# Spectral representation for Schrödinger operators with magnetic vector potentials 

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## 1. Introduction

The present paper is concerned with the Schrödinger operator

$$
\begin{equation*}
L=-\sum_{j=1}^{n}\left(\frac{\partial}{\partial x_{j}}+i b_{j}\right)^{2}+V \quad \text { in } \quad \boldsymbol{R}^{n}, \tag{1.1}
\end{equation*}
$$

where $b_{j}$ and $V$ denote the multiplication operators by real-valued functions $b_{j}(x)$ and $V(x) . \quad b(x)=\left(b_{1}(x), \ldots, b_{n}(x)\right)$ is the magnetic vector potential (thus rot $b$ represents the magnetic field when $n=3$ ) and $V(x)$ is the electric scalar potential. In their classical work [5], Ikebe and Kato have proved the essential self-adjointness of $L$ on $C_{0}^{\infty}\left(\boldsymbol{R}^{n}\right)$ in the Hilbert space $\mathscr{H}=L_{2}\left(\boldsymbol{R}^{n}\right)$ for an arbitrary $b(x)$ which is continuously differentiable and for $V(x)$ in an appropriate class of functions. Recently, the condition on $b(x)$ has been improved considerably by Leinfelder-Simader [9]. In the present paper, for simplicity, we impose the differentiability condition on $b(x)$ and somewhat strong conditions on the local behavior of $V(x)$ as well as on its decay rate at infinity, which guarantee the uniqueness of the self-adjoint realization of $L$ in $\mathscr{H}$, and we shall denote it by $H$. Moreover, the magnetic field is assumed to tend to zero at infinity (for the study of the Schrödinger operators with constant magnetic fields, see Avron-Herbst-Simon [2]). We assume the following conditions throughout the paper:
( $V$ ) $\quad V(x)$ is a real-valued measurable function and there exist positive constants $C_{0}, \delta$ such that $|V(x)| \leqq C_{0}(1+|x|)^{-1-\delta}$ for all $x \in \boldsymbol{R}^{n}$.
(b) $\quad b_{j}(x)(j=1, \ldots, n)$ are real-valued $C^{2}$ functions and there exist positive constants $C_{0}, \delta$ such that $B_{j k}=\frac{\partial b_{k}}{\partial x_{j}}-\frac{\partial b_{j}}{\partial x_{k}}(j, k=1, \ldots, n)$ satisfies $\left|B_{j k}(x)\right| \leqq C_{0}(1+|x|)^{-3 / 2-\delta}$ for all $x \in \boldsymbol{R}^{n}$.
Logically, the constants $C_{0}$ and $\delta$ in $(V)$ and (b) are different. But we may and do assume that they are identical. Moreover, we can take $\delta$ so that $0<\delta<1$,
$\delta \neq \frac{1}{2}$. It is known that these assumptions imply the following properties of $H$ (Ikebe-Saitō [6], Kuroda [8]):
(i) The essential spectrum of $H$ is $[0, \infty)$,
(ii) $E((0, \infty)) H$ is an absolutely continuous operator, where $E$ denotes the spectral measure associated with $H$.
Our purpose in the present paper is to show that $H$ admits a spectral representation, which needs a stronger assumption on $b(x)$, i.e., the following:
( $b^{\prime}$ ) In addition to (b), $\left|\frac{\partial B_{j k}}{\partial x_{j}}(x)\right| \leqq C_{0}(1+|x|)^{-2-\delta}$ for $j, k=1, \ldots, n$ and for all $x \in \boldsymbol{R}^{n}$.
Namely, we shall establish the existence of an unitary operator $\mathscr{F}$ from the subspace $\mathscr{H}_{a c}$ of absolute continuity for $H$ onto $\hat{\mathscr{H}}=L_{2}\left((0, \infty) ; L_{2}(\Omega)\right)(\Omega$ denotes the unit sphere in $\boldsymbol{R}^{n}$ ), which diagonalizes $H$ (Theorem 4.2).

Let us make a brief sketch of some well-known results about the spectral representation in the case $n=3$ and $b=0$, i.e., for $H=-\Delta+V$ (see Ikebe [3], [4]). Let $k$ be a non zero vector in $\boldsymbol{R}^{3}$. A generalized eigenfunction $\psi_{k}(x)$ for $-\Delta+V$, which behaves asymptotically like the plane wave $e^{i k \cdot x} \equiv \phi_{k}(x)$, is obtained by solving the Lippmann-Schwinger equation

$$
\begin{equation*}
\psi_{k}(x)=\phi_{k}(x)-\frac{1}{4 \pi} \int_{\mathbf{R}^{3}} \frac{e^{i|k||x-y|}}{|x-y|} V(y) \psi_{k}(y) d y \tag{1.2}
\end{equation*}
$$

The spectral representation for $-\Delta+V$ can be obtained in terms of generalized Fourier transforms

$$
\mathscr{F} f(k)=(2 \pi)^{-3 / 2} \int_{\mathbf{R}^{3}} \overline{\psi_{k}(x)} f(x) d x
$$

However, this procedure works only for $V$ which decays faster than $|x|^{-2-\delta}$ for some positive $\delta$. Agmon [1] has used a version of the Lippmann-Schwinger equation to construct the generalized eigenfunctions and has obtained the spectral representation in the case of short-range $V$ (i.e. $V(x)=O\left(|x|^{-1-\delta}\right)$ ). On the other hand, in the case of long-range $V$ which satisfies $V(x)=O\left(|x|^{-1 / 2-\delta}\right)$, Ikebe [4] has obtained the spectral representation by considering the following limit, instead of using the generalized eigenfunctions explicitly:

$$
\begin{equation*}
\underset{m \rightarrow \infty}{\mathrm{~s}-\lim _{m}} r_{m}^{(n-1) / 2} e^{-i K\left(r_{m} \cdot, \lambda\right)} R(\lambda+i 0) f\left(r_{m} \cdot\right) \tag{1.3}
\end{equation*}
$$

(strong limit in $L_{2}(\Omega)$ ), where $\left\{r_{m}\right\}$ is a sequence tending to infinity as $m \rightarrow \infty, r_{m}$. stands for $r_{m} \omega(\omega \in \Omega), K(x, \lambda)$ is a real-valued function which behaves like $\lambda|x|$ at infinity and $R(\lambda+i 0)$ denotes the boundary value of the resolvent of $H$ on the upper side of the positive real axis (see Theorem 2.3 for details).

The spectral representation for Schrödinger operators with long-range potentials have been investigated by several authors since [4] (e.g. Isozaki [7], Saitō [10]). But it seems that, except for the case of constant magnetic field, the spectral representation for $H$ with magnetic vector potentials has not been studied yet.

In the present paper, $K(x, \lambda)$ is of the form $\lambda|x|-A(x)$, where $A(x)$ is a certain function depending only on $b(x)$, which will be constructed in $\S 2$. This function $A(x)$ has been utilized in Kuroda [8] and, as noticed there, is closely related to the gauge transformation, which changes the magnetic potential $b$ into $b-\operatorname{grad} A$, but does not change the magnetic field. Our assumption ( $b$ ) implies that $A(x)$ can be chosen so that $|b(x)-\operatorname{grad} A(x)| \leqq|x|^{-1 / 2-\delta}$ for some positive $\delta$.
$\S 2$ is a preliminary section including the construction of $A(x)$ and the limiting absorption theorem. In $\S 3$, we study the asymptotic behavior of $R(\lambda+i 0) f$, that is, the existence of the limit (1.3) for any sequence $\left\{r_{m}\right\}$ tending to infinity. For this purpose, we need further the following assumption:
$\left(V^{\prime}\right)$ In addition to $(V), V(x)$ satisfies $|V(x)| \leqq C_{0}(1+|x|)^{-3 / 2-\delta}$ for all $x \in \boldsymbol{R}^{n}$. Theorem 3.9 asserts that the limit (1.3) exists for $f \in L_{2,1}$ (i.e. $\left.(1+|x|) f(x) \in L_{2}\left(\boldsymbol{R}^{n}\right)\right)$ without taking subsequences if the assumptions $\left(V^{\prime}\right)$ and $(b)$ are fulfilled.

It must be noted that, for obtaining our final result, the spectral representation theorem, it suffices to show that the limit (1.3) exists for certain specified sequences $\left\{r_{m}\right\}$. This is, in fact, what we are going to do in $\S 4$ under the assumptions ( $V$ ) and ( $b^{\prime}$ ).

## 2. Preliminaries.

Throughout the paper we use the following notations:

$$
\begin{aligned}
& B_{r}=\left\{x \in \boldsymbol{R}^{n}| | x \mid<r\right\}, \\
& E_{r}=\left\{x \in \boldsymbol{R}^{n}| | x \mid>r\right\}, \\
& S_{r}=\left\{x \in \boldsymbol{R}^{n}| | x \mid=r\right\}, \quad(r>0) .
\end{aligned}
$$

For $\alpha \in \boldsymbol{R}$ and a domain $G \subset \boldsymbol{R}^{n}$, let $L_{2, \alpha}(G)$ denote the Hilbert space of all measurable functions over $G$ such that

$$
\|u\|_{\alpha, G}^{2}=\int_{G}(1+|x|)^{2 \alpha}|u(x)|^{2} d x<\infty .
$$

The $L_{2}$ inner product over $G$ will be denoted by

$$
(u, v)_{G}=\int_{G} u(x) \overline{v(x)} d x
$$

which makes sense if $u \in L_{2, \alpha}(G)$ and $v \in L_{2, \beta}(G)$ with $\alpha+\beta \geqq 0$. When $u$ and $v$ are vector-valued, we also write

$$
\begin{aligned}
& (u, v)_{G}=\sum_{j=1}^{n}\left(u_{j}, v_{j}\right)_{G}, \\
& \|u\|_{\alpha, G}^{2}=\sum_{j=1}^{n}\left\|u_{j}\right\|_{\alpha, G}^{2} .
\end{aligned}
$$

If $\alpha=0$ or if $G=\boldsymbol{R}^{n}$, the subscript $\alpha$ or $G$ will be omitted.
$H_{2, \text { loc }}$ is the set of all locally $L_{2}$ functions on $\boldsymbol{R}^{n}$ with locally $L_{2}$ distribution derivatives up to the second order.

Let $\partial_{j}=\partial / \partial x_{j}(j=1, \ldots, n), \operatorname{grad} u=\left(\partial_{1} u, \ldots, \partial_{n} u\right), r=|x|, \tilde{x}=x / r=\left(\tilde{x}_{1}, \ldots, \tilde{x}_{n}\right)$, $\langle x, y\rangle=\sum_{j=1}^{n} x_{j} y_{j}$ and $\partial_{r} u=\langle\tilde{x}, \operatorname{grad} u\rangle$. Then, for $z \in \boldsymbol{C} \backslash \boldsymbol{R}^{+}$and $\lambda \in \boldsymbol{R}^{+}\left(\boldsymbol{R}^{+}\right.$is the set of all positive real numbers), several differential operators are defined as follows:

$$
\begin{aligned}
& D_{j}=\partial_{j}+i b_{j}, \\
& D u=\left(D_{1} u, \ldots, D_{n} u\right)=\operatorname{grad} u+i b u, \\
& D_{r} u=\langle\tilde{x}, D u\rangle=\partial_{r} u+i\langle\tilde{x}, b\rangle u, \\
& D_{T} u=D u-\tilde{x} D_{r} u, \\
& \mathscr{D}(z) u=D u+\left(\frac{n-1}{2 r}-i \sqrt{z}\right) \tilde{x} u, \\
& \mathscr{D}_{r}(z) u=\langle\tilde{x}, \mathscr{D}(z) u\rangle, \\
& \mathscr{D}_{ \pm} u=\mathscr{D}(\lambda \pm i 0) u=D u+\frac{n-1}{2 r} \tilde{x} u \mp i \sqrt{\lambda} \tilde{x} u, \\
& \mathscr{D}_{ \pm, r} u=\left\langle\tilde{x}, \mathscr{D}_{ \pm} u\right\rangle=D_{r} u+\frac{n-1}{2 r} u \mp i \sqrt{\lambda} u,
\end{aligned}
$$

where $\sqrt{z}$ is the square root of $z$ such that $\mathscr{I}_{m} \sqrt{z}>0\left(\mathscr{I}_{m}=\right.$ the imaginary part $)$.
Proposition 2.1. Suppose the assumption (b) is satisfied. Define

$$
A(x)=\sum_{j=1}^{n} \int_{0}^{|x|} b_{j}(s \tilde{x}) \tilde{x}_{j} d s \quad\left(x \in \boldsymbol{R}^{n}\right) .
$$

Then $A$ is a $C^{2}$ function and the following assertions hold:
(1) There exists a constant $C$ such that $\left|D\left(e^{-i A}\right)\right|=|b-\operatorname{grad} A| \leqq C(1+|x|)^{-1 / 2-\delta}$.
(2) If the assumption ( $b^{\prime}$ ) is satisfied, then (1) holds and, in addition, there exists a constant $C$ such that $\left|L_{0}\left(e^{-i A}\right)\right| \leqq C(1+|x|)^{-1-\delta}$, where $L_{0}=-\sum_{j=1}^{n} D_{j}^{2}=-$ $\sum_{j=1}^{n}\left(\partial_{j}-i b_{j}\right)^{2}$.

Proof. Since $A(x)=\sum_{k=1}^{n} x_{k} \int_{0}^{1} b_{k}(t x) d t, A(x)$ is a $C^{2}$ function and, differentiating this by $x_{j}$ and integrating by parts, we obtain

$$
\begin{aligned}
\partial_{j} A(x) & =\int_{0}^{1} b_{j}(t x) d t+\sum_{k=1}^{n} x_{k} \int_{0}^{1}\left(\partial_{k} b_{j}(t x)+B_{j k}(t x)\right) t d t \\
& =\int_{0}^{1} b_{j}(t x) d t+\int_{0}^{1} t \frac{d}{d t}\left(b_{j}(t x)\right) d t+\int_{0}^{1} \Phi_{j}(t x) d t \\
& =b_{j}(x)+\int_{0}^{1} \Phi_{j}(t x) d t,
\end{aligned}
$$

where $\Phi_{j}(x)=\sum_{k=1}^{n} x_{k} B_{j k}(x)$. From the assumption (b) [resp. ( $\left.b^{\prime}\right)$ ], we obtain the following estimate:

$$
\begin{aligned}
& \left|\Phi_{j}(x)\right| \leqq C(1+|x|)^{-1 / 2-\delta} \\
& {\left[\left|\partial_{j} \Phi_{j}(x)\right|=\left|\sum_{k} x_{k}\left(\partial_{j} B_{j k}\right)(x)\right| \leqq C(1+|x|)^{-1-\delta}\right] .}
\end{aligned}
$$

(Here we have used $B_{j j}=0$.) Hence we have

$$
\begin{aligned}
&\left|\partial_{j} A-b_{j}\right| \leqq \int_{0}^{1}\left|\Phi_{j}(t x)\right| d t \\
& \leqq C \int_{0}^{1}(1+t|x|)^{-1 / 2-\delta} d t \\
&=C \frac{1}{|x|} \int_{0}^{|x|}(1+s)^{-1 / 2-\delta} d s \leqq C^{\prime}(1+|x|)^{-1 / 2-\delta} \\
& {\left[\left|\partial_{j}\left(\partial_{j} A-b_{j}\right)\right|=\left|\int_{0}^{1}\left(\partial_{j} \Phi_{j}(t x)\right) t d t\right| \leqq C^{\prime}(1+|x|)^{-1-\delta}\right] . }
\end{aligned}
$$

Consequently, we obtain the required inequalities by noting that

$$
\begin{aligned}
& D\left(e^{-i A}\right)=i(b-\operatorname{grad} A) e^{-i A}, \\
& L_{0}\left(e^{-i A}\right)=\sum_{j}\left\{\left(\partial_{j} A-b_{j}\right)^{2}-i \partial_{j}\left(\partial_{j} A-b_{j}\right)\right\} e^{-i A} .
\end{aligned}
$$

Proposition 2.2. Suppose that the assumption (b) is satisfied. Let $\lambda$ be a positive number, $\phi$ a smooth function on $\Omega$ (the unit sphere in $\boldsymbol{R}^{n}$ ) and $\rho_{0}(r)$ a smooth function such that $\rho_{0}(r)=1(r \geqq 1)$ and $\rho_{0}(r)=0(r \leqq 1 / 2)$. For $x \in$ $\boldsymbol{R}^{n}\left(r=|x|\right.$ and $\left.\omega=\frac{x}{r}\right)$, let $v_{\phi}(x, \lambda)$ be defined by

$$
\begin{equation*}
v_{\phi}(x, \lambda)=C(\lambda)^{-1} r^{-(n-1) / 2} e^{i(\sqrt{\lambda} r-\lambda(r \omega))} \phi(\omega) \rho_{0}(r), \tag{2.1}
\end{equation*}
$$

where $A$ is as in Proposition 2.1 and $C(\lambda)=\pi^{-1 / 2} \lambda^{1 / 4}$. Then the following assertions hold with a constant $C$ which can be taken uniformly bounded when $\lambda$ varies in a compact set in $\boldsymbol{R}^{+}$:

$$
\begin{align*}
& \left|v_{\phi}(r \omega)\right| \leqq C|\phi(\omega)| r^{-(n-1) / 2}  \tag{1}\\
& \left|\mathscr{D}_{+} v_{\phi}(r \omega)\right| \leqq C\left(r^{-1 / 2-\delta}|\phi(\omega)|+r^{-1}\left|\operatorname{grad}_{\Omega} \phi\right|\right) r^{-(n-1) / 2}
\end{align*}
$$

where $\operatorname{grad}_{\Omega}$ denotes the gradient on $\Omega: \operatorname{grad}_{\Omega}=r\left(\operatorname{grad}-\tilde{x} \partial_{r}\right)$. Hence $v_{\phi} \in$ $L_{2,-(1+\varepsilon) / 2}, \mathscr{D}_{+} v_{\phi} \in L_{2}\left(E_{1}\right)$ for an arbitrary positive $\varepsilon$.
(2) If the assumption ( $b^{\prime}$ ) is satisfied, then (1) holds and in addition

$$
\left|\left(L_{0}-\lambda\right) v_{\phi}(r \omega)\right| \leqq C\left(|\phi|+\left|\operatorname{grad}_{\Omega} \phi\right|+|\Lambda \phi|\right) r^{-\frac{n-1}{2}-1-\delta},
$$

where $\Lambda$ is the Laplace-Beltrami operator on $\Omega$. Hence $\left(L_{0}-\lambda\right) v_{\phi} \in L_{2,(1+\varepsilon) / 2}$ for sufficiently small $\varepsilon(0<\varepsilon<\delta / 2)$.

Proof. First note that, by simple calculation,

$$
\begin{align*}
& \mathscr{D}_{+}\left(r^{-(n-1) / 2} e^{i \sqrt{\lambda} r}\{\cdot\}\right)=r^{-(n-1) / 2} e^{i \sqrt{\lambda} r} D\{\cdot\},  \tag{2.2}\\
& D\left(e^{-i A}\right)=i(b-\operatorname{grad} A) e^{-i A},
\end{align*}
$$

and, since $A$ is so constructed that $\partial_{r} A=\langle\tilde{x}, b\rangle$,

$$
\begin{equation*}
\left\langle\tilde{x}, D\left(e^{-i A}\right)\right\rangle=D_{r}\left(e^{-i A}\right)=0 \tag{2.3}
\end{equation*}
$$

The estimate for $v_{\phi}$ in (1) is immediate from (2.1). We have for $r \geqq 1$, using (2.2),

$$
\begin{align*}
C(\lambda) \mathscr{D}_{+} v_{\phi} & =r^{-(n-1) / 2} e^{i \sqrt{\lambda} r} D\left(e^{-i A} \phi\right)  \tag{2.4}\\
& =r^{-(n-1) / 2} e^{i \sqrt{\lambda} r}\left\{D\left(e^{-i A}\right) \phi+\frac{e^{-i A}}{r} \operatorname{grad}_{\Omega} \phi\right\} .
\end{align*}
$$

Hence, the estimate for $\mathscr{D}_{+} v_{\phi}$ in (1) follows from (2.4) combined with (1) of Proposition 2.1.

Let $v_{0, \phi}$ denote the function $r^{-(n-1) / 2} e^{i \sqrt{\lambda} r} \phi(\omega)$. Then we have by direct computation

$$
\begin{aligned}
& \operatorname{grad} v_{0, \phi}=\left(-\frac{n-1}{2 r}+i \sqrt{\lambda}\right) \tilde{x} v_{0, \phi}+r^{-(n-1) / 2} e^{i \sqrt{\lambda} \frac{1}{r}} \operatorname{grad}_{\Omega} \phi, \\
& (\Delta+\lambda) v_{0, \phi}=-\frac{(n-1)(n-3)}{4 r^{2}} v_{0, \phi}-r^{-(n-1) / 2} e^{i \sqrt{\lambda} r} \frac{\Lambda \phi}{r^{2}}
\end{aligned}
$$

Hence, by noting $v_{\phi}=C(\lambda)^{-1} e^{-i A} v_{0, \phi}$ and by the use of (2.3), we obtain

$$
\begin{align*}
C(\lambda)\left(L_{0}-\lambda\right) v_{\phi}= & v_{0, \phi} L_{0}\left(e^{-i A}\right)-2\left\langle\operatorname{grad} v_{0, \phi}, D\left(e^{-i A}\right)\right\rangle  \tag{2.5}\\
& \quad\left(\Delta v_{0, \phi}+\lambda v_{0, \phi}\right) e^{-i A} \\
=r^{-(n-1) / 2} e^{i \sqrt{\lambda} r}[\{ & \left\{\begin{array}{l}
\left.-i A \frac{(n-1)(n-3)}{r^{2}}+L_{0}\left(e^{-i A}\right)\right\} \phi- \\
\\
\\
\left.\quad-\frac{2 i}{r}\left\langle\operatorname{grad}_{\Omega} \phi, D\left(e^{-i A}\right)\right\rangle+\frac{\Lambda \phi}{r^{2}} e^{-i A}\right] .
\end{array}\right.
\end{align*}
$$

The required estimate in (2) follows from (2.5) combined with Proposition 2.1.
The following theorem which has been established in Ikebe-Saito [6] is fundamental to this paper and will be stated without proof. In what follows $\varepsilon$ will denote a positive constant smaller than $\delta / 2$.

Theorem 2.3 (Limiting absorption principle). Let the assumptions ( $V$ ) and (b) be satisfied. Then the following assertions hold:
(1) Let $K$ be a bounded domain in $\boldsymbol{C} \backslash \boldsymbol{R}$ such that $\bar{K}$, the closure of $K$ in $\boldsymbol{C}$, does not intersect $(-\infty, 0]$, and let $f \in L_{2 .(1+\varepsilon) / 2}$. Then, if we denote the resolvent of $H$ by $R(z)(z \in K), u \equiv R(z) f \in L_{2} \subset L_{2,-(1+\varepsilon) / 2}$ satisfies the inequalities

$$
\begin{aligned}
& \|u\|_{-(1+\varepsilon) / 2} \leqq C\|f\|_{(1+\varepsilon) / 2}, \\
& \|\mathscr{D}(z) u\|_{(-1+\varepsilon) / 2, E_{1}} \leqq C\|f\|_{(1+\varepsilon) / 2},
\end{aligned}
$$

where $C$ is a domain constant independent of $f$.
(2) $R(z) f$ is continuous in $L_{2,-(1+\varepsilon) / 2}$ with respect to $z \in \boldsymbol{C} \backslash \boldsymbol{R}$ and $f \in L_{2,(1+\varepsilon) / 2}$, and for any $\lambda>0$, the limit

$$
R(\lambda \pm i 0) f=\underset{\substack{\mathrm{s}-\lim _{\begin{subarray}{c}{z \rightarrow \lambda \\
\pm=z>0} }} R(z) f} \\
{\hline=2>0}\end{subarray}}{ }
$$

exists in $L_{2,-(1+\varepsilon) / 2}$ in such a way that $R(z) f$ can be extended to a continuous map from $\boldsymbol{C}^{ \pm} \cup \boldsymbol{R}^{+}\left(\boldsymbol{C}^{ \pm}=\left\{z \in \boldsymbol{C} \mid \pm \mathscr{I}_{m z}>0\right\}\right)$ to $L_{2,-(1+\varepsilon) / 2}$. The inequlaities in (1) are satisfied with $u=R(\lambda \pm i 0) f$ and $\mathscr{D}(z)$ replaced by $\mathscr{D}(\lambda \pm i 0)$ when $\lambda \in \bar{K}$.
(3) Given $\lambda>0$ and $f \in L_{2,(1+\varepsilon) / 2}, R(\lambda \pm i 0) f$ in (2) solves the following problem uniquely:

$$
\left\{\begin{array}{l}
L u-\lambda u=f, \quad u \in H_{2, l o c} \cap L_{2,-(1+\varepsilon) / 2}  \tag{2.6}\\
\mathscr{D}(\lambda \pm i 0) u \in L_{2,(-1+\varepsilon) / 2}\left(E_{1}\right) .
\end{array}\right.
$$

(4) For.f, $g \in L_{2,(1+\varepsilon) / 2}$ and any Borel set $B$ in $\boldsymbol{R}^{+}$, we have

$$
\begin{aligned}
(E(B) f, g) & =\frac{1}{2 \pi i} \int_{B}(R(\lambda+i 0) f-R(\lambda-i 0) f, g) d \lambda \\
& =\frac{1}{2 \pi i} \int_{B}\{(R(\lambda+i 0) f, g)-(f, R(\lambda+i 0) g)\} d \lambda,
\end{aligned}
$$

where $E$ is the spectral measure associated with $H$. The part of $H$ in $E((0, \infty))$ is absolutely continuous.

Remark. When $\mathscr{D}_{+} u \in L_{2,(-1+\varepsilon) / 2}\left[\mathscr{D}_{-} u \in L_{2,(-1+\varepsilon) / 2}\right], u$ is said to satisfy the outgoing [resp. incoming] radiation condition, and $R(\lambda+i 0) f[R(\lambda-i 0) f]$ is called the outgoing [resp. incoming] solution for the equation $L u-\lambda u=f$. For example, if the assumptions $(V)$ and $\left(b^{\prime}\right)$ are fulfilled and if $g_{\phi}$ denotes $(L-\lambda) v_{\phi}$, where $v_{\phi}$ is a function as in Proposition 2.2, $v_{\phi}$ is the outgoing solution for $(L-\lambda) v_{\phi}=g_{\phi}$ (thus $\left.R(\lambda+i 0) g_{\phi}=v_{\phi}\right)$, since $v_{\phi} \in L_{2,-(1+\varepsilon) / 2}, \quad \mathscr{D}_{+} v_{\phi} \in L_{2}\left(E_{1}\right)$ and $g_{\phi}=\left(L_{0}-\lambda\right) v_{\phi}+V v_{\phi} \in$ $L_{2,(1+\varepsilon) / 2}$ as noticed in Proposition 2.2.

## 3. Asymptotic behavior of outgoing solutions

As has been seen in (4) of Theorem 2.3, the following quantity is important in investigating the spectral representation for $H$ :

$$
\begin{equation*}
(R(\lambda+i 0) f-R(\lambda-i 0) f, f) \quad \text { for } \quad f \in L_{2,(1+\varepsilon) / 2} . \tag{3.1}
\end{equation*}
$$

We shall utilize the following Green's formula for computing (3.1):

$$
\begin{equation*}
\int_{B_{r}}(u \bar{f}-f \bar{u}) d x=\int_{S_{r}}\left[\left(\mathscr{D}_{+, r} u\right) \bar{u}-u\left(\overline{\mathscr{D}_{+, r} u}\right)\right] d S+2 i \sqrt{\lambda} \int_{S_{r}}|u|^{2} d S, \tag{3.2}
\end{equation*}
$$

where $u=R(\lambda+i 0) f$. The left-hand side of (3.2) converges to (3.1) as $r \rightarrow \infty$. As remarked after Theorem 2.3, $v_{\phi}$ in Proposition 2.2 is the outgoing solution for $(L-\lambda) v_{\phi}=g_{\phi}$ under the assumptions ( $V$ ) and ( $b^{\prime}$ ), and, by letting $f=g_{\phi}$ and $r \rightarrow \infty$ in (3.2), we have

$$
\begin{equation*}
\frac{1}{2 \pi i}\left(R(\lambda+i 0) g_{\phi}-R(\lambda-i 0) g_{\phi}, g_{\phi}\right)=\lim _{r \rightarrow \infty} \frac{\sqrt{\lambda}}{\pi} \int_{S_{r}}\left|v_{\phi}\right|^{2} d S=\|\phi\|_{L_{2}(\Omega)}^{2} \tag{3.3}
\end{equation*}
$$

where $\left\|\|_{L_{2}(\Omega)}\right.$ is the norm of $L_{2}(\Omega)\left(=\right.$ the $L_{2}$ space over the unit sphere $\left.\Omega\right)$. In the next lemma and the succeeding propositions, we are going to prove that $R(\lambda+i 0) f$ behaves like $v_{\phi}$ near infinity and the analogue of (3.3) holds for $f$ satisfying an appropriate condition under the assumptions ( $V$ ) and (b).

Definition 3.1. Let the operator $\mathscr{F}(\lambda, r): L_{2,(1+\varepsilon) / 2} \rightarrow L_{2}(\Omega)$ be defined by

$$
\mathscr{F}(\lambda, r) f(\omega)=C(\lambda) r^{(n-1) / 2} e^{-i \sqrt{\lambda} r+i A(r \omega)} R(\lambda+i 0) f(r \omega),
$$

where $\omega \in \Omega, f \in L_{2,(1+\varepsilon) / 2}$ and $C(\lambda)=\pi^{-1 / 2} \lambda^{1 / 4}$.
Lemma 3.2. Suppose that the assumptions $(V)$ and (b) are satisfied, $f \in$ $L_{2,(1+\varepsilon) / 2}$ and $\mathscr{D}(\lambda+i 0) R(\lambda+i 0) f \in L_{2}\left(E_{1}\right)$. Then there exists the following strong limit in $L_{2}(\Omega)$ :

$$
\mathscr{F}(\lambda ; f)=\lim _{r \rightarrow \infty} \mathscr{F}(\lambda, r) f .
$$

For the proof of this lemma, we need some formulae and propositions. To begin with, we consider:

$$
\begin{equation*}
\Gamma_{ \pm, r}(u, v) \equiv \int_{S_{r}}\left(\mathscr{D}_{ \pm, r} u\right) \bar{v} d S \tag{3.4}
\end{equation*}
$$

where $u, v \in H_{2, \text { loc }}$. First, we obtain

$$
\begin{align*}
& \frac{d}{d r} \Gamma_{ \pm, r}(u, v)  \tag{3.5}\\
& \quad=\frac{d}{d r}\left(r^{(n-1) / 2} e^{ \pm i \sqrt{\lambda} r+i A}\left(\mathscr{D}_{ \pm, r} u\right), \quad r^{(n-1) / 2} e^{ \pm i \sqrt{\lambda} r+i A} v\right)_{L_{2}(\Omega)} \\
& \quad=\int_{S_{r}}\left(\mathscr{D}_{\mp, r} \mathscr{D}_{ \pm, r} u\right) \bar{v} d S+\int_{S_{r}}\left(\mathscr{D}_{ \pm, r} u\right)\left(\mathscr{D}_{\mp, r}\right) d S,
\end{align*}
$$

where we have used

$$
\partial_{r}\left(r^{(n-1) / 2} e^{ \pm i \sqrt{\lambda} r+i \Lambda} u\right)=r^{(n-1) / 2} e^{ \pm i \sqrt{\lambda} r+i A} \mathscr{D}_{\mp, r} u .
$$

Moreover, we have

$$
\begin{align*}
& \mathscr{D}_{\mp_{,}, r} \mathscr{D}_{ \pm, r}=D_{r}^{2}+\frac{n-1}{r}-D_{r}+\frac{(n-1)(n-3)}{4 r^{2}}+\lambda,  \tag{3.6}\\
& \int_{S_{r}}\left(L_{0} u+D_{r}^{2} u+\frac{n-1}{r} D_{r} u\right) \bar{v} d S  \tag{3.7}\\
& \quad=\int_{S_{r}}\langle D u, \overline{D v}\rangle d S-\int_{S_{r}} D_{r} u \overline{D_{r} v} d S \\
& \quad=\int_{S_{r}}\left\langle D_{T} u, \overline{D_{T} v}\right\rangle d S,
\end{align*}
$$

where $L_{0}=-\sum_{j} D_{j}^{2}:(3.6)$ is obtained by straightforward calculation, and (3.7) is obtained by differentiating Green's formula

$$
\left(L_{0} u, v\right)_{B_{r}}=(D u, D v)_{B_{r}}-\int_{S_{r}}\left(D_{r} u\right) \bar{v} d S
$$

with respect to $r$ and by noting that $D=D_{T}+\tilde{x} D_{r}$ is an orthogonal sum decomposition. By the use of (3.5), (3.6) and (3.7) we obtain

$$
\begin{gather*}
\frac{d}{d r} \Gamma_{ \pm, r}(u, v)=\int_{S_{r}}\left\langle D_{T} u, \overline{D_{T} v}\right\rangle d S+\int_{S_{r}}(\tilde{V} u-(L-i) u) \bar{v} d S  \tag{3.8}\\
\\
+\int_{S_{r}}\left(\mathscr{D}_{ \pm, r} u\right)\left(\mathscr{D}_{\mp, r} v\right) d S
\end{gather*}
$$

where $\tilde{V}=V+\frac{(n-1)(n-3)}{4 r^{2}}$ Further, taking the upper side of the double sign of (3.8) and using the relation $\mathscr{D}_{-, r}=\mathscr{D}_{+, r}+2 \sqrt{\bar{i} i}$ and the orthogonal sum decomposition $\mathscr{D}_{+}=D_{T}+\tilde{x} \mathscr{D}_{+, r}$, we have

$$
\begin{equation*}
\left(\frac{d}{d r}+2 \sqrt{\bar{\lambda} i}\right) \Gamma_{+, r}(u, v)=\int_{S_{r}}\left\langle\mathscr{D}_{+} u, \overline{\mathscr{D}_{+} v}\right\rangle d S+\int_{S_{r}}(\bar{V} u-(L-\lambda) u) \bar{v} d S \tag{3.9}
\end{equation*}
$$

Proposition 3.3. Let $u \in L_{2,-(1+\varepsilon) / 2}$ and $\mathscr{D}_{+} u \in L_{2,(-1+\eta) / 2}\left(E_{1}\right)$. Then there exists a sequence $\left\{r_{m}\right\}$ of positive numbers diverging to infinity as $m \rightarrow \infty$ such that

$$
\begin{align*}
& r_{m}^{-\varepsilon} \int_{S_{r_{m}}}|u|^{2} d S \longrightarrow 0 \quad \text { as } \quad m \longrightarrow \infty  \tag{3.10}\\
& r_{m}^{\eta} \int_{S_{r_{m}}}\left|\mathscr{D}_{+} u\right|^{2} d S \longrightarrow 0 \text { as } m \longrightarrow \infty \tag{3.11}
\end{align*}
$$

Proof. It is not difficult to verify that, for an integrable function $g$ over $\boldsymbol{R}^{n}$, there exists a sequence $\left\{r_{m}\right\}$ diverging to infinity as $m \rightarrow \infty$ such that

$$
r_{m} \int_{S_{r_{m}}}|g(x)| d S \longrightarrow 0
$$

Considering this fact, we have the assertion of the proposition.
Proposition 3.4. Let the assumption of lemma 3.2 be satisfied. Let $v \in$ $L_{2,-(1+\varepsilon) / 2} \cap H_{2, \text { loc }}$ with $\mathscr{D}_{+} v \in L_{2}\left(E_{1}\right)$. Then we have

$$
\begin{equation*}
\int_{S_{r}}\left(\mathscr{D}_{+, r} u\right) \bar{v} d S \longrightarrow 0 \quad \text { as } \quad r \longrightarrow \infty \quad(u=R(\lambda+i 0) f) \tag{3.12}
\end{equation*}
$$

In particular, we have by (3.12) with $v=u$, (3.2) and Theorem 2.3

$$
\begin{equation*}
\lim _{r \rightarrow \infty}\|\mathscr{F}(\lambda, r) f\|_{L_{2}(\Omega)}^{2}=\frac{1}{2 \pi i}(R(i+i 0) f-R(i-i 0) f, f) \tag{3.13}
\end{equation*}
$$

Proof. Let $F(r)$ be the right-hand side of (3.9). Then, $F(r)$ is integrable over $(1, \infty)$ because by assumption $\mathscr{D}_{+} u, \mathscr{D}_{+} v \in L_{2}\left(E_{1}\right), \tilde{V} u,(L-\lambda) u=f \in L_{2,(1+\varepsilon) / 2}$, $v \in L_{2,-(1+\varepsilon) / 2}$. Since we can rewrite (3.9) as

$$
\frac{d}{d r}\left(e^{2 i \sqrt{\lambda} r} \Gamma_{+, r}\right)=e^{2 i \sqrt{\lambda} r} F(r)
$$

we have

$$
\begin{equation*}
e^{2 i \sqrt{\lambda} r} \Gamma_{+, r}=-\int_{r}^{\infty} e^{2 i \sqrt{\lambda} r} F(r) d r . \tag{3.14}
\end{equation*}
$$

Here we have used the existence of a sequence $\left\{r_{m}\right\}$ such that $r_{m} \rightarrow \infty$ and $\Gamma_{+, r_{m} \rightarrow 0}$ as $m \rightarrow \infty$, which can be verified by an argument similar to the proof of Proposition 3.3, since $r^{-(1+\varepsilon) / 2}\left(\mathscr{D}_{+, r} u\right) \bar{v}$ is integrable over $E_{1}$ and $\Gamma_{+, r}=\int_{S_{r}}\left(\mathscr{D}_{+, r} u\right) \bar{v} d S((3.4))$. (3.12) follows from (3.4) and (3.14).

Proposition 3.5. Under the assumption of Lemma 3.2, there exists a function $\alpha(r)$ such that $\alpha(r) \downarrow 0$ as $r \rightarrow \infty$, and for all $r$ and s satisfying $1 \leqq r \leqq s$ and for $\phi \in$ $H_{1}(\Omega)\left(=\right.$ the set of $L_{2}$ functions over $\Omega$ with $L_{2}$ distribution derivatives up to the first order), the following inequality holds:

$$
\begin{aligned}
& \left|\left(w_{+}(r), \phi\right)_{L_{2}(\Omega)}\right|+\left|\left(w_{-}(r)-w_{-}(s), \phi\right)_{L_{2}(\Omega)}\right| \\
& \quad \leqq \alpha(r)\left(\|\phi\|_{L_{2}(\Omega)}+r^{-1 / 2}\left\|\operatorname{grad}_{\Omega} \phi\right\|_{L_{2}(\Omega)}\right),
\end{aligned}
$$

where $w_{ \pm}(r) \in L_{2}(\Omega)$ are defined by

$$
\begin{aligned}
& w_{ \pm}(r)(\cdot)=C(\lambda) r^{(n-1) / 2} e^{-i, ~ \hat{\lambda} r+i A(r \cdot)} \mathscr{D}_{ \pm, r} u(r \cdot) \\
& \left(u=R(\lambda+i 0) f, \quad C(\lambda)=\pi^{-1 / 2} \lambda^{1 / 4}\right) .
\end{aligned}
$$

Proof. Let $\phi \in C^{\infty}(\Omega), v_{\phi}$ be as in Proposition 2.2 and $F(r)$ be the right-hand side of (3.9) with $v$ replaced by $v_{\phi}$, i.e.,

$$
\begin{equation*}
F(r)=\int_{S_{r}}\left\langle\mathscr{D}_{+} u, \mathscr{D}_{+} v_{\phi}\right\rangle d S+\int_{S_{r}}(\tilde{V} u-f) \overline{v_{\phi}} d S . \tag{3.15}
\end{equation*}
$$

Then, since $v_{\phi} \in L_{2,-(1+\varepsilon) / 2} \cap H_{2, \text { loc }}, \mathscr{D}_{+} v_{\phi} \in L_{2}\left(E_{1}\right)$ as noted in Proposition 2.2, the argument in the proof of Proposition 3.4 is applicable to the case $v=v_{\phi}$. That is, $F(r)$ is integrable over $(1, \infty)$, and we have

$$
\begin{equation*}
e^{2 i \sqrt{ } / r} \Gamma_{+, r}\left(u, v_{\phi}\right)=-\int_{r}^{\infty} e^{2 i \sqrt{\lambda} r} F(r) d r . \tag{3.16}
\end{equation*}
$$

Moreover, taking the lower side of the double sign of (3.8) and replacing $v$ by $v_{\phi}$, we have

$$
\begin{align*}
& \frac{d}{d r} \Gamma_{-, r}\left(u, v_{\phi}\right)=\int_{S_{r}}\left\langle D_{T} u, \overline{D_{T}} \bar{v}_{\phi}\right\rangle d S+\int_{S_{r}}\left(\mathscr{D}_{-, r} u\right)\left(\overline{\mathscr{D}_{+, r} v_{\phi}}\right) d S  \tag{3.17}\\
&+\int_{S_{r}}(\tilde{V} u-f) \overline{v_{\phi}} d S .
\end{align*}
$$

The right-hand side of (3.17) coincides with $F(r)$ because we have $\mathscr{D}_{+,} v_{\phi}=0$ by straightforward calculation and the orthogonal sum decomposition $\mathscr{D}_{+}=D_{T}+$ $\tilde{x} \mathscr{D}_{+, r}$. Hence, we obtain by integrating (3.17),

$$
\begin{equation*}
\Gamma_{-, s}\left(u, v_{\phi}\right)-\Gamma_{-, r}\left(u, v_{\phi}\right)=\int_{r}^{s} F(r) d r \quad(1 \leqq r \leqq s) . \tag{3.18}
\end{equation*}
$$

On the other hand, (3.4) and the definition of $w_{ \pm}(r)$ and $v_{\phi}$ leads to

$$
\begin{equation*}
\left(w_{ \pm}(r), \phi\right)_{L_{2}(\Omega)}=C(\lambda)^{2} \Gamma_{ \pm, r}\left(u, v_{\phi}\right), \tag{3.19}
\end{equation*}
$$

from which, in view of (3.16) and (3.18), the following estimate can be obtained for $1 \leqq r \leqq s:$

$$
\begin{aligned}
& \left|\left(w_{+}(r), \phi\right)_{L_{2}(\Omega)}\right|+\left|\left(w_{-}(r)-w_{-}(s), \phi\right)_{L_{2}(\Omega)}\right| \\
& \quad \leqq C \int_{r}^{\infty}|F(r)| d r \\
& \leqq C \int_{r}^{\infty} d r \int_{S_{r}}\left|\left\langle\mathscr{D}_{+} u, \overline{\mathscr{D}_{+} v_{\phi}}\right\rangle\right| d S+C \int_{r}^{\infty} d r \int_{S_{r}}\left|\bar{V} u \overline{v_{\phi}}\right| d S+C \int_{r}^{\infty} d r \int_{S_{r}}\left|f \overline{v_{\phi}}\right| d S \\
& \equiv \quad+\quad I_{2} \quad+\quad I_{3},
\end{aligned}
$$

where $C$ is a constant depending only on $\lambda$. According to Schwarz' inequality and Proposition 2.2 (1), we have

$$
\begin{aligned}
& I_{1} \leqq C \int_{r}^{\infty} r^{(n-1) / 2}\left\|\mathscr{D}_{+} u\right\|_{L_{2}(\Omega)}\left(r^{-1 / 2-\delta}\|\phi\|_{L_{2}(\Omega)}+r^{-1}\left\|\operatorname{grad}_{\Omega} \phi\right\|_{L_{2}(\Omega)}\right) d r \\
& \leqq C \sqrt{\int_{r}^{\infty} r^{n-1}\left\|\mathscr{D}_{+} u\right\|_{L_{2}(\Omega)}^{2} d r \times} \\
& \quad \times\left(\sqrt{\left.\int_{r}^{\infty} r^{-1-2 \delta} d r\|\phi\|_{L_{2}(\Omega)}+\sqrt{\int_{r}^{\infty}} r^{-2} d r\left\|\operatorname{grad}_{\Omega} \phi\right\|_{L_{2}(\Omega)}\right)}\right. \\
&=C\left\|\mathscr{D}_{+} u\right\|_{E_{r}}\left(\frac{r^{-\delta}}{\sqrt{2 \delta}}\|\phi\|_{L_{2}(\Omega)}+r^{-1 / 2}\left\|\operatorname{grad}_{\Omega} \phi\right\|_{L_{2}(\Omega)}\right), \\
& I_{2} \leqq C \int_{r}^{\infty} r^{(n-1) / 2} r^{-1-\delta}\|u\|_{L_{2}(\Omega)} d r\|\phi\|_{L_{2}(\Omega)} \\
& \leqq C r^{-\varepsilon / 2}\|u\|_{-(1+\varepsilon) / 2}\|\phi\|_{L_{2}(\Omega)}, \\
& I_{3} \leqq C \int_{r}^{\infty} r^{(n-1) / 2}\|f\|_{L_{2}(\Omega)} d r\|\phi\|_{L_{2}(\Omega)} \\
& \leqq C r^{-\varepsilon / 2}\|f\|_{(1+\varepsilon) / 2}\|\phi\|_{L_{2}(\Omega)} .
\end{aligned}
$$

Then, if we put $\alpha(r)=C\left\|\mathscr{D}_{+} u\right\|_{E_{r}}+C r^{-\varepsilon / 2}\left(\|u\|_{-(1+\varepsilon) / 2}+\|f\|_{(1+\varepsilon) / 2}\right)$, the required inequality holds for $\phi \in C^{\infty}(\Omega)$, and $\alpha(r) \downarrow 0$ as $r \rightarrow \infty$ since $\mathscr{D}_{+} u=\mathscr{D}_{+} R(\lambda+i 0) f \in$ $L_{2}\left(E_{1}\right)$ by assumption. Finally, the inequality is obtained for $\phi \in H_{1}(\Omega)$ by approximating it by smooth functions.

Proposition 3.6. Let the assumptions $(V)$ and $(b)$ be satisfied and $f \in L_{2,(1+\varepsilon) / 2}$. Then, $\left\{r_{m}^{-1 / 2} \operatorname{grad}_{\Omega} \mathscr{F}\left(\lambda, r_{m}\right) f\right\}_{m}$ is bounded in $L_{2}(\Omega)$ for any sequence $\left\{r_{m}\right\}$ which satisfies (3.10), (3.11) with $u=R(\lambda+i 0) f$ and $\eta=1$.

Proof. By straightforward computation, we have

$$
\begin{aligned}
& \frac{1}{C(\lambda)} \operatorname{grad}_{\Omega} \mathscr{F}(\lambda, r) f \\
& \quad=r^{(n+1) / 2} e^{-i \sqrt{\bar{\lambda}} r}\left\{\operatorname{grad}\left(e^{i A} u\right)-\tilde{x} \partial_{r}\left(e^{i A} u\right)\right\} \\
& =r^{(n+1) / 2} e^{-i \sqrt{\bar{\lambda} r+i A}\left\{D u-i(b-\operatorname{grad} A) u-\tilde{x} D_{r} u\right\}} \\
& =r^{(n+1) / 2} e^{-i \sqrt{\lambda} r+i A}\left\{D_{T} u-i(b-\operatorname{grad} A) u\right\} .
\end{aligned}
$$

Hence, by using Proposition 2.1 (1), we have

$$
\begin{aligned}
& \left\|\operatorname{grad}_{\Omega} \mathscr{F}\left(\lambda, r_{m}\right) f\right\|_{L_{2}(\Omega)} \\
& \quad \leqq C\left\{r_{m}^{(n+1) / 2}\left\|D_{T} u\left(r_{m} \cdot\right)\right\|_{L_{2}(\Omega)}+r_{m}^{n / 2-\delta}\left\|u\left(r_{m} \cdot\right)\right\|_{L_{2}(\Omega)}\right\} \\
& \quad \leqq C r_{m}^{1 / 2}\left\{\sqrt{r_{m}} \int_{S_{r_{m}}}\left|\mathscr{D}_{+} u\right|^{2} d S+\sqrt{r_{m}^{-2 \delta}} \int_{S_{r_{m}}}|u|^{2} d S\right.
\end{aligned} .
$$

Consequently, since $\varepsilon$ denotes a positive constant smaller than $\delta / 2$, we have the assertion of the proposition.

Proof of Leimma 3.2. Seeing that $\mathscr{F}(\lambda, r) f=(2 i \sqrt{\lambda})^{-1}\left(w_{-}(r)-w_{+}(r)\right)$, we have by Proposition 3.5

$$
\begin{align*}
& \left|(\mathscr{F}(\lambda, r) f-\mathscr{F}(\lambda, s) f, \phi)_{L_{2}(\Omega)}\right|  \tag{3.20}\\
& \quad \leqq \lambda^{-1 / 2} \alpha(r)\left(\|\phi\|_{L_{2}(\Omega)}+r^{-1 / 2}\left\|\operatorname{grad}_{\Omega} \phi\right\|_{L_{2}(\Omega)}\right) .
\end{align*}
$$

for $\phi \in H_{1}(\Omega)$ and $r, s(1 \leqq r \leqq s)$. Consequently, $(\mathscr{F}(\lambda, r) f, \phi)_{L_{2}(\Omega)}$ is convergent when $r \rightarrow \infty$. Since $\{\mathscr{F}(\lambda, r) f\}$ is bounded with respect to $r$ (Proposition 3.4) and $H_{1}(\Omega)$ is dense in $L_{2}(\Omega)$, we have the weak convergence in $L_{2}(\Omega)$ of $\{\mathscr{F}(\lambda, r) f\}$. Let $\mathscr{F}(\lambda ; f)$ denote this weak limit. Then, letting $s \rightarrow \infty$ in (3.20), we have

$$
\begin{align*}
& \left|(\mathscr{F}(\lambda, r) f-\mathscr{F}(\lambda ; f), \phi)_{L_{2}(\Omega)}\right|  \tag{3.21}\\
& \quad \leqq \lambda^{-1 / 2} \alpha(r)\left(\|\phi\|_{L_{2}(\Omega)}+r^{-1 / 2}\left\|\operatorname{grad}_{\Omega} \phi\right\|_{L_{2}(\Omega)}\right) .
\end{align*}
$$

Putting $r=r_{m}$ and $\phi=\mathscr{F}\left(\lambda, r_{m}\right) f$ in (3.21) and using the boundedness of $\{\mathscr{F}(\lambda, r) f\}$ (Proposition 3.4) and Proposition 3.6 for a sequence $\left\{r_{m}\right\}$ which satifises (3.10) and (3.11) with $u=R(\lambda+i 0) f$ and $\eta=1$ (the existence of such a sequence is guaranteed by Proposition 3.3 since $u \in L_{2,-(1+\varepsilon) / 2}$ by Theorem 2.3 and $\mathscr{D}_{+} u \in L_{2}\left(E_{1}\right)$ by assumption), one can see that there exists a constant $M$ independent of $m$ such that the following inequality holds:

$$
\left|\left(\mathscr{F}\left(\lambda, r_{m}\right) f-\mathscr{F}(\lambda ; f), \mathscr{F}\left(\lambda, r_{m}\right) f\right)\right| \leqq \alpha\left(r_{m}\right) M .
$$

According to this and the weak convergence of $\{\mathscr{F}(\lambda, r) f\}$ to $\mathscr{F}(\lambda ; f)$, we have

$$
\lim _{m \rightarrow \infty}\left\|\mathscr{F}\left(\lambda, r_{m}\right) f\right\|_{L_{2}(\Omega)}^{2}=\|\mathscr{F}(\lambda ; f)\|_{L_{2}(\Omega)}^{2} .
$$

for some sequence $\left\{r_{m}\right\}$, Since $\|\mathscr{F}(i, r) f\|_{L_{2}(\Omega)}$ converges when $r \rightarrow \infty$ by Proposition 3.4, we have further

$$
\lim _{r \rightarrow \infty}\|\mathscr{F}(\lambda, r) f\|_{L_{2}(\Omega)}^{2}=\|\mathscr{F}(\lambda ; f)\|_{L_{L_{2}}(\Omega)}^{2} .
$$

From this and the weak convergence, the strong convergence of $\mathscr{F}(\lambda, r) f$ to $\mathscr{F}(\lambda ; f)$ follows.

Lemma 3.7. Suppose that the assumptions $\left(V^{\prime}\right)$ and (b) are satisfied. Let $f \in L_{2,1} \subset L_{2,(1+\varepsilon) / 2}$. Then, in addition to all the statements of Theorem 2.3, we have

$$
\begin{aligned}
& \mathscr{D}(\lambda \pm i 0) R(\lambda \pm i 0) f \in L_{2}\left(E_{1}\right) \\
& \|\mathscr{D}(\lambda \pm i 0) R(\lambda \pm i 0) f\|_{E_{1}} \leqq C\|f\|_{1},
\end{aligned}
$$

where $C$ is a constant independent of $f$ and remains bounded when $i$ varies in $a$ compact set of $\boldsymbol{R}^{+}$.

For the proof, we need the following proposition which is a version of Lemma 2.1 in Ikebe-Saitō [6].

Proposition 3.8. Assume that $V$ is a bounded measurable function and $b_{j}$ are continuously differentiable. Then, for any $u \in H_{2, \text { loc }}$ and $R>0$, the following inequality holds with a positive constant $C$ independent of $u$ and $R$ :

$$
\begin{equation*}
\|D u\|_{B_{R}}^{2} \leqq C\left(\|u\|_{B_{2 R}}^{2}+\|L u\|_{B_{2 R}}^{2}\right) . \tag{3.22}
\end{equation*}
$$

Proof. For $\psi$, a real-valued smooth function with compact support, we have by partial integration

$$
\begin{aligned}
\|D(\psi u)\|^{2}= & -\left(\psi u, \sum_{j} D_{j} D_{j}(\psi u)\right) \\
= & -\left(\psi u, \psi\left(\sum_{j} D_{j} D_{j} u\right)\right)-2(\psi u,\langle\operatorname{grad} \psi, D u\rangle)-(\psi u,(\Delta \psi) u) \\
= & (\psi u, \psi L u)-(\psi u, \psi V u)-2(u \operatorname{grad} \psi, D(\psi u))+ \\
& \quad+2\|u \operatorname{grad} \psi\|^{2}-(\psi u,(\Delta \psi) u) .
\end{aligned}
$$

Hence, noting that $2|(u \operatorname{grad} \psi, D(\psi u))| \leqq 2\|u \operatorname{grad} \psi\|^{2}+\frac{1}{2}\|D(\psi u)\|^{2}$ by the use of Schwarz' inequality, we have

$$
\frac{1}{2}\|D(\psi u)\|^{2} \leqq|(\psi u, \psi L u)|+|(\psi u, \psi V u)|+4\|u \operatorname{grad} \psi\|^{2}+|(\psi u,(\Delta \psi) u)| .
$$

(3.22) is obtained from this inequality by taking $\psi(x)=\rho(x / R)$ where $\rho$ is a smooth function on $\boldsymbol{R}^{n}$ such that $\rho(x)=1(|x| \leqq 1)$ and $\rho(x)=0(|x| \geqq 2)$.

Proof of Lemma 3.7. Let $\left\{z_{m}\right\}$ be a sequence of complex numbers such that $\mathscr{I}_{m} z_{m}>0$ and $z_{m} \rightarrow \lambda$ as $m \rightarrow \infty$. Then for $u_{m} \equiv R\left(z_{m}\right) f \in L_{2}$ we have by Theorem 2.3

$$
\begin{align*}
& \left\|u_{m}\right\|_{-(1+\varepsilon) / 2} \leqq C\|f\|_{(1+\varepsilon) / 2} \leqq C\|f\|_{1},  \tag{3.23}\\
& u_{m} \longrightarrow u=R(\lambda+i 0) f \text { strongly in } L_{2,-(1+\varepsilon) / 2} . \tag{3.24}
\end{align*}
$$

Since $f=\left(L-z_{m}\right) u_{m}$, applying Proposition 3.8 with $u=u_{m}$, we have

$$
\begin{equation*}
\int_{1<|x|<R}\left|\mathscr{D}\left(z_{m}\right) u_{m}\right|^{2} d x \leqq C\left\{\|f\|_{B_{2 R}}^{2}+\left\|u_{m}\right\|_{B_{2 R}}^{2}\right\} \tag{3.25}
\end{equation*}
$$

with $C$ independent of $R, f$ and $m$. Since $u_{m} \in L_{2}$ and $f \in L_{2,1}, \mathscr{D}\left(z_{m}\right) u_{m} \in L_{2}\left(E_{1}\right)$ by (3.25). Let $\rho$ be a smooth function of $r=|x|$ on $\boldsymbol{R}^{n}$ such that $\rho(1)=0$ and $\rho(r)=r(r \geqq 2)$. Then, by multiplying $\left(L-z_{m}\right) u_{m}=f$ by $\rho \mathscr{D}_{r}\left(z_{m}\right) u_{m}$ and integrating by parts, the following equality can be obtained:

$$
\begin{aligned}
& \int_{E_{1}}\left[\left\{\left(\mathscr{I} m \sqrt{z_{m}}\right) \rho+\frac{1}{2} \frac{\partial \rho}{\partial r}\right\}\left|\mathscr{D}\left(z_{m}\right) u_{m}\right|^{2}+\right. \\
& \left.\quad+\left(\frac{\rho}{r}-\frac{\partial \rho}{\partial r}\right)\left(\left|\mathscr{D}\left(z_{m}\right) u_{m}\right|^{2}-\left|\mathscr{D}_{r}\left(z_{m}\right) u_{m}\right|^{2}\right)\right] d x \\
& =-\mathscr{R}_{e}\left[\int_{E_{1}} \rho \tilde{V} u_{m} \mathscr{\mathscr { D }}_{r}\left(z_{m}\right) u_{m} d x\right]+\mathscr{R}_{e}\left[\int_{E_{1}} \rho f \overline{\mathscr{O}_{r}}\left(z_{m}\right) u_{m} d x\right]- \\
& \quad-\mathscr{I}_{m}\left[\int_{E_{1}} \sum_{j \cdot k=1}^{n} \rho B_{j k} \mathscr{D}_{j}\left(z_{m}\right) u_{m} \tilde{x}_{k} \overline{u_{m}} d x\right] .
\end{aligned}
$$

(For the details of the computation see Ikebe-Saitō [6].) From this equality, noting that $\mathscr{I}_{m} \sqrt{z_{m}}>0, \frac{\partial \rho}{\partial r}=1$ and $\frac{\rho}{r}-\frac{\partial \rho}{\partial r}=0$ if $r \geqq 2,\left|\mathscr{D}_{r}\left(z_{m}\right) u_{m}\right|^{2} \leqq\left|\mathscr{D}\left(z_{m}\right) u_{m}\right|^{2}$, $\rho \widetilde{V}=O\left(|x|^{-1 / 2-\delta}\right)$ and $\rho B_{j k}=O\left(|x|^{-1 / 2-\delta}\right)$, we have by the use of Schwarz' inequality

$$
\begin{aligned}
& \left\|\mathscr{D}\left(z_{m}\right) u_{m}\right\|_{E_{1}}^{2} \\
& =\int_{1<|x|<2}\left|\mathscr{D}\left(z_{m}\right) u_{m}\right|^{2} d x+\int_{E_{2}}\left|\mathscr{D}\left(z_{m}\right) u_{m}\right|^{2} d x \\
& \leqq \int_{1<|x|<2}\left|\mathscr{D}\left(z_{m}\right) u_{m}\right|^{2} d x+2 \int_{E_{2}}\left(\left(\mathscr{I} m \sqrt{z_{m}}\right) \rho+\frac{1}{2} \frac{\partial \rho}{\partial r}\right)\left|\mathscr{D}\left(z_{m}\right) u_{m}\right|^{2} d x \\
& \leqq C \int_{1<|x|<2}\left|\mathscr{D}\left(z_{m}\right) u_{m}\right|^{2} d x+C\left\|(1+|x|)^{-1 / 2-\delta} u_{m}\right\|\left\|\mathscr{D}\left(z_{m}\right) u_{m}\right\|_{E_{1}}+ \\
& \quad+C\|(1+|x|) f\|\left\|\mathscr{D}\left(z_{m}\right) u_{m}\right\|_{E_{1}} \\
& \leqq C \int_{1<|x|<2}\left|\mathscr{D}\left(z_{m}\right) u_{m}\right|^{2} d x+C^{\prime}\left(\left\|u_{m}\right\|^{2}-(1+\varepsilon) / 2+\|f\|_{1}^{2}\right)+ \\
& \quad+\frac{1}{2}\left\|\mathscr{D}\left(z_{m}\right) u_{m}\right\|_{E_{1}}^{2},
\end{aligned}
$$

where $C, C^{\prime}$ are constants depending only on $n, \rho$ and $C_{0}$ in the assumptions ( $V^{\prime}$ ) and (b). Hence, applying (3.25) with $R=2$ and using (3.23), we have

$$
\left\|\mathscr{D}\left(z_{m}\right) u_{m}\right\|_{E_{1}}^{2} \leqq C^{\prime \prime}\|f\|_{1}^{2}
$$

Consequently, there exists a subsequence $\left\{m_{j}\right\}$ such that $\left\{\mathscr{D}\left(z_{m_{j}}\right) u_{m_{j}}\right\}$ is weakly convergent in $L_{2}\left(E_{1}\right)$ to some $w \in L_{2}\left(E_{1}\right)$. On the other hand, according to (3.24), $\left\{\mathscr{D}\left(z_{m}\right) u_{m}\right\}$ converges to $\mathscr{D}_{+} u$ in the distributional sense. Therefore, $\mathscr{D}_{+} u$ coincides with $w$ and we have

$$
\left\|\mathscr{D}_{+} u\right\|_{E_{1}} \leqq \liminf _{j \rightarrow \infty}\left\|\mathscr{D}\left(z_{m_{j}}\right) u_{m_{j}}\right\|_{E_{1}} \leqq C^{\prime \prime}\|f\|_{1} .
$$

We have thus concluded the proof of Lemma 3.7.
We conclude this section with a theorem which can be obtained by combining Lemma 3.7 with Lemma 3.2.

Theorem 3.9. Suppose that the assumptions $\left(V^{\prime}\right)$ and $(b)$ are satisfied. Then, there exists an operator $\mathscr{F}(\lambda): L_{2,1} \rightarrow L_{2}(\Omega)(\lambda>0)$ such that the following assertions hold:

$$
\begin{align*}
& \mathscr{F}(i) f=s-\lim _{r \rightarrow \infty} \mathscr{F}(\lambda, r) f \quad \text { in } \quad L_{2}(\Omega) \quad\left(f \in L_{2,1}\right)  \tag{1}\\
& \|\mathscr{F}(\lambda) f\|_{L_{2}(\Omega)}^{2}=\frac{1}{2 \pi i}(R(\lambda+i 0) f-R(\lambda-i 0) f, f) \quad\left(f \in L_{2,1}\right) . \tag{2}
\end{align*}
$$

Proof. Under the assumptions $\left(V^{\prime}\right)$ and $(b), f \in L_{2,1}$ satisfies the assumption of Lemma 3.2 for every positive number $\lambda$ according to Lemma 3.7. Hence, we can define $\mathscr{F}(\lambda) f$ by $\mathscr{F}(\lambda ; f)=\lim _{r \rightarrow \infty} \mathscr{F}(\lambda, r) f$ of Lemma 3.2 which obviously satisfies $(1)$. (2) follows from (1) and (3.13) of Proposition 3.4.

## 4. Spectral representation for $\mathbf{H}$

In this section, the spectral representation for $H$ is obtained by means of the next lemma, a version of Theorem 3.9, where we impose a stronger condition on but, instead, relax the condition on $V$.

Lemma 4.1. Suppose that the assumptions (V) and ( $b^{\prime}$ ) are satisfied. Then there exists a bounded operator $\mathscr{F}(\lambda): L_{2,(1+\varepsilon) / 2} \rightarrow L_{2}(\Omega)(\lambda>0)$ such that the following assertions hold:

$$
\begin{equation*}
\mathscr{F}(\lambda) f=\underset{m \rightarrow \infty}{s-\lim _{m \rightarrow \infty} \mathscr{F}\left(\lambda, r_{m}\right) f \quad\left(f \in L_{2,(1+\varepsilon) / 2}\right), ~, ~} \tag{1}
\end{equation*}
$$

where $\left\{r_{m}\right\}$ is any sequence satisfying

$$
\left\{\begin{array}{l}
r_{m} \longrightarrow \infty  \tag{4.1}\\
r_{m}^{-c} \int_{S_{r_{m}}}|R(\lambda+i 0) f|^{2} d S \longrightarrow 0 \\
r_{m}^{\varepsilon} \int_{S_{r_{m}}}|\mathscr{D}(\lambda+i 0) R(\lambda+i 0) f|^{2} d S \longrightarrow 0 \quad \text { as } \quad m \longrightarrow \infty
\end{array}\right.
$$

$$
\begin{equation*}
\|\mathscr{F}(\lambda) f\|_{L_{2}(\Omega)}^{2}=\frac{1}{2 \pi i}(R(\lambda+i 0) f-R(\lambda-i 0) f, f) \quad\left(f \in L_{2,(1+\varepsilon) / 2}\right) \tag{2}
\end{equation*}
$$

(3) $\mathscr{F}(\lambda) f$ is strongly continuous in $\lambda$ for any $f \in L_{2,(1+\varepsilon) / 2}$.

Proof. Let $f \in L_{2,(1+\varepsilon) / 2}$. By Theorem 2.3, $u \equiv R(\lambda+i 0) f \in L_{2,-(1+\varepsilon) / 2}$ and $\mathscr{D}_{+} u \in L_{2,(-1+\varepsilon) / 2}$. Therefore, the existence of a sequence $\left\{r_{m}\right\}$ satisfying (4.1) is guaranteed by Proposition 3.3 with $\eta=\varepsilon$.

Let $\phi \in C^{\infty}(\Omega)$ and $v_{\phi}=v_{\phi}(x, \lambda)$ be defined as in Proposition 2.2. Then we have by Green's formula

$$
\begin{align*}
& \left((L-\lambda) u, v_{\phi}\right)_{B_{r}}-\left(u,(L-\lambda) v_{\phi}\right)_{B_{r}}  \tag{4.2}\\
& \quad=\int_{S_{r}} u \overline{\mathscr{D}_{+, r}} v_{\phi} d S-\int_{S_{r}} \mathscr{D}_{+, r} u \overline{v_{\phi}} d S+2 i \sqrt{\lambda} \int_{S_{r}} u \bar{v}_{\phi} d S
\end{align*}
$$

Hence, since $(\mathscr{F}(\lambda, r) f, \phi)_{L_{2}(\Omega)}=C(i)^{2} \int_{S_{r}} u \overline{v_{\phi}} d S\left(C(\lambda)=\pi^{-1 / 2} \lambda^{1 / 4}\right)$ in view of the definition of $v_{\phi}$ and $\mathscr{F}(\lambda, r) f((2.1)$ and Definition 3.1), we obtain from (4.2) with $r=r_{m}$

$$
\left.\begin{array}{rl}
\left(\mathscr{F}\left(\lambda, r_{m}\right) f, \phi\right)_{L_{2}(\Omega)}=\frac{1}{2 \pi i} & \left\{\int_{S_{r_{m}}} u \mathscr{D}_{+, r} v_{\phi}\right. \tag{4.3}
\end{array} d S-\int_{S_{r_{m}}} \mathscr{D}_{+, r} u \overline{v_{\phi}} d S-\right] .
$$

From (4.1) and the estimates in Proposition 2.2, we have

$$
\begin{equation*}
\int_{S_{r_{m}}} u \overline{\mathscr{D}_{+, r} v_{\phi}} d S \longrightarrow 0, \quad \int_{S_{r_{m}}} \mathscr{D}_{+, r} u \overline{v_{\phi}} d S \longrightarrow 0 \quad \text { as } \quad m \longrightarrow \infty \tag{4.4}
\end{equation*}
$$

Because $u \in L_{2,-(1+\varepsilon) / 2}, \quad(L-\lambda) u=f \in L_{2,(1+\varepsilon) / 2}, \quad v_{\phi} \in L_{2,-(1+\varepsilon) / 2} \quad$ and $(L-\lambda) v_{\phi} \in$ $L_{2,(1+\varepsilon) / 2}$ (Proposition 2.2), we have in view of (4.4) the following equality by letting $m \rightarrow \infty$ in (4.3):

$$
\begin{equation*}
\lim _{m \rightarrow \infty}\left(\mathscr{F}\left(\lambda, r_{m}\right) f, \phi\right)_{L_{2}(\Omega)}=\frac{1}{2 \pi i}\left\{-\left(f, v_{\phi}\right)+\left(u,(L-\lambda) v_{\phi}\right)\right\} \tag{4.5}
\end{equation*}
$$

Similarly, by letting $r=r_{m}$ and $m \rightarrow \infty$ in (3.2), we have

$$
\begin{equation*}
\lim _{m \rightarrow \infty}\left\|\mathscr{F}\left(\lambda, r_{m}\right) f\right\|_{L_{2}(\Omega)}^{2}=-\frac{1}{2 \pi i}(R(\lambda+i 0) f-R(\lambda-i 0) f, f) \tag{4.6}
\end{equation*}
$$

Therefore, $\left\{\mathscr{F}\left(\lambda, r_{m}\right) f\right\}_{m}$ is bounded in $L_{2}(\Omega)$. Hence, since (4.5) holds for $\phi \in C^{\infty}(\Omega)$, which is a dense subspace of $L_{2}(\Omega)$, the weak convegence of $\left\{\mathscr{F}\left(\lambda, r_{m}\right) f\right\}$ in $L_{2}(\Omega)$ follows. Note that this weak limit is independent of the choice of the sequence $\left\{r_{m}\right\}$ because the right-hand side of $(4.5)$ is so.

Let an operator $\mathscr{F}(\lambda): L_{2,(1+\varepsilon) / 2} \rightarrow L_{2}(\Omega)$ be defined by
for $f \in L_{2,(1+\varepsilon) / 2}$, where $\left\{r_{m}\right\}$ is any sequence satisfying (4.1). Then we have by (4.6) and Theorem 2.3 (2)

$$
\begin{equation*}
\|\mathscr{F}(\lambda) f\|_{L_{2}(\Omega)}^{2} \leqq \frac{1}{2 \pi i}(R(\lambda+i 0) f-R(\lambda-i 0) f, f) \leqq C\|f\|_{(1+\varepsilon) / 2}^{2} \tag{4.8}
\end{equation*}
$$

where $C$ is a constant independent of $f$ which can be taken uniformly bounded when $\lambda$ varies in a compact set in $\boldsymbol{R}^{+}$. Hence $\mathscr{F}(\lambda)$ is a bounded operator: $L_{2,(1+\varepsilon) / 2} \rightarrow L_{2}(\Omega)$.

Let us prove that (4.7) is a strong limit.
Let $H_{0}$ be the self-adjoint realization of $L_{0}=-\sum_{j=1}^{n} D_{j}^{2}$ in $\mathscr{H}=L_{2}\left(\boldsymbol{R}^{n}\right)$. Then the argument developed so far is applicable to $H_{0}$. That is, if $R_{0}(\lambda+i 0)$ denotes the
boundary value of the resolvent of $H_{0}$ and $\mathscr{F}_{0}(\lambda, r)$ is defined as $\mathscr{F}(\lambda, r)$ in Definition 3.1 with $R(\lambda+i 0)$ replaced by $R_{0}(\lambda+i 0)$, we have

$$
\begin{align*}
& \lim _{m \rightarrow \infty}\left\|\mathscr{F}_{0}\left(\lambda, r_{m}\right) f^{\prime}\right\|_{L_{2}(\Omega)}^{2}=\frac{1}{2 \pi i}\left(R_{0}(\lambda+i 0) f-R_{0}(\lambda-i 0) f, f\right), \\
& \mathscr{F}_{0}(\lambda) f=\underset{m \rightarrow \infty}{\mathrm{w}-\lim _{\boldsymbol{\prime}}} \mathscr{F}_{0}\left(\lambda, r_{m}\right) f \tag{}
\end{align*}
$$

for $f \in L_{2,(1+\varepsilon) / 2}$ and for a sequence $\left\{r_{m}\right\}$ satisfying (4.1) with $R(\lambda+i 0)$ replaced by $R_{0}(\lambda+i 0)$, and $\mathscr{F}_{0}(\lambda)$ is a bounded operator: $L_{2,(1+\varepsilon) / 2} \rightarrow L_{2}(\Omega)$. Since Theorem 3.9 applies to $H_{0},\left\{\mathscr{F}_{0}(\lambda, r) f\right\}$ converges strongly in $L_{2}(\Omega)$ when $r \rightarrow \infty$ if $f \in L_{2,1}$. Thus, from (4.6') and (4.7'), we have

$$
\begin{equation*}
\left\|\mathscr{F}_{0}(\lambda) f\right\|_{L_{2}(\Omega)}^{2}=\frac{1}{2 \pi i}\left(R_{0}(\lambda+i 0) f-R_{0}(\lambda-i 0) f, f\right) \tag{4.9}
\end{equation*}
$$

for $f \in L_{2,1}$. But since both sides of (4.9) are continuous in $f \in L_{2,(1+\varepsilon) / 2}$ and since $L_{2,1}$ is dense in $L_{2,(1+\varepsilon) / 2}$, (4.9) holds for all $f \in L_{2,(1+\varepsilon) / 2}$. (4.9) combined with $\left(4.6^{\prime}\right),\left(4.7^{\prime}\right)$ leads to the strong convergence of $\left\{\mathscr{F}_{0}\left(\lambda, r_{m}\right) f\right\}$ to $\mathscr{F}_{0}(\lambda) f$ for $f \in$ $L_{2,(1+\varepsilon) / 2}$.

Next, noting that $H=H_{0}+V$ and $V$ is a bounded linear operator: $L_{2,-(1+\varepsilon) / 2} \rightarrow$ $L_{2,(1+e) / 2}$ by the assumption ( $V$ ), we have, passing to the limit in the usual resolvent equation,

$$
\begin{equation*}
R(\lambda+i 0)-R_{0}(\lambda+i 0)=-R_{0}(\lambda+i 0) V R(\lambda+i 0) . \tag{4.10}
\end{equation*}
$$

Hence, applying (410) to $f \in L_{i,(1+\varepsilon) / 2}$, we have the following equality for $f \in L_{2,(1+\varepsilon) / 2}$ and $u=R(\lambda+i 0) f \in L_{2,-(1+\varepsilon) / 2}$ :

$$
u=R_{0}(\lambda+i 0)(f-V u) .
$$

From the definition of $\mathscr{F}_{0}(\lambda, r)$ and $\mathscr{F}(\lambda, r)$, we have

$$
\mathscr{F}(\lambda, r) f=\mathscr{F}_{0}(\lambda, r)(f-V u) .
$$

Noting that the condition (4.1) on $\left\{r_{m}\right\}$ is concerned only with $u=R(\lambda+i 0) f=$ $R_{0}(\lambda+i 0)(f-V u)$ and applying the result obtained for $H_{0}$ to $f-V u \in L_{2,(1+\varepsilon) / 2}$, we have the strong convergence in $L_{2}(\Omega)$ of the sequence $\left\{\mathscr{F}\left(\lambda, r_{m}\right) f\right\}=\left\{\mathscr{F}_{0}\left(\lambda, r_{m}\right)\right.$ ( $f-V u$ ) $\}$ for $f \in L_{2,(1+e) / 2}$ and for any sequence $\left\{r_{m}\right\}$ satisfying (4.1). Thus, we have proved that the weak limit (4.7) is also a strong limit. From (4.6) and the strong convergence of $\left\{\mathscr{F}\left(\lambda, r_{m}\right) f\right\}$, (2) of the lemma follows. Thus we have proved (1) and (2) of the lemma.

Finally, for obtaining the continuity in $\lambda$ of $\mathscr{F}(\lambda) f$, it suffices to show the continuity of $(\mathscr{F}(\lambda) f, \phi)_{L_{2}(\Omega)}$ for all $\phi \in C^{\infty}(\Omega)$ since $\|\mathscr{F}(\lambda) f\|_{L_{2}(\Omega)}$ is continuous in $\lambda$ as is seen from the right-hand side of (2) of the lemma and Theorem 2.3. We have by (4.5)

$$
\begin{align*}
& (\mathscr{F}(\lambda) f, \phi)_{L_{2}(\Omega)}  \tag{4.11}\\
& \quad=\frac{1}{2 \pi i}\left\{-\left(f(\cdot), v_{\phi}(\cdot, \lambda)\right)+\left(R(\lambda+i 0) f(\cdot),(L-\lambda) v_{\phi}(\cdot, \lambda)\right)\right\} .
\end{align*}
$$

Let $\lambda$ vary in an interval $\left[a_{1}, a_{2}\right]\left(0<a_{1}<a_{2}<\infty\right)$. Then we have the following pointwise estimates by Proposition 2.2:

$$
\begin{aligned}
& \left|v_{\phi}(x, \lambda)\right| \leqq C(1+|x|)^{-(n-1) / 2} \\
& \left|\left(L_{0}-\lambda\right) v_{\phi}(x, \lambda)\right| \leqq C(1+|x|)^{-(n-1) / 2-1-\delta}
\end{aligned}
$$

with $C$ independent of $\lambda \in\left[a_{1}, a_{2}\right]$. Hence, since $|V(x)| \leqq C_{0}(1+|x|)^{-1-\delta}$ and $L=L_{0}+V$, we have

$$
\begin{aligned}
& (1+|x|)^{-(1+\varepsilon)}\left|v_{\phi}\left(x, \lambda^{\prime}\right)-v_{\phi}(x, \lambda)\right|^{2} \leqq 2 C(1+|x|)^{-n-\varepsilon} \\
& (1+|x|)^{1+\varepsilon}\left|\left(L-\lambda^{\prime}\right) v_{\phi}\left(x, \lambda^{\prime}\right)-(L-\lambda) v_{\phi}(x, \lambda)\right|^{2} \leqq C^{\prime}(1+|x|)^{-n-2 \delta+\varepsilon}
\end{aligned}
$$

for $\lambda, \lambda^{\prime} \in\left[a_{1}, a_{2}\right]$. Therefore, noting that according to (2.1) $v_{\phi}\left(x, \lambda^{\prime}\right) \rightarrow v_{\phi}(x, \lambda)$ and $\left(L-\lambda^{\prime}\right) v_{\phi}\left(x, \lambda^{\prime}\right) \rightarrow(L-\lambda) v_{\phi}(x, \lambda)$ as $\lambda^{\prime} \rightarrow \lambda$ for each $x \in \boldsymbol{R}^{n}$, we have by the use of the Lebesgue dominated convergence theorem

$$
\begin{aligned}
& \left\|v_{\phi}\left(\cdot, \lambda^{\prime}\right)-v_{\phi}(\cdot, \lambda)\right\|_{-(1+\varepsilon) / 2} \longrightarrow 0 \\
& \left\|\left(L-\lambda^{\prime}\right) v_{\phi}\left(\cdot, \lambda^{\prime}\right)-(L-\lambda) v_{\phi}(\cdot, \lambda)\right\|_{(1+\varepsilon) / 2} \longrightarrow 0
\end{aligned}
$$

as $\lambda^{\prime} \rightarrow \lambda\left(\lambda^{\prime}, \lambda \in\left[a_{1}, a_{2}\right]\right)$. Thus, $v_{\phi}(\cdot, \lambda)$ and $(L-\lambda) v_{\phi}(\cdot, \lambda)$ are continuous for $\lambda>0$ in $L_{2,-(1+\varepsilon) / 2}$ and $L_{2,(1+\varepsilon) / 2}$, respectively. Hence, since $R(\lambda+i 0) f$ is continuous in $\lambda$ in $L_{2,-(1+\varepsilon) / 2}$ (Theorem 2.3 (2)), we obtain the continuity of $(\mathscr{F}(\lambda) f$, $\phi)_{L_{2}(\Omega)}$ from (4.11). This completes the proof of (3) and thus Lemma 4.1.

We leave the proof of the next theorem, the spectral representation for $H$, to the reader, because it can be obtained in the same way as in the proof of Theorem 2.8 and Theorem 3.1 of Ikebe [4], by using Theorem 2.3, Lemma 4.1 and Proposition 2.2.

Theorem 4.2. Suppose that the assumptions ( $V$ ) and ( $b^{\prime}$ ) are satisfied. Then the following assertions hold:
(1) Let $P_{a c}=E(0, \infty)$ be the projection onto the subspace $\mathscr{H}_{a c}$ of absolute continuity for $H$. Let $\hat{\mathscr{H}}=L_{2}\left((0, \infty) ; L_{2}(\Omega)\right)$ be the Hilbert space of all $L_{2}(\Omega)$-valued square integrable functions over $(0, \infty)=\boldsymbol{R}^{+}$. For $f \in L_{2,(1+\varepsilon) / 2}$, we define a mapping $\mathscr{F} f: \boldsymbol{R}^{+} \rightarrow L_{2}(\Omega)$ by

$$
\mathscr{F} f(\lambda) \equiv \mathscr{F}(\lambda) f \quad(\lambda>0),
$$

where $\mathscr{F}(\lambda)$ is given in Lemma 4.1. Then for $f, g \in L_{2,(1+\varepsilon) / 2}$ and for any Borel subset $B$ of $\boldsymbol{R}^{+}$, we have

$$
\begin{equation*}
(E(B) f, g)=\int_{B}(\mathscr{F} f(\lambda), \mathscr{F} g(\lambda))_{L_{2}(\Omega)} d \lambda \tag{4.12}
\end{equation*}
$$

where $E$ is the spectral measure for $H$. In particular, by letting $B=\boldsymbol{R}^{+}$in (4.12), we have $\mathscr{F} f \in \stackrel{\mathscr{H}}{ }$ and

$$
\left(P_{a c} f, g\right)=(\mathscr{F} f, \mathscr{F} g)_{\hat{\mathscr{F}}}=\int_{0}^{\infty}(\mathscr{F} f(\lambda), \mathscr{F} g(\lambda))_{L_{2}(\Omega)} d \lambda
$$

(2) The operator $\mathscr{F}$ defined above on $L_{2,(1+\varepsilon) / 2}$ can be uniquely extended to whole $\mathscr{H}$ (this will be denoted by $\mathscr{F}$ also). $\mathscr{F}$ is a partial isometry with the initial set $\mathscr{H}_{a c}$ and the final set $\hat{\mathscr{H}}$ (i.e. $\mathscr{F}$ is an unitary operator from $\mathscr{H}_{\text {uc }}$ to $\hat{\mathscr{H}}$ ). For a bounded Borel measureable function $\alpha(\lambda)$ ) on $\boldsymbol{R}^{+}$, we have for all $f \in \mathscr{H}_{a c}$

$$
(\mathscr{F} \alpha(H) f)(\lambda)=\alpha(\lambda) \mathscr{F} f(\lambda) \quad \text { a.e. } \quad \lambda>0 .
$$

(3) Let B be a relatively compact Borel subset of $\boldsymbol{R}^{+}$. Then $\mathscr{F}_{B}^{*}$ is defined by

$$
\mathscr{F}_{B}^{*} \hat{f}=\int_{\boldsymbol{B}} \mathscr{F}(\lambda)^{*} \hat{f}(\lambda) d \lambda \quad \text { for } \quad \hat{f} \in \hat{\mathscr{H}},
$$

which is a partial isometry from $\hat{\mathscr{H}}$ to $\mathscr{H}_{a c}$ and we have

$$
\mathscr{F}_{B}^{*}=E(B) \mathscr{F}^{*}=(\mathscr{F} E(B))^{*} .
$$

The following inversion formula holds:

$$
P_{a c} f=\mathrm{s}-\lim \int_{N \rightarrow \infty} \int_{1 / N}^{N} \mathscr{F}(\lambda)^{*}(\mathscr{F} f)(\lambda) d \lambda .
$$

(4) $\mathscr{F}(\lambda)^{*}: L_{2}(\Omega) \rightarrow L_{2,-(1+\varepsilon) / 2}$ is an eigenoperator of $H$ with eigenvalue $\lambda$ in the sense that for any $\phi \in L_{2}(\Omega), L \mathscr{F}(\lambda)^{*} \phi=\lambda \mathscr{F}(\lambda)^{*} \phi$ in the distributional sense.

Remark. One can obtain the spectral representation under the assumptions $\left(V^{\prime}\right)$ and $(b)$ on the basis of $\mathscr{F}(\lambda)$ defined in Theorem 3.9 except for the unitraity assertion.

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