{2.7}\\
& -\frac{1}{4 \pi} \int_{E} \frac{\mathrm{e}^{-\lambda \mid x-y_{\mid}} \mid}{|x-y|} c(y) u_{k}(y, \lambda) d y \quad(\operatorname{Re} \lambda>0 ; k=1,2)
\end{align*}
$$

3. Integral equations. In this section we shall study the unique solvability of equation (2.7) in $B_{\delta}$, which leads to the analytical extension to $\operatorname{Re} \lambda \leqslant 0$ of $u_{k}(x, \lambda)$ in (2.5) and then to the meromorphical one of $w_{k}(x, \lambda)$ in $(2.5)(k=1,2)$.

Now let us introduce a domain $D_{\delta}$ defined by $D_{\delta}=\{\lambda ; \operatorname{Re} \lambda>-\delta / 2\}$.
Proposition 3.1. Under assumption $\left(A_{1}\right)$ the integral operator $T_{\lambda}$, defined by

$$
T_{\lambda} f(x)=-\frac{1}{4 \pi} \int_{E} \frac{\mathrm{e}^{-\lambda|x-y|}}{|x-y|} c(y) f(y) d y
$$

is completely continuous on $B_{\delta}$ into itself for each fixed $\lambda \in D_{\delta}$.
Proof. Let $f \in B_{\delta}$ and $\lambda \in D_{\delta}$. When $\operatorname{Re} \lambda<0$, we first have

$$
\left|T_{\lambda} f(x)\right| \leqslant C| | f\left\|_{\delta} \int_{E} \frac{\mathrm{e}^{-\operatorname{Re} \lambda\left|x-y_{\mid}-(\delta / 2 \mid)\right| y_{\mid}}}{|x-y|} d y=C\right\| f \|_{\delta} J .
$$

To estimate $J$ we note that $|x-y| \leqslant|x|+|y|$, and $|x-y| \geqslant|y|$ when $|y| \leqslant|x| / 2$, and $|y| \geqslant(1 / 3)|x-y|$ when $|y| \geqslant|x| / 2$. We then obtain

$$
\begin{aligned}
J & \leqslant \mathrm{e}^{-\operatorname{Re} \lambda|x|}\left(\int_{|y| \leqslant|x| / 2} \frac{\mathrm{e}^{-(1 / 2)(\delta+2 \operatorname{Re} \lambda) \mid y_{\mid}}}{|y|} d y+\int_{|y| \geqslant|x| / 2} \frac{\mathrm{e}^{-(1 / 6)(\delta+2 \operatorname{Re} \lambda)|x-y|}}{|x-y|} d y\right) \\
& \leqslant C(\delta+2 \operatorname{Re} \lambda)^{-2} \mathrm{e}^{-\operatorname{Re} \lambda|x|} .
\end{aligned}
$$

Hence we have

$$
\begin{equation*}
\left|T_{\lambda} f(x)\right| \leqslant C(\delta+2 \operatorname{Re} \lambda)^{-2}| | f \|_{\delta} \mathrm{e}^{-\operatorname{Re} \lambda|x|} \quad(\operatorname{Re} \lambda<0), \tag{3.1}
\end{equation*}
$$

where $C$ is independent of $x$ and $\lambda$. When $\operatorname{Re} \lambda \geqslant 0$, the estimate for $T_{\lambda} f(x)$ can be worked out with more ease in a similar fashion by using the inequality $\left|\mathrm{e}^{-\lambda \mid x-y}\right| \leqslant 1$, and we get

$$
\begin{equation*}
\left|T_{\lambda} f(x)\right| \leqslant C\|f\|_{\delta} \quad(\operatorname{Re} \lambda \geqslant 0) \tag{3.2}
\end{equation*}
$$

where $C$ does not depend on $x$ nor $\lambda$. From (3.1) and (3.2) it follows that

$$
\mathrm{e}^{-(\delta / 2)|x|}\left|T_{\lambda} f(x)\right| \leqslant\left\{\begin{array}{ll}
C(\delta+2 \operatorname{Re} \lambda)^{-2}| | f \|_{\delta} \mathrm{e}^{-(\delta / 2+\operatorname{Re} \lambda) \mid x_{\mid}}  \tag{3.3}\\
C\|f\|_{\delta} \mathrm{e}^{-(\delta / 2)|x|} \quad & (\operatorname{Re} \lambda \geqslant 0)
\end{array} \quad(\operatorname{Re} \lambda<0)\right.
$$

Next we proceed to show the continuity of $T_{\lambda} f(x)$ in $x$ for each fixed $\lambda \in D_{\delta}$. For this purpose we consider the difference

$$
\begin{aligned}
T_{\lambda} f(x)-T_{\lambda} f\left(x^{\prime}\right)= & -\frac{1}{4 \pi} \int_{E} \frac{\mathrm{e}^{-\lambda|x-y|}-\mathrm{e}^{-\lambda \mid x^{\prime}-y_{\mid}}}{|x-y|} c(y) f(y) d y- \\
& -\frac{1}{4 \pi} \int_{E}\left(\frac{1}{|x-y|}-\frac{1}{\left|x^{\prime}-y\right|}\right) \mathrm{e}^{-\lambda\left|x^{\prime}-y\right|} c(y) f(y) d y \\
= & J_{1}+J_{2}
\end{aligned}
$$

Considering ( $A_{1}$ ) and the inequality

$$
\left|\mathrm{e}^{-\lambda|x-y|}-\mathrm{e}^{-\lambda\left|x^{\prime}-y\right|}\right| \leqslant\left\{\begin{array}{l}
\mathrm{e}^{-\operatorname{Re} \lambda \max \left(|x-y|,\left|x^{\prime}-y\right|\right)}|\lambda|\left|x^{\prime}-x\right| \quad(\operatorname{Re} \lambda<0) \\
|\lambda|\left|x^{\prime}-x\right| \quad(\operatorname{Re} \lambda \geqslant 0)
\end{array}\right.
$$

we have, when $\operatorname{Re} \lambda<0$,

$$
\begin{aligned}
&\left|J_{1}\right| \leqslant C\|f\|_{\delta}\left|x^{\prime}-x\right|\left(\int_{\left|x^{\prime}-y\right| \leq|x-y|} \frac{\mathrm{e}^{-\operatorname{Re} \lambda\left|x-y_{\mid}\right|-(\delta / 2 \mid)|y|}}{|x-y|} d y+\right. \\
&\left.+\int_{\left|x^{\prime}-y\right| \geqslant|x-y|} \frac{\mathrm{e}^{-\operatorname{Re} \lambda\left|x^{\prime}-y_{\mid}-(\delta / 2)\right| y \mid}}{|x-y|} d y\right) \\
& \leqslant C\|f\|_{\delta}\left|x^{\prime}-x\right|\left(1+\mathrm{e}^{-\operatorname{Re} \lambda\left|x^{\prime}-x_{1}\right|}\right) \int_{E} \frac{\mathrm{e}^{-\operatorname{Re} \lambda\left|x-y_{\mid}-(\delta / 2)\right| y \mid}}{|x-y|} d y \\
& \leqslant C\|f\|_{\delta}\left|x^{\prime}-x\right|\left(1+\mathrm{e}^{-\operatorname{Re} \lambda\left|x^{\prime}-x\right|}\right) \mathrm{e}^{-\operatorname{Re} \lambda|x|},
\end{aligned}
$$

$$
\begin{aligned}
\left|J_{2}\right| & \leqslant C| | f \|_{\delta}\left|x^{\prime}-x\right|\left(\int_{E} \frac{\mathrm{e}^{-2 \operatorname{Re} \lambda\left|x^{\prime}-y_{\mid}-(1 / 2)(\delta-2 \operatorname{Re} \lambda)\right| y_{\mid}}}{\left|x^{\prime}-y\right|^{2}} d y\right)^{1 / 2} \times \\
& \times\left(\int_{E} \frac{\mathrm{e}^{-(1 / 2)(\delta+2 \operatorname{Re} \lambda)\left|y_{1}\right|}}{|x-y|^{2}} d y\right)^{1 / 2} \leqslant C\|f\|_{\delta}\left|x^{\prime}-x\right| \mathrm{e}^{-\operatorname{Re} \lambda\left|x^{\prime}\right|}
\end{aligned}
$$

in a way similar to the one used for obtaining (3.1). Similarly, when $\operatorname{Re} \lambda \geqslant 0$, we get

$$
\begin{aligned}
\left|J_{1}\right| & \leqslant C\|f\|_{\delta}\left|x^{\prime}-x\right| \int_{E} \frac{\mathrm{e}^{-(\delta / 2) \mid y_{\mid}}}{|x-y|} d y \leqslant C\|f\|_{\delta}\left|x^{\prime}-x\right|, \\
\left|J_{2}\right| & \leqslant C\|f\|_{\delta}\left|x^{\prime}-x\right|\left(\int_{E} \frac{\mathrm{e}^{-(\delta / 2) \mid y_{\mid}}}{|x-y|^{2}} d y\right)^{1 / 2}\left(\int_{E} \frac{\mathrm{e}^{-(\delta / 2) \mid y_{\mid}}}{\left|x^{\prime}-y\right|^{2}} d y\right)^{1 / 2} \\
& \leqslant C\|f\|_{\delta}\left|x^{\prime}-x\right| .
\end{aligned}
$$

So it follows that in $E \times D_{\delta}$

$$
\left|T_{\lambda} f(x)-T_{\lambda} f\left(x^{\prime}\right)\right| \leqslant\left\{\begin{array}{ll}
C\|f\|_{\delta}\left|x^{\prime}-x\right| \mathrm{e}^{-\operatorname{Re} \lambda\left(\left|x_{\mid}+\left|x^{\prime}-x\right|\right)\right.}  \tag{3.4}\\
C| | f \|_{\delta}\left|x^{\prime}-x\right| & (\operatorname{Re} \lambda \geqslant 0)
\end{array} \quad(\operatorname{Re} \lambda<0)\right.
$$

where $C$ 's depend only on $\lambda$. By (3.3) and (3.4) $T_{\lambda}$ is a bounded linear operator on $B_{\delta}$ into itself.

Let $S$ be any bounded set in $B_{\delta}$. We want to show that $T_{\lambda} S$ is relatively compact in $B_{\delta}$, which deduces the complete cotinuity of $T_{\lambda}$. Let $\left\{g_{n}\right\}(n=1,2$, $\cdots$ ) be any sequence chosen from $T_{\lambda} S$. By (3.1), (3.2), (3.4) and the boundedness of $S,\left\{g_{n}(x)\right\}$ must be a uniformly bounded and equicontinuous family of continuous functions on any compact domain of $E$. Employing the AscoliArzelà selection theorem we can find out a subsequence $\left\{g_{n^{\prime}}(x)\right\}$ converging to a continuous function $g(x)$ uniformly on any compact domain of $E$. In view of (3.3), $\mathrm{e}^{-(\delta / 2) \mid x} g_{n^{\prime}}(x)$ tends to 0 uniformly in $n^{\prime}$ as $|x| \rightarrow \infty$. So does $\mathrm{e}^{-(\delta / 2)|x|} g(x)$. Then $\left\|g_{n^{\prime}}-g\right\|_{\delta} \rightarrow 0$ as $n^{\prime} \rightarrow \infty$, and $g \in B_{\delta}$, which is the desired result.

Proposition 3.2. Assume $\left(A_{1}\right),\left(A_{2}\right)$ and $\left(A_{3}\right)$. Then there exists a positive number $\alpha(<\delta / 2)$ such that for any $g \in B_{\delta}$ the equation $\left(I-T_{\lambda}\right) f=g$ admits one and only one solution $f \in B_{\delta}$ for each fixed $\lambda \in D=\{\lambda ; \operatorname{Re} \lambda>-\alpha\}$, which is analytic in $D$, and $\left(I-T_{\lambda}\right)^{-1}$ has the estimate

$$
\left\|\left(I-T_{\lambda}\right)^{-1}\right\|_{\delta} \leqslant C
$$

uniformly in $D$, where $\|T\|_{\delta}=\sup _{\|f\|_{\delta}=1}\|T f\|_{\delta}$ for an operator $T$ in $B_{\delta}$ and $I$ stands for the identity operator in $B_{\delta}$.

Proposition 3.2 will be proved by the following four lemmas.
Lemma 3.1. Suppose $\left(A_{1}\right),\left(A_{2}\right)$ and $\left(A_{3}\right)$. Let $\operatorname{Re} \lambda \geqslant 0$ and $g \in B_{\delta}$. Then the integral equation

$$
\begin{equation*}
f=g+T_{\lambda} f \tag{3.5}
\end{equation*}
$$

has a unique solution $f \in B_{\delta}$, which is analytic on $\operatorname{Re} \lambda \geqslant 0$.
Proof. Let $f \in B_{\delta}$ be a solution of the homogeneous equation $f=T_{\lambda} f$. When $\operatorname{Re} \lambda>0$, we get

$$
|f(x)|=\left|T_{\lambda} f(x)\right| \leqslant C\|f\|_{\delta} \int_{E} \frac{\mathrm{e}^{-\operatorname{Re} \lambda\left|x-y_{\mid}-(\delta / 2)\right| y_{\mid}}}{|x-y|} d y=C\|f\|_{\delta} J,
$$

and, noting that $|x-y| \geqslant|y|$ and $|x-y| \geqslant \frac{|x|}{2}$ when $|y| \leqslant \frac{|x|}{2}$,

$$
\begin{aligned}
J & \leqslant \mathrm{e}^{-(\operatorname{Re} \lambda / 4)|x|} \int_{|y| \leqslant|x| / 2} \frac{\mathrm{e}^{-(\operatorname{Re} \lambda / 2+\delta / 2)\left|y_{\mid}\right|}}{|y|} d y+\mathrm{e}^{-(\delta / 4)|x|} \int_{|y| \geqslant \mid x / / 2} \frac{\mathrm{e}^{-\operatorname{Re} \lambda|x-y|}}{|x-y|} d y \\
& \leqslant C\left(\mathrm{e}^{-(\operatorname{Re} \lambda / 4)|x|}\left(\frac{\operatorname{Re} \lambda}{2}+\frac{\delta}{2}\right)^{-2}+\mathrm{e}^{-(\delta / 4)|x|}(\operatorname{Re} \lambda)^{-2}\right),
\end{aligned}
$$

where $C$ 's are independent of $x$ and $\lambda$. So we can see that

$$
|f(x)|=\left|T_{\lambda} f(x)\right| \leqslant C\|f\|_{\delta}\left(\mathrm{e}^{-(\operatorname{Re} \lambda / 4)|x|}+\mathrm{e}^{-(\delta / 4)|x|}\right),
$$

where $C$ is dependent only on $\lambda$. Therefore $f \in L^{2}$ for $\lambda$ with $\operatorname{Re} \lambda>0$. Also when $\operatorname{Re} \lambda=0$ and $\lambda \neq 0, f \in L^{2}$. This follows from A. Ya. Povzner [11], Chapter 2, Lemmas 1, 2, 5 and 6. As $f$ fulfills the equation $\left(L+\lambda^{2}\right) f=0,-\lambda^{2}$ ( $\operatorname{Re} \lambda \geqslant 0, \lambda \neq 0$ ) cannot be an eigenvalue of $L$ by $\left(A_{2}\right)$ and the statement on the spectrum of $L$ in Section 1. Hence, if we also note $\left(A_{3}\right)$ for $\lambda=0$, the equation $f=T_{\lambda} f$ implies $f=0$ in $B_{\delta}$ for $\lambda$ with $\operatorname{Re} \lambda \geqslant 0$. So by the Riesz-Schauder theory together with Proposition 3.1, equation (3.5) is uniqely solvable in $B_{\delta}$ for any $g \in B_{\delta}$ and the operator $I-T_{\lambda}$ has a bounded inverse in $B_{\delta}$. Moreover, from its definition in Proposition $3.1 T_{\lambda}$ is seen to be analytic for $\operatorname{Re} \lambda \geqslant 0$ (cf. e.g. K. Yosida [13], Chapter 5, Section 3). Thus $\left(I-T_{\lambda}\right)^{-1}$ is also analytic for $\operatorname{Re} \lambda \geqslant 0$, which was to be proved (cf. e.g. ibid., Chapter 8, Section 2).

Lemma 3.2. Assume $\left(A_{1}\right),\left(A_{2}\right)$ and $\left(A_{3}\right)$. Then for any $N^{\prime}>0$ one can find a positive $\alpha^{\prime}(<\delta / 2)$ such that for $\lambda \in D^{\prime}=\left\{\lambda ;-\alpha^{\prime}<\operatorname{Re} \lambda,|\lambda|<N^{\prime}\right\}$ equation (3.5) has a unique solution in $B_{\delta}$, which is analytic in $D^{\prime}$, and the estimate

$$
\left\|\left(I-T_{\lambda}\right)^{-1}\right\|_{\delta} \leqslant C
$$

holds in $D^{\prime}$, where $C$ is dependent only on $D^{\prime}$.
Proof. Clear from Lemma 3.1 and the Heine-Borel theorem.
Lemma 3.3. Assume $\left(A_{1}\right)$ and let $0<b<1$. Then, for any $\varepsilon>0$ there exists a positive number $N$ independent of $(x, y) \in E \times E$ such that the kernel of the operator
$T_{\lambda}{ }^{2}$

$$
\tau(x, y ; \lambda)=(4 \pi)^{-2} c(y) \int_{E} \frac{\mathrm{e}^{-\lambda|x-s|-\lambda|s-y|}}{|x-s||s-y|} c(s) d s
$$

has the estimate

$$
\mathrm{e}^{-(\delta / 2)|x|}|\tau(x, y ; \lambda)|<\varepsilon \mathrm{e}^{-(1 / 2+(1-b) / 4) \delta|y|}
$$

for $(x, y, \lambda) \in E \times E \times\left\{\lambda ; \operatorname{Re} \lambda>-\frac{b}{2} \delta,|\lambda|>N\right\}$.
Proof. Let $\operatorname{Re} \lambda>-\frac{b}{2} \delta$. We first assume $(x, y) \in\left(E-K_{R}\right) \times E$ or $K_{R} \times\left(E-K_{R}\right)$, where $K_{R}$ means the sphere of radius $R$ with its center at the origin. Since by an estimation similar to the one utilized for having (3.1) we have

$$
\begin{aligned}
|\tau(x, y ; \lambda)| & \leqslant C \mathrm{e}^{-\delta|y|}\left(\int_{E} \frac{\mathrm{e}^{b \delta|x-s|-\delta|s|}}{|x-s|^{2}} d s\right)^{1 / 2}\left(\int_{E} \frac{\mathrm{e}^{b \delta\left|s-y_{1}-\delta\right| s \mid}}{|s-y|^{2}} d s\right)^{1 / 2} \\
& \leqslant C \mathrm{e}^{-(1-b / 2) \delta|y|} \mathrm{e}^{(b / 2) \delta|x|},
\end{aligned}
$$

we can see that

$$
\begin{align*}
& \mathrm{e}^{-(\delta / 2)|x|}|\tau(x, y ; \lambda)| \leqslant C \mathrm{e}^{-(1 / 2)(1-b) \delta R} \mathrm{e}^{-(1-b / 2) \delta \mid y_{1}} \quad\left((x, y) \in\left(E-K_{R}\right) \times E\right)  \tag{3.6}\\
& \mathrm{e}^{-(\delta / 2)|x|}|\tau(x, y ; \lambda)| \leqslant C \mathrm{e}^{-(1 / 4)(1-b) \delta R} \mathrm{e}^{-(1 / 2+(1-b) / 4) \delta|y|} \quad\left((x, y) \in K_{R} \times\left(E-K_{R}\right)\right) . \tag{3.7}
\end{align*}
$$

Next, assuming $(x, y) \in K_{R} \times K_{R}$, we similarly have

$$
\begin{align*}
& \mathrm{e}^{-(\delta / 2)|x|}|\tau(x, y ; \lambda)| \leqslant C \mathrm{e}^{-(\delta / 2)|x|-\delta|y|}\left|\int_{K_{R}} \frac{\mathrm{e}^{-\lambda|x-s|-\lambda|s-y|}}{|x-s||s-y|} c(s) d s\right|+  \tag{3.8}\\
&+C \mathrm{e}^{-(\delta / 2)|x|-\delta|y|}\left|\int_{B-K_{R}}\right|=J_{1}+J_{2}
\end{align*}
$$

$$
\begin{align*}
J_{2} & \leqslant C \mathrm{e}^{-(1 / 2)(1-b) \delta R} \mathrm{e}^{-(1-b / 2) \delta|y|}\left(\int_{E} \frac{\mathrm{e}^{-(1 / 2)(1-b) \delta|s|}}{|x-s|^{2}} d s\right)^{1 / 2}\left(\int_{E} \frac{\mathrm{e}^{-(1 / 2)(1-b) \delta|s|}}{|s-y|^{2}} d s\right)^{1 / 2}  \tag{3.9}\\
& \leqslant C \mathrm{e}^{-(1 / 2)(1-b) \delta R} \mathrm{e}^{-(1-b / 2) \delta|y|} .
\end{align*}
$$

To estimate $J_{1}$ we consider the ellipsoid $|x-s|+|s-y| \leqslant \xi$ with the foci $x$ and $y$ for each fixed $(x, y) \in K_{R} \times K_{R}$ and a positive $\xi$. Let us denote such an ellipsoid by $E_{\xi}(x, y)$. Then, noting that $E_{4 R}(x, y)$ contains $K_{R}$ for any $(x, y) \in K_{R} \times K_{R}$, we have

$$
\begin{equation*}
J_{1}=\left.C \mathrm{e}^{-(\delta / 2)|x|-\delta|y|}\right|_{E_{|x-y|+\rho}} \int \frac{\mathrm{e}^{-\lambda \mid x, y)}}{|x-s||s|-\lambda|s-y|} c(s) d s+ \tag{3.10}
\end{equation*}
$$

$$
\begin{gathered}
+\int_{E_{4 R}(x, y)-E_{|x-y|+\rho}(x, y)}-\int_{E_{4 R}(x, y)-K_{R}} \mid \\
\leqslant \mathrm{e}^{-\delta|y|}\left(\left|C \int_{E_{|x-y|+\rho}(x, y)}\right|+\left|C \int_{E_{4 R}(x, y)-E_{|x-y|+\rho}(x, y)}\right|\right)+C \mathrm{e}^{-(\delta / 2)|x|-\delta|y|} \times \\
\times \int_{E_{4 R}(x, y)-K_{R}} \mid=\mathrm{e}^{-\delta|y|}\left(J_{11}+J_{12}\right)+J_{13} \quad(0<\rho<2 R),
\end{gathered}
$$

$$
\begin{equation*}
J_{13} \leqslant C \mathrm{e}^{-(\delta / 2)|x|-\delta|y|}\left|\int_{E-K_{R}}\right|=J_{2} \leqslant C \mathrm{e}^{-(1 / 2)(1-b) \delta R^{-(1-b / 2) \delta|y|}} \tag{3.11}
\end{equation*}
$$

Here $C$ 's appearing in (3.6) to (3.11) all depend only on $b$. Now for any $\varepsilon>0$ we can choose a sufficiently large $R$ such that $C \mathrm{e}^{-(1 / 2)(1-b) \delta R}$ in (3.6), $C \mathrm{e}^{-(1) 4)(1-b) \delta R}$ in (3.7), and $C \mathrm{e}^{-(1 / 2)(1-b) \delta R}$ in (3.9) and (3.11) are all smaller than $\varepsilon / 4$. Furthermore

$$
J_{11} \leqslant C \int_{\substack{\left.E_{|x-x|} \mid(x, x) y\right) \\|x| s|s| s-y}} \frac{d s}{|x-s|^{2}},
$$

where $C$ depends only on $b$ and $R$. In $E_{|x-y|+\rho}(x, y)$ we introduce a cylindrical coordinate system $(r, \theta, t)$ such that $(0, \theta, 0)$ corresponds to $x$ and the $t$-axis is directed from $x$ to $y$. Since $|x-s|^{2}=r^{2}+t^{2}$ for $s=(r, \theta, t)$ and the Jacobian for the coordinate transformation becomes $r$, a simple evaluation gives

$$
\begin{equation*}
J_{11} \leqslant C \int_{0}^{2 \pi} d \theta \int_{-R}^{R} d t \int_{0}^{C v_{\bar{\rho}}} \frac{r}{\bar{t}^{2}+r^{2}} d r<\frac{\varepsilon}{4} \tag{3.12}
\end{equation*}
$$

for a sufficiently small $\rho(<2 R)$ uniformly in $(x, y) \in K_{R} \times K_{R} . \quad C$ 's in (3.12) never depend on $\rho$. For the estimation of $J_{12}$ we first assign the 3-dimensional orthogonal coordinates $\left(s_{1}, s_{2}, s_{3}\right)$ to each point $s$ in $E_{4 R}(x, y)-E_{|x-y|+\rho}(x, y)$ in such a way that the origin 0 is the middle point of $x$ and $y$, and the $s_{1}$-axis the line directed from 0 to $y$. Secondly, considering $s$ to be a radius vector which starts from 0 , we introduce two angles $\theta$ and $\varphi ; \theta$ is measured from the $s_{3}$-axis toward $s$, and $\varphi$ from the $s_{1}$-axis toward the projection of $s$ to the $\left(s_{1}, s_{2}\right)$-plane. Finally we set $\xi=|x-s|+|s-y|$. Thus for each fixed $(x, y) \in K_{R} \times K_{R}$ every point $s=\left(s_{1}, s_{2}, s_{3}\right)$ in $E_{4 R}(x, y)-E_{|x-y|+\rho}(x, y)$ can be expressed in terms of new coordinates $(\xi, \theta, \varphi)(|x-y|+\rho \leqslant \xi \leqslant 4 R, 0 \leqslant \theta \leqslant \pi, 0 \leqslant \varphi \leqslant 2 \pi)$. Now we have

$$
\begin{aligned}
& |x-s| \geqslant \frac{\rho}{2}, \quad|s-y| \geqslant \frac{\rho}{2} ; \quad \frac{\partial}{\partial \xi}|x-s| \leqslant C, \quad \frac{\partial}{\partial \xi}|s-y| \leqslant C ; \\
& |J(\xi, \theta, \varphi)| \leqslant C, \quad \frac{\partial}{\partial \xi}|J(\xi, \theta, \varphi)| \leqslant C
\end{aligned}
$$

where $J(\xi, \theta, \varphi)$ is the Jacobian for the coordinate transformation, and $C$ 's depend only on $R$ and $\rho$. By integration by parts in consideration of $\left(A_{1}\right)$ and these inequalities, $J_{12}$ can be estimated as follows:

$$
\begin{aligned}
J_{12}= & \left|C \int_{\substack{|x-y|+\rho \leqslant \xi \leqslant 4 \\
0 \leqslant \theta \leqslant \pi, 0<\varphi \leqslant 2 \pi}} \frac{\mathrm{e}^{-\lambda \xi} c(\xi, \theta, \varphi)|J(\xi, \theta, \varphi)|}{|x-s(\xi, \theta, \varphi)||s(\xi, \theta, \varphi)-y|} d \xi d \theta d \varphi\right| \\
= & \left\lvert\,-\frac{C}{\lambda} \int d \theta d \varphi\left[\frac{\mathrm{e}^{-\lambda \xi} c(\xi, \theta, \varphi)|J(\xi, \theta, \varphi)|}{|x-s||s-y|}\right]_{\xi=\mid x-x^{1}+\rho}^{\xi=4 R}+\right. \\
& \left.+\frac{C}{\lambda} \int d \theta d \varphi \int \mathrm{e}^{-\lambda \xi} \frac{\partial}{\partial \xi}\left(\frac{c(\xi, \theta, \varphi)|J(\xi, \theta, \varphi)|}{|x-s||s-y|}\right) d \xi \right\rvert\, \leqslant \frac{C}{|\lambda|} .
\end{aligned}
$$

Therefore for any given $\varepsilon>0$ we can find an $N>0$ independent of $(x, y) \in K_{R}$ $\times K_{R}$ such that

$$
\begin{equation*}
J_{12}<\frac{\varepsilon}{4} \tag{3.13}
\end{equation*}
$$

holds for $(x, y, \lambda) \in K_{R} \times K_{R} \times\left\{\lambda ; \operatorname{Re} \lambda>-\frac{b}{2} \delta,|\lambda|>N\right\}$. Considering (3.6), (3.7) and (3.8), together with (3.9) to (3.13) and the statement just below (3.11), we have Lemma 3.3.

Lemma 3.4. Assume $\left(A_{1}\right)$. Let $0<b<1$. Then there exists an $N^{\prime \prime}>0$ such that the equation $f=g+T_{\lambda} f$ has one and only one solution in $B_{\delta}$ for any $g \in B_{\delta}$ and $\lambda \in D^{\prime \prime}=\left\{\lambda ; \operatorname{Re} \lambda>-\frac{b}{2} \delta,|\lambda|>N^{\prime \prime}\right\}$, which is analytic in $D^{\prime \prime}$, and $\left(I-T_{\lambda}\right)^{-1}$ is uniformly bounded there in the operator norm.

Proof. By Lemma 3.3, for any $\eta(0<\eta<1)$ we can find out an $N^{\prime \prime}>0$ such that the inequality

$$
\begin{equation*}
\left\|T_{\lambda}^{2}\right\|_{\delta} \leqslant \eta \tag{3.14}
\end{equation*}
$$

holds in $D^{\prime \prime}=\left\{\lambda ; \operatorname{Re} \lambda>-\frac{b}{2} \delta,|\lambda|>N^{\prime \prime}\right\} . \quad$ By (3.14) the series

$$
\left(I+T_{\lambda}\right)\left(I+T_{\lambda}^{2}+T_{\lambda}^{4}+\cdots+T_{\lambda}^{2 n}+\cdots\right)
$$

converges to a bounded linear operator in $B_{\delta}$ uniformly in $D^{\prime \prime}$ in the operator norm. Moreover, multiplication of the series by $I-T_{\lambda}$ on the left or right gives $I$, so that the series actually represents $\left(I-T_{\lambda}\right)^{-1}$. Its uniform boundedness and analyticity in $D^{\prime \prime}$ follow from the series, (3.14) and the analyticity of $T_{\lambda}$, which proves the assertion of Lemma 3.4.

Combining Lemmas 3.4 and 3.2 under the assumptions of Proposition 3.2, we obtain Proposition 3.2.

Let $a^{(1)}(x, \lambda)$ and $a^{(2)}(x, \lambda)$ be functions given by

$$
\begin{aligned}
& a^{(1)}(x, \lambda)=\frac{1}{4 \pi} \int_{E} \frac{\mathrm{e}^{-\lambda|x-y|}}{|x-y|} q(y) d y, \\
& a^{(2)}(x, \lambda)=\frac{1}{4 \pi} \int_{E} \frac{\mathrm{e}^{-\lambda|x-y|}}{|x-y|} g(y) d y
\end{aligned}
$$

Then the integral equation (2.7) can be rewritten in the form

$$
\begin{equation*}
\left(I-T_{\lambda}\right) u_{k}(\cdot, \lambda)=a^{(1)}(\cdot, \lambda)+a^{(2)}(\cdot, \lambda) f_{k}^{-1}(\lambda) \quad(k=1,2) . \tag{3.15}
\end{equation*}
$$

Proposition 3.3a. Suppose $\left(B_{1}\right)$. Let $D_{\gamma}=\{\lambda ; \operatorname{Re} \lambda>-\gamma\}$. Then the inequality

$$
\left|a^{(1)}(x, \lambda)\right| \leqslant \begin{cases}C(\gamma+\operatorname{Re} \lambda)^{-2} \mathrm{e}^{-\operatorname{Re} \lambda|x|} & (\operatorname{Re} \lambda<0) \\ C & (\operatorname{Re} \lambda \geqslant 0)\end{cases}
$$

holds in $E \times D_{\gamma}$, where C's are independent of $x$ and $\lambda$.
Proof. Let $\lambda \in D_{\gamma}$. By virtue of the estimation used for obtaining (3.1) and (3.2), we clearly have, when $\operatorname{Re} \lambda<0$,

$$
\begin{aligned}
\left|a^{(1)}(x, \lambda)\right| & \leqslant C \int_{E} \frac{\mathrm{e}^{-\operatorname{Re} \lambda|x-y|-\gamma|y|}}{|x-y|} d y \\
& \leqslant C \mathrm{e}^{-\operatorname{Re} \lambda|x|}(\gamma+\operatorname{Re} \lambda)^{-2}
\end{aligned}
$$

and, when $\operatorname{Re} \lambda \geqslant 0$,

$$
\left|a^{(1)}(x, \lambda)\right| \leqslant C \int_{E} \frac{\mathrm{e}^{-\gamma|y|}}{|x-y|} d y \leqslant C
$$

where $C$ 's are all independent of $x$ and $\lambda$. This completes the proof.
Proposition 3.3b. Assume $\left(B_{1}\right)$ and $\left(B_{2}\right)$. Let $D_{\gamma, \nu}=\{\lambda ; \operatorname{Re} \lambda>-\min (\gamma, \nu)\}$. Then $a^{(1)}(x, \lambda)$ has the estimate
$\left|a^{(1)}(x, \lambda)\right| \leqslant\left\{\begin{aligned} & C \mathrm{e}^{-\operatorname{Re} \lambda|x|}\left(1+|x|+(\gamma+\operatorname{Re} \lambda)^{-1}+(\nu+\operatorname{Re} \lambda)^{-1}|x|+\right. \\ &\left.+(\nu+\operatorname{Re} \lambda)^{-2}\right)|\lambda|^{-1} \\ & C(1+|x|)|\lambda|^{-1} \quad(\operatorname{Re} \lambda \geqslant 0)(\operatorname{Re} \lambda<0)\end{aligned}\right.$
in $E \times D_{\gamma, \nu}$, where $C$ 's are not dependent on any of $x$ and $\lambda$.
Proof. Let $\lambda \in D_{\gamma, \nu}$. Introducing spherical coordinates, we can express $a^{(1)}(x, \lambda)$ as

$$
\begin{align*}
a^{(1)}(x, \lambda) & =\frac{1}{4 \pi} \int_{0}^{2 \pi} d \varphi \int_{0}^{\pi} \sin \theta d \theta \int_{0}^{\infty} \mathrm{e}^{-\lambda \rho} \rho q(x+t) d \rho  \tag{3.16}\\
& =\frac{1}{4 \pi} \int_{0}^{2 \pi} d \varphi \int_{0}^{\pi} \sin \theta d \theta \cdot J \quad(y-x=t,|t|=\rho) .
\end{align*}
$$

Integrating by parts in view of $\left(B_{1}\right)$ and $\left(B_{2}\right)$, we obtain

$$
\begin{align*}
J & =\frac{1}{\lambda} \int_{0}^{\infty} \mathrm{e}^{-\lambda \rho} \frac{\partial}{\partial \rho}(\rho q(x+t)) d \rho \\
& =\frac{1}{\lambda}\left(\int_{0}^{\infty} \mathrm{e}^{-\lambda \rho} q(x+t) d \rho+\int_{0}^{\infty} \mathrm{e}^{-\lambda \rho} \rho \frac{\partial}{\partial \rho} q(x+t) d \rho\right), \\
|J| & \leqslant C\left(\int_{0}^{\infty} \mathrm{e}^{-\operatorname{Re} \lambda \rho} \mathrm{e}^{-\gamma|\rho-|x||} d \rho+\int_{0}^{\infty} \mathrm{e}^{-\operatorname{Re} \lambda \rho} \rho \mathrm{e}^{-\nu|\rho-|x||} d \rho\right)|\lambda|^{-1}  \tag{3.17}\\
& =C\left(J_{1}+J_{2}\right)|\lambda|^{-1},
\end{align*}
$$

where

$$
\begin{aligned}
& J_{1}=\mathrm{e}^{-\gamma|x|} \int_{0}^{|x|} \mathrm{e}^{(\gamma-\operatorname{Re} \lambda) \rho} d \rho+\mathrm{e}^{\gamma|x|} \int_{|x|}^{\infty} \mathrm{e}^{-(\gamma+\operatorname{Re} \lambda) \rho} d \rho \\
& J_{2}=\mathrm{e}^{-\nu|x|} \int_{0}^{|x|} \mathrm{e}^{(\nu-\operatorname{Re} \lambda) \rho} \rho d \rho+\mathrm{e}^{\nu|x|} \int_{|x|}^{\infty} \mathrm{e}^{-(\nu+\operatorname{Re} \lambda) \rho} \rho d \rho
\end{aligned}
$$

When $\operatorname{Re} \lambda<0$, we get

$$
\begin{aligned}
& J_{1} \leqslant C \mathrm{e}^{-\operatorname{Re} \lambda|x|}\left(1+(\gamma+\operatorname{Re} \lambda)^{-1}\right), \\
& J_{2} \leqslant C \mathrm{e}^{-\operatorname{Re} \lambda|x|}\left(|x|+(\nu+\operatorname{Re} \lambda)^{-1}|x|+(\nu+\operatorname{Re} \lambda)^{-2}\right) .
\end{aligned}
$$

Hence, when $\operatorname{Re} \lambda<0$, (3.17) becomes

$$
\begin{equation*}
|J| \leqslant C \mathrm{e}^{-\operatorname{Re} \lambda|x|}\left(1+|x|+(\gamma+\operatorname{Re} \lambda)^{-1}+(\nu+\operatorname{Re} \lambda)^{-1}|x|+(\nu+\operatorname{Re} \lambda)^{-2}\right)|\lambda|^{-1}, \tag{3.18}
\end{equation*}
$$ where $C$ is independent of $x$ and $\lambda$. When $\operatorname{Re} \lambda \geqslant 0$, we have

$$
\begin{aligned}
& J_{1} \leqslant \begin{cases}\mathrm{e}^{-\gamma|x|} \int_{0}^{|x|} d \rho+\mathrm{e}^{\gamma|x|} \int_{|x|}^{\infty} \mathrm{e}^{-\gamma \rho} d \rho \leqslant C & (\operatorname{Re} \lambda \geqslant \gamma) \\
\mathrm{e}^{-\gamma|x|} \int_{0}^{|x|} \mathrm{e}^{\gamma \rho} d \rho+\mathrm{e}^{\gamma|x|} \int_{|x|}^{\infty} \mathrm{e}^{-\gamma \rho} d \rho \leqslant C & (\operatorname{Re} \lambda<\gamma),\end{cases} \\
& J_{2} \leqslant \begin{cases}\mathrm{e}^{-\nu|x| x \mid} \int_{0}^{|x|} \rho d \rho+\mathrm{e}^{\nu|x|} \int_{|x|}^{\infty} \mathrm{e}^{-\nu \rho} \rho d \rho \leqslant C(1+|x|) & (\operatorname{Re} \lambda \geqslant \nu) \\
\mathrm{e}^{-\nu|x|} \int_{0}^{|x|} \mathrm{e}^{\nu \rho} \rho d \rho+\mathrm{e}^{\nu|x|} \int_{|x|}^{\infty} \mathrm{e}^{-\nu \rho} \rho d \rho \leqslant C(1+|x|) & (\operatorname{Re} \lambda<\nu) .\end{cases}
\end{aligned}
$$

Consequently, when $\operatorname{Re} \lambda \geqslant 0$, we can see that

$$
J_{1} \leqslant C, \quad J_{2} \leqslant C(1+|x|),
$$

and

$$
\begin{equation*}
|J| \leqslant C(1+|x|)|\lambda|^{-1}, \tag{3.19}
\end{equation*}
$$

where $C$ 's are independent of $x$ and $\lambda$.
Thus (3.16), (3.18) and (3.19) give the required estimate for $a^{(1)}(x, \lambda)$, which proves Proposition 3.3b.

Proposition 3.4. Suppose (C). Let $D_{\mu}=\{\lambda ; \operatorname{Re} \lambda>-\mu\}$. Then $a^{(2)}(x, \lambda)$ satisfies

$$
\begin{align*}
& \text { (i) }\left|a^{(2)}(x, \lambda)\right| \leqslant \begin{cases}C(\mu+\operatorname{Re} \lambda)^{-2} \mathrm{e}^{-\operatorname{Re} \lambda|x|} & (\operatorname{Re} \lambda<0) \\
C & (\operatorname{Re} \lambda \geqslant 0),\end{cases}  \tag{i}\\
& \text { (ii) }\left|a^{(2)}(x, \lambda)\right| \leqslant \begin{cases}C \mathrm{e}^{-\operatorname{Re} \lambda|x|}\left(1+|x|+(\mu+\operatorname{Re} \lambda)^{-1}(1+|x|)+(\mu+\operatorname{Re} \lambda)^{-2}\right)|\lambda|^{-2} \\
& (\operatorname{Re} \lambda<0) \\
C(1+|x|)|\lambda|^{-2} & (\operatorname{Re} \lambda \geqslant 0)\end{cases}
\end{align*}
$$

in $E \times D_{\mu}$, where C's never depend on $x$ nor $\lambda$.
Proof. Considering ( $C$ ) and the definition of $a^{(2)}(x, \lambda)$, (i) is obtained by an estimation similar to the one in the proof of Proposition 3.3a.

Proceeding as in the proof of Proposition 3.3b, we can establish (ii) in $E \times$ $D_{\mu}$. First we rewrite $a^{(2)}(x, \lambda)$ as

$$
\begin{align*}
a^{(2)}(x, \lambda) & =\frac{1}{4 \pi} \int_{0}^{2 \pi} d \varphi \int_{0}^{\pi} \sin \theta d \theta \int_{0}^{\infty} \mathrm{e}^{-\lambda \rho} \rho g(x+t) d \rho  \tag{3.20}\\
& =\frac{1}{4 \pi} \int_{0}^{2 \pi} d \varphi \int_{0}^{\pi} \sin \theta d \theta \cdot J \quad(y-x=t,|t|==\rho)
\end{align*}
$$

Then, integrating by parts twice in view of $(C)$, we have

$$
\begin{aligned}
J & =\frac{1}{\lambda} \int_{0}^{\infty} \mathrm{e}^{-\lambda \rho} \frac{\partial}{\partial \rho}(\rho g(x+t)) d \rho \\
& =\frac{1}{\lambda^{2}} g(x)+\frac{1}{\lambda^{2}}\left(2 \int_{0}^{\infty} \mathrm{e}^{-\lambda \rho} \frac{\partial}{\partial \rho} g(x+t) d \rho+\int_{0}^{\infty} \mathrm{e}^{-\lambda \rho} \rho \frac{\partial^{2}}{\partial \rho^{2}} g(x+t) d \rho\right),
\end{aligned}
$$

$$
\begin{align*}
|J| & \leqslant C\left(1+\int_{0}^{\infty} \mathrm{e}^{-\operatorname{Re} \lambda \rho} \mathrm{e}^{-\mu|\rho-|x||} d \rho+\int_{0}^{\infty} \mathrm{e}^{-\operatorname{Re} \lambda \rho} \rho \mathrm{e}^{-\mu|\rho-|x||} d \rho\right)|\lambda|^{-2}  \tag{3.21}\\
& \leqslant\left\{\begin{array}{l}
C \mathrm{e}^{-\operatorname{Re} \lambda|x|}\left(1+|x|+(\mu+\operatorname{Re} \lambda)^{-1}(1+|x|)+(\mu+\operatorname{Re} \lambda)^{-2}\right)|\lambda|^{-2}(\operatorname{Re} \lambda<0) \\
C(1+|x|)|\lambda|^{-2} \quad(\operatorname{Re} \lambda \geqslant 0),
\end{array}\right.
\end{align*}
$$

where $C$ 's do not depend on any of $x$ and $\lambda$. (3.20) and (3.21) immediately prove (ii). This completes the proof of Proposition 3.4.

Proposition 3.5. Assume $\left(A_{1}\right),\left(A_{2}\right),\left(A_{3}\right),\left(B_{1}\right)$ and $(C)$. Put $D_{1}=$ $\{\lambda ; \operatorname{Re} \lambda>-\min (\alpha, \gamma, \mu)\}$ with $\alpha$ in Proposition 3.2. Then it follows that
(i) $a^{(1)}(\cdot, \lambda) \in B_{\delta}, a^{(2)}(\cdot, \lambda) \in B_{\delta}$ for every $\lambda \in D_{1}$,
(ii) Equation (3.15) admits a unique solution in $B_{\delta}$ for each $\lambda \in D_{1}$, which is analytic in $D_{1}$, with the estimate $\left\|\left(I-T_{\lambda}\right)^{-1}\right\|_{\delta} \leqslant C$ uniformly in $D_{1}$,
(iii) In terms of this unique solution, $u_{k}(x, \lambda)(k=1,2)$ in (2.5) is extended analytically to $D_{1}$.

Proof. In view of Propositions $3.3 a$ and 3.4 (i) we have only to prove the
continuity of $a^{(1)}(x, \lambda)$ and $a^{(2)}(x, \lambda)$ in $x \in E$ for every $\lambda \in D_{1}$ in order to establish (i). This is proved by using $\left(B_{1}\right)$ and $(C)$ if we proceed as in the estimation for having (3.4).

Now (ii) follows directly from Proposition 3.2 with (i) and the analyticity of $a^{(1)}(\cdot, \lambda)$ and $a^{(2)}(\cdot, \lambda)$ in $D_{1}$. Let $\tilde{u}_{k}(x, \lambda)$ be the unique solution of (3.15) ( $k=1,2$ ).

Furthermore, $\tilde{u}_{k}(\cdot, \lambda) \in L^{2}$ for $\operatorname{Re} \lambda>0$, noting in (3.15) that $a^{(1)}(\cdot, \lambda)$, $a^{(2)}(\cdot, \lambda)$ and $T_{\lambda} \tilde{u}_{k}(\cdot, \lambda)$ all belong to $L^{2}$ for $\operatorname{Re} \lambda>0$ by $\left(B_{1}\right),(C)$ and $\left(A_{1}\right)$ $(k=1,2)$ (cf. e.g. the estimation of $T_{\lambda} f(x)$ in the proof of Lemma 3.1). On the other hand, we have already stated in Section 2 that $u_{k}(\cdot, \lambda)(k=1,2)$ in (2.5), which is analytic in $\operatorname{Re} \lambda>0$, belongs to $L^{2}$ and satisfies (3.15) for $\operatorname{Re} \lambda>0$. By $\left(A_{2}\right)$ and the note on the spectrum of $L$ in Section $1, f=T_{\lambda} f\left(f \in L^{2}, \operatorname{Re} \lambda>0\right)$ implies that $f=0$. Then $u_{k}(x, \lambda)=\tilde{u}_{k}(x, \lambda)(k=1,2)$ in $E \times\{\lambda ; \operatorname{Re} \lambda>0\}$ in consideration of their continuity in $x$ for $\operatorname{Re} \lambda>0$. Thus we have (iii), and the proof is complete.

Proposition 3.6a. Suppose $\left(A_{1}\right),\left(A_{2}\right),\left(A_{3}\right),\left(B_{1}\right)$ and $(C)$. Let $D_{1}=$ $\{\lambda ; \operatorname{Re} \lambda>-\min (\alpha, \gamma, \mu)\}$ with $\alpha$ in Proposition 3.2. Then $w_{1}(x, \lambda)$ in (2.5) can be defined for $(x, \lambda) \in E \times D_{1}$ in the form $w_{1}(x, \lambda)=u_{1}(x, \lambda) f_{1}(\lambda)$ as a meromorphic function in $D_{1}$ with the sole simple pole $\lambda=i \omega$, and has the estimate

$$
\left|w_{1}(x, \lambda)\right| \leqslant\left\{\begin{array}{cc}
C \mathrm{e}^{(\delta / 2)|x|}\left((\gamma+\operatorname{Re} \lambda)^{-2}|\lambda-i \omega|^{-2}+\left(1+(\mu+\operatorname{Re} \lambda)^{-1}+\right.\right. \\
& \left.\left.+(\mu+\operatorname{Re} \lambda)^{-2}\right)|\lambda|^{-2}\right)
\end{array} \quad(\operatorname{Re} \lambda<0)\right.
$$

where $C$ 's are not dependent on $x$ nor $\lambda$.
Proof. Putting $w_{1}(x, \lambda)=u_{1}(x, \lambda) f_{1}(\lambda)$ for $\operatorname{Re} \lambda \leqslant 0, w_{1}(x, \lambda)$ in (2.5) can be extended meromorphically to $D_{1}$ by Proposition 3.5 (iii). By Proposition 3.5 (ii) $u_{1}(x, \lambda)$ may be expressed as

$$
\begin{equation*}
u_{1}(\cdot, \lambda)=\left(I-T_{\lambda}\right)^{-1}\left[a^{(1)}(\cdot, \lambda)+a^{(2)}(\cdot, \lambda) f_{1}^{-1}(\lambda)\right] \quad\left(\lambda \in D_{1}\right), \tag{3.22}
\end{equation*}
$$

where $\left\|\left(I-T_{\lambda}\right)^{-1}\right\|_{\delta} \leqslant C$ uniformly in $D_{1}$. Using Propositions 3.5 (i), $3.3 a$ and 3.4 (ii) we have, in $D_{1}$,

$$
\begin{align*}
& \left\|a^{(1)}(\cdot, \lambda)\right\|_{\delta} \leqslant \begin{cases}C(\gamma+\operatorname{Re} \lambda)^{-2} & (\operatorname{Re} \lambda<0) \\
C & (\operatorname{Re} \lambda \geqslant 0),\end{cases}  \tag{3.23}\\
& \left\|a^{(2)}(\cdot, \lambda)\right\|_{\delta} \leqslant \begin{cases}C\left(1+(\mu+\operatorname{Re} \lambda)^{-1}+(\mu+\operatorname{Re} \lambda)^{-2}\right)|\lambda|^{-2} & (\operatorname{Re} \lambda<0) \\
C|\lambda|^{-2} & (\operatorname{Re} \lambda \geqslant 0)\end{cases} \tag{3.24}
\end{align*}
$$

where $C$ 's are all independent of $\lambda$. Hence, noting (3.22) together with (3.23) and (3.24), we get, in $E \times D_{1}$,

$$
\begin{aligned}
\left|u_{1}(x, \lambda)\right| & \leqslant C \mathrm{e}^{(\delta / 2)|x|}| |\left|a^{(1)}(\cdot, \lambda)\left\|_{\delta}+\left|\left|a^{(2)}(\cdot, \lambda)\right| \|_{\delta}\right| f_{1}^{-1}(\lambda) \mid\right)\right. \\
& \leqslant\left\{\begin{aligned}
& C \mathrm{e}^{(\delta / 2)|x|}\left((\gamma+\operatorname{Re} \lambda)^{-2}+\left(1+(\mu+\operatorname{Re} \lambda)^{-1}+(\mu+\operatorname{Re} \lambda)^{-2}\right) \times\right. \\
&\left.\times\left|f_{1}^{-1}(\lambda)\right||\lambda|^{-2}\right) \\
&(\operatorname{Re} \lambda<0) \\
& C \mathrm{e}^{(\delta / 2)|x|}\left(1+\left|f_{1}^{-1}(\lambda)\right||\lambda|^{-2}\right)(\operatorname{Re} \lambda \geqslant 0),
\end{aligned}\right.
\end{aligned}
$$

where $C$ 's do not depend on any of $x$ and $\lambda$. Now, in view of $w_{1}(x, \lambda)=$ $u_{1}(x, \lambda) f_{1}(\lambda)$, we have proved the desired estimate for $w_{1}(x, \lambda)$ in $E \times D_{1}$. This completes the proof of the proposition.

Proposition 3.6b. Assume $\left(B_{2}\right)$ in addition to the conditions in Proposition 3.6a. Let $D_{2}=\{\lambda ; \operatorname{Re} \lambda>-\min (\alpha, \gamma, \nu, \mu)\}$ with $\alpha$ in Proposition 3.2. Then $w_{2}(x, \lambda)$ in (2.5) may be extended meromorphically to $D_{2}$ in the form $w_{2}(x, \lambda)=$ $u_{2}(x, \lambda) f_{2}(\lambda)$ with the only simple pole $\lambda=i \omega$. Besides, in $E \times D_{2} w_{2}(x, \lambda)$ has the estimate

$$
\left|w_{2}(x, \lambda)\right| \leqslant \begin{cases}C \mathrm{e}^{(\delta / 2)|x|}\left(C_{1}|\lambda-i \omega|^{-1}|\lambda|^{-1}+C_{2}|\lambda|^{-2}\right) & (\operatorname{Re} \lambda<0) \\ C \mathrm{e}^{(\delta / 2)|x|}\left(|\lambda-i \omega|^{-1}|\lambda|^{-1}+|\lambda|^{-2}\right) & (\operatorname{Re} \lambda \geqslant 0)\end{cases}
$$

where

$$
\left\{\begin{array}{l}
C_{1}=1+(\gamma+\operatorname{Re} \lambda)^{-1}+(\nu+\operatorname{Re} \lambda)^{-1}+(\nu+\operatorname{Re} \lambda)^{-2} \\
C_{2}=1+(\mu+\operatorname{Re} \lambda)^{-1}+(\mu+\operatorname{Re} \lambda)^{-2}
\end{array}\right.
$$

and where $C$ 's are independent of $x$ and $\lambda$.
Proof. Considering that our assumptions in this proposition are more stringent than those in Proposition 3.6a, we have (3.24), (3.22) for $u_{2}(\cdot, \lambda)$ and $f_{2}^{-1}(\lambda)$ in $D_{2}$, and the meromorphical extension of $w_{2}(x, \lambda)$ in (2.5) to $D_{2}$, in a way similar to the one in the proof of the preceding proposition. Now by Propositions 3.5 (i) and $3.3 b$ we have the estimate

$$
\left\|a^{(1)}(\cdot, \lambda)\right\|_{\delta} \leqslant\left\{\begin{array}{lr}
C\left(1+(\gamma+\operatorname{Re} \lambda)^{-1}+(\nu+\operatorname{Re} \lambda)^{-1}+(\nu+\operatorname{Re} \lambda)^{-2}\right)|\lambda|^{-1}  \tag{3.25}\\
& (\operatorname{Re} \lambda<0) \\
C|\lambda|^{-1} & (\operatorname{Re} \lambda \geqslant 0)
\end{array}\right.
$$

in $D_{2}$, where $C$ 's do not depend on $\lambda$. From (3.24), (3.25) and (3.22) for $u_{2}(\cdot, \lambda)$ and $f_{2}^{-1}(\lambda)$ it follows that in $E \times D_{2}$

$$
\begin{aligned}
\left|u_{2}(x, \lambda)\right| & \leqslant C \mathrm{e}^{\delta / 2)|x|}\left(\left\|a^{(1)}(\cdot, \lambda)\right\|_{\delta}+\left\|a^{(2)}(\cdot, \lambda)\left|\|_{\delta}\right| f_{2}^{-1}(\lambda) \mid\right)\right. \\
& \leqslant \begin{cases}C \mathrm{e}^{(\delta / 2)|x|}\left(C_{1}|\lambda|^{-1}+C_{2}\left|f_{2}^{-1}(\lambda)\right||\lambda|^{-2}\right) & (\operatorname{Re} \lambda<0) \\
C\left(|\lambda|^{-1}+\left|f_{2}^{-1}(\lambda)\right||\lambda|^{-2}\right) & (\operatorname{Re} \lambda \geqslant 0)\end{cases}
\end{aligned}
$$

where $C_{1}$ and $C_{2}$ are the same as in the proposition, and where $C$ 's are independent of $x$ and $\lambda$. Noting that $w_{2}(x, \lambda)=u_{2}(x, \lambda) f_{2}(\lambda)$, we complete the proof.
4. Proof of the theorems. Now we assume $\left(A_{1}\right),\left(A_{2}\right),\left(A_{3}\right),\left(B_{1}\right)$ and
$(C)$ for the proof of Theorem 1, and, moreover, $\left(B_{2}\right)$ for that of Theorem 2. On the complex plane we choose a rectangular path $\Gamma_{k}$ with vertices $\sigma-i A$, $\sigma+i A,-\alpha_{k}+i A$ and $-\alpha_{k}-i A$, where $\sigma>\beta$ (see Proposition 2.1), $0<\alpha_{1}<$ $\min (\alpha, \gamma, \mu)$ and $0<\alpha_{2}<\min (\alpha, \gamma, \nu, \mu)$ with $\alpha$ in Proposition 3.2, and where $A$ is a positive number large enough for $\Gamma_{k}$ to contain the simple pole $\lambda=i \omega$ of $w_{k}(x, \lambda)(k=1,2)$. Propositions $3.6 a \quad(k=1)$ and $3.6 b(k=2)$ enable us to apply the residue theorem to the contour integral of $w_{k}(x, \lambda) \mathrm{e}^{\lambda t}$ along the positively oriented closed path $\Gamma_{k}$. That is

$$
\int_{\Gamma_{k}} w_{k}(x, \lambda) \mathrm{e}^{\lambda t} d \lambda=2 \pi i \underset{\lambda=i \omega}{\operatorname{Res}}\left[w_{k}(x, \lambda) \mathrm{e}^{\lambda t}\right] \quad(k=1,2) .
$$

Meanwhile, the left-hand side can be divided as

$$
\int_{\Gamma_{k}} w_{k}(x, \lambda) \mathrm{e}^{\lambda t} d \lambda=\int_{\sigma-i A}^{\sigma+i A} w_{k}(x, \lambda) \mathrm{e}^{\lambda t} d \lambda-\int_{-\alpha_{k-i} A}^{-\alpha_{k}+i A}+\int_{\sigma+i A}^{-\alpha_{k}+i A}+\int_{-a_{k}-i A}^{\sigma-i A}
$$

Propositions $3.6 a(k=1)$ and $3.6 b(k=2)$ also assert that the third and the fourth integrlas on the right-hand side tend to zero as $A \rightarrow \infty$ for every $(x, t) \in E \times[0, \infty)$, and that

$$
\operatorname{Res}_{\lambda=i \omega}\left[w_{k}(x, \lambda) \mathrm{e}^{\lambda t}\right]=u_{k}(x, i \omega) \mathrm{e}^{i \omega t} \quad(k=1,2)
$$

In view of equation (3.15) for $\lambda=i \omega$ we can set

$$
u_{k}(x, i \omega)=u(x) \quad(k=1,2)
$$

which is a solution of the reduced wave equation (1.2) by Proposition 3.5 (ii). Thus, recalling (2.3), we have

$$
\begin{equation*}
v_{k}(x, t)=\lim _{A \rightarrow \infty} \frac{1}{2 \pi i} \int_{-\omega_{k-i A}}^{-\infty_{k} k^{+i A}} w_{k}(x, \lambda) e^{\lambda t} d \lambda+u(x) \mathrm{e}^{i \omega t} \quad(k=1,2) \tag{4.1}
\end{equation*}
$$

Furthermore, by Propositions $3.6 a(k=1)$ and $3.6 b(k=2)$ we get

$$
\left|\int_{-\alpha_{k-i A}}^{-\alpha_{k}+i A} w_{k}(x, \lambda) \mathrm{e}^{\lambda t} d \lambda\right| \leqslant\left\{\begin{array}{l}
C \mathrm{e}^{(\delta / 2|x| x \mid} \mathrm{e}^{-\alpha_{1} t} \int_{-A}^{A}\left(\frac{1}{(p-\omega)^{2}+\alpha_{1}^{2}}+\frac{1}{p^{2}+\alpha_{1}^{2}}\right) d p \quad(k=1) \\
C \mathrm{e}^{(\delta / 2)|x|} \mathrm{e}^{-\alpha_{2} t} \int_{-A}^{A}\left(\frac{1}{\sqrt{(p-\omega)^{2}+\alpha_{2}^{2}} \sqrt{p^{2}+\alpha_{2}^{2}}}+\frac{1}{p^{2}+\alpha_{2}^{2}}\right) d p \\
(k=2),
\end{array}\right.
$$

where $p$ is the imaginary part of $\lambda$, and where the first $C$ depends on $\alpha_{1}$ and the second on $\alpha_{2}$, but both are independent of $(x, t) \in E \times[0, \infty)$. As the integrals on the right-hand side converge as $A \rightarrow \infty$, (4.1) can be rewritten as

$$
\left|v_{k}(x, t)-u(x) \mathrm{e}^{i \omega t}\right| \leqslant C \mathrm{e}^{\delta / 2)|x|} \mathrm{e}^{-\omega_{k} t}
$$

where $C$ depends on $\alpha_{k}$ and $\omega$, but not on $x$ nor $t(k=1,2)$. Therefore we have (i) and (ii) in both theorems.

By its definition $u(x)$ satisfies the integral equation

$$
u(x)=\frac{1}{4 \pi} \int_{E} \frac{\mathrm{e}^{-i \omega|x-y|}}{|x-y|} q(y) d y-\frac{1}{4 \pi} \int_{E} \frac{\mathrm{e}^{-i \omega|x-y|}}{|x-y|} c(y) u(y) d y
$$

where $|q(y)| \leqslant Q_{1} \mathrm{e}^{-\gamma|y|}$ and $|c(y) u(y)| \leqslant C \mathrm{e}^{-\left(\delta / 2\left|y^{\mid}\right|\right.}$by means of $\left(B_{1}\right),\left(A_{1}\right)$ and Proposition 3.5 (ii). Then, for the proof of (iii) in both theorems it is sufficient to refer to A. Ya. Povzner [11], Chapter 2, Lemmas 1 and 2.

Thus we have proved Theorems 1 and 2.
Appendix. As was mentioned in Section 1 we remark here that the assumption of $\left(A_{3}\right)$ has a justifiable ground. For we can give an example of $c(x)$ satisfying $\left(A_{1}\right)$ and $\left(A_{2}\right)$, but not $\left(A_{3}\right)$, which means that $\left(A_{3}\right)$ is independent of the others on $c(x)$ ( $c f$. also T. Shirota and K. Asano [12]).

Consider

$$
c_{0}(x)=\left\{\begin{array}{cc}
-1 & \left(|x|<\frac{\pi}{2}\right) \\
0 & \left(|x| \geqslant \frac{\pi}{2}\right)
\end{array}, \quad f_{0}(x)= \begin{cases}\frac{\sin |x|}{|x|} & \left(|x|<\frac{\pi}{2}\right) \\
\frac{1}{|x|} & \left(|x| \geqslant \frac{\pi}{2}\right)\end{cases}\right.
$$

It follows from a simple computation that $f_{0}(x)$ satisfies the equation

$$
f_{0}(x)=-\frac{1}{4 \pi} \int_{E} \frac{c_{0}(y)}{|x-y|} f_{0}(y) d y
$$

However $c_{0}(x)$ is not smooth. So let $\rho(x)=\rho(|x|)$ be a $C^{\infty}$ function such that $\rho(x) \geqslant 0, \rho(x)=0(|x| \geqslant 1)$ and $\int_{E} \rho(x) d x=1$. Moreover, denoting by $*$ the convolution, we put

$$
\begin{equation*}
c_{1}(x)=\frac{\rho * c_{0} f_{0}(x)}{\rho * f_{0}(x)}, \quad f_{1}(x)=\rho * f_{0}(x) . \tag{1}
\end{equation*}
$$

Then $f_{1}(x)=f_{1}(|x|)$ is a strictly postitive solution in $B_{\delta} \cap C^{\infty}(E)$ of the equation

$$
f_{1}(x)=-\frac{1}{4 \pi} \int_{E} \frac{c_{1}(y)}{|x-y|} f_{1}(y) d y
$$

but it is not in $L^{2}$ because it equals $|x|^{-1}$ for $|x| \geqslant \frac{\pi}{2}+1$. Furthermore, the numerator of $c_{1}(x)$ is a $C^{\infty}$ function with support in the sphere $|x| \leqslant \frac{\pi}{2}+1$, and
so is $c_{1}(x)=c_{1}(|x|)$. Hence $c_{1}(x)$ satisfies $\left(A_{1}\right)$, but not $\left(A_{3}\right)$.
Now we are proving that the operator $L_{1}=-\Delta+c_{1}(x)$ has no negative eigenvalue with an $L^{2}$ eigenfunction. Let $\lambda<0$. Let $f \in L^{2}$ be a solution of the equation

$$
\begin{equation*}
L_{1} f=\lambda f \tag{2}
\end{equation*}
$$

If we put, in spherical coordinates,

$$
\begin{equation*}
g(r, \theta, \varphi)=r f(r, \theta, \varphi) \quad(|x|=r) \tag{3}
\end{equation*}
$$

$g$ may be expressed in terms of the series

$$
g(r, \theta, \varphi)=\sum_{n=0}^{\infty} \sum_{m=0}^{2 n} b_{n, m}(r) Y_{n, m}(\theta, \varphi),
$$

where $\left\{Y_{n, m}\right\}$ is a complete orthonormal system of spherical harmonics, and where

$$
\begin{equation*}
b_{n, m}(r)=\int_{0}^{2 \pi} d \varphi \int_{0}^{\pi} g(r, \theta, \varphi) Y_{n, m}(\theta, \varphi) \sin \theta d \theta \tag{4}
\end{equation*}
$$

which fulfills the equation

$$
\begin{equation*}
\frac{d^{2}}{d r^{2}} b_{n, m}(r)+\left(\lambda-c_{1}(r)-\frac{n(n+1)}{r^{2}}\right) b_{n, m}(r)=0 \tag{5}
\end{equation*}
$$

Since $\frac{n(n+1)}{r^{2}}(n \geqslant 1)$ represents a positive definite operator, we have only to prove, in view of (5), (3) and (4), that $b(r) \in L^{2}(0, \infty)$ must identically vanish if $b(r)$ satisfies

$$
\begin{gather*}
b^{\prime \prime}(r)+\left(\lambda-c_{1}(r)\right) b(r)=0,  \tag{6}\\
b(0)=0 \tag{7}
\end{gather*}
$$

By the statement following (1), $b_{0}(r)=r f_{1}(r)$ fulfills

$$
\begin{align*}
& b_{0}(r)>0(r>0), \quad b_{0}(r)=1\left(r \geqslant \frac{\pi}{2}+1\right),  \tag{8}\\
& b_{0}^{\prime \prime}(r)-c_{1}(r) b_{0}(r)=0  \tag{9}\\
& b_{0}(0)=0, \quad b_{0}^{\prime}(0)>0 \tag{10}
\end{align*}
$$

Any solution of equtaion (6) with (7) may be determined except for a constant multiple. Hence we can assume

$$
\begin{equation*}
b^{\prime}(0)>b_{0}{ }^{\prime}(0) . \tag{11}
\end{equation*}
$$

By (7), (10) and (11) we have, for sufficiently small positive values of $r$,

$$
\begin{equation*}
b(r)>b_{0}(r) . \tag{12}
\end{equation*}
$$

Assume that $b(r)$ and $b_{0}(r)$ have the first common value at $r=r_{0}$ except the origin. Then we obtain

$$
\begin{equation*}
b^{\prime}\left(r_{0}\right) \leqslant b_{0}{ }^{\prime}\left(r_{0}\right) . \tag{13}
\end{equation*}
$$

Multiplying (6) and (9) by $b_{0}(r)$ and $b(r)$ respectively, and subtracting, we have

$$
\lambda b_{0}(r) b(r)=b(r) b_{0}{ }^{\prime \prime}(r)-b_{0}(r) b^{\prime \prime}(r) .
$$

Integrating the above equation over $\left[0, r_{0}\right]$, we get, by (7) and (10),

$$
\lambda \int_{0}^{r_{0}} b_{0}(r) b(r) d r=b_{0}\left(r_{0}\right)\left(b_{0}{ }^{\prime}\left(r_{0}\right)-b^{\prime}\left(r_{0}\right)\right) \quad(\lambda<0),
$$

where the left-hand side is negative by (8) and (12), while the right-hand side is non-negative by (8) and (13). This is a contradiction. So by (12) we have

$$
b(r)>b_{0}(r) \quad(r>0) .
$$

Hence, by (8) $b(r)$ does not lie in $L^{2}(0, \infty)$. That is, $b(r) \in L^{2}(0, \infty)$ satisfying (6) and (7) must be identically zero, which was to be proved. Therefore $\left(A_{2}\right)$ is the case with $c_{1}(x)$.

Thus $c_{1}(x)$ is a required example satisfying $\left(A_{1}\right)$ and $\left(A_{2}\right)$, but not $\left(A_{3}\right)$. In other words, $\left(A_{3}\right)$ is not too unnatural a restriction on the potential function $c(x)$.

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