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# Elliptic Regularity and Essential Self-adjointness of Dirichlet Operators on $\mathbb{R}^{n}$ 

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One of the classical problems in mathematical physics is the problem of essential self-adjointness for Dirichlet operators

$$
L:=\Delta+\beta \cdot \nabla
$$

with domain $C_{0}^{\infty}\left(\mathbb{R}^{n}\right)$ (:= all infinitely differentiable functions on $\mathbb{R}^{n}$ with compact support) on $L^{2}\left(\mathbb{R}^{n}, \mu\right)$, where $\mu$ is a measure on $\mathbb{R}^{n}$ with density $\rho:=\varphi^{2}$, with $\varphi \in H_{\mathrm{loc}}^{2,1}\left(\mathbb{R}^{n}\right)$ and $\beta:=\nabla \rho / \rho$. (By definition $\beta(x)=0$ if $\rho(x)=0$ ). The results obtained in [1], [8], [9], [11], [14], [25] have been important steps in the investigation of this problem. One motivation to study this problem is that the operator $-L$ is unitary equivalent to the Schrödinger operator $H:=-\Delta+V, V:=\Delta \varphi / \varphi$, considered on $L^{2}\left(\mathbb{R}^{n}, d x\right)$ (see, e.g., [1], [5]) where $d x$ denotes Lebesgue measure on $\mathbb{R}^{n}$. The corresponding isomorphism $L^{2}\left(\mathbb{R}^{n}, \mu\right) \rightarrow L^{2}\left(\mathbb{R}^{n}, d x\right)$ is given by $f \mapsto \varphi \cdot f$. Conversely, if $H=-\Delta+V$ is a Schrödinger operator on $L^{2}\left(\mathbb{R}^{n}, d x\right)$ with lower bounded spectrum $\sigma(H)$ whose minimum is an eigenvalue $E$, then the isomorphism above holds for the potential $V-E$ (and $\varphi:=$ the ground state). Since this unitary equivalence only holds for sufficiently regular $\varphi$, Dirichlet operators are also sometimes called generalized Schrödinger operators. We emphasize that under the above isomorphism in general domains change drastically. Hence known results on the essential self-adjointness of $H$ with domain $C_{0}^{\infty}\left(\mathbb{R}^{n}\right)$ on $L^{2}\left(\mathbb{R}^{n}, d x\right)$ do not apply. On the contrary in many cases the essential self-adjointness of Dirichlet operators implies this property for Schrödinger operators (see e.g. [16, pp. 217, 218]).

There are basically two different types of sufficient conditions known for the essential self-adjointness of Dirichlet operators: global and local. A typical global condition obtained in [14] is: $|\beta| \in L^{4}\left(\mathbb{R}^{n}, \mu\right)$ (provided $\rho>0$ a.e.). The best local condition obtained so far has been found in [25] where $\rho$ has been required to be locally Lipschitzian and strictly positive if $n \geq 2$ (and
with even weaker conditions if $n=1$, cf. Remark 2 below). In particular, this means that $\beta$ is locally bounded. One of our main results in this paper (cf. Theorem 7 below) says that $L$ is essentially self-adjoint provided that $\rho$ is merely locally bounded and locally uniformly positive (cf. below) and $|\beta| \in$ $L_{\mathrm{loc}}^{\gamma}\left(\mathbb{R}^{n}, \mu\right)$ for some $\gamma>n$ (which as we shall show below, is equivalent to $|\beta| \in L_{\mathrm{loc}}^{\gamma}\left(\mathbb{R}^{n}, d x\right)$; cf. Corollary 8 ). The proof of Theorem 7 is based on an elliptic regularity result (which is the first main result of this paper) giving $H_{\mathrm{loc}}^{\gamma, 1}$ regularity of distributional solutions of the elliptic equation $L^{*} F=0$, where $L f:=\Delta f+\langle B, \nabla f\rangle+c f$. This result is formulated as Theorem 1 below. As a consequence one gets $H_{\mathrm{loc}}^{\gamma, 1}$-regularity of invariant measures for diffusion processes with drifts satisfying certain mild local integrability conditions (which extends a result from [3], [4]). Finally, we note that for the above mentioned special applications to Schrödinger operators $H=-\Delta+V$, of course, one still needs corresponding information about the ground state $\varphi$ to ensure that $|\beta|=2|\nabla \varphi / \varphi| \in L_{\mathrm{loc}}^{\gamma}\left(\mathbb{R}^{n} ; \mu\right)$.

Throughout this paper, $\Omega$ is a (fixed) open subset of $\mathbb{R}^{n}$, and for $r \in$ $(-\infty, \infty)$ and $p \geq 1, H_{\text {loc }}^{p, r}(\Omega)$ denotes the class of (generalized) functions $u$ on $\Omega$, such that $(1-\Delta)^{r / 2} \psi u \in L^{p}\left(\mathbb{R}^{n}, d x\right)$ for every $\psi \in C_{0}^{\infty}(\Omega)$. These spaces coincide with the usual Sobolev spaces for integer $r \geq 1$. All properties of these spaces which are needed below can be found, for instance, in [23]. If $v$ is a signed measure, then by definition $\int f d \nu=\int f \chi d|\nu|$, where $\chi:=d \nu / d|\nu|$, and $L^{p}(\Omega, \nu):=L^{p}(\Omega,|\nu|)$. If, in addition, $\nu \ll d x$, then we write $v$ instead of $\frac{d \nu}{d x}$. Furthermore, $\langle$,$\rangle denotes the Euclidean inner product on \mathbb{R}^{n}$ and $|\cdot|$ the corresponding norm.

Theorem 1. Let $n \geq 2$ and let $\mu$, ve (signed) Radon measures on $\Omega$. Let $B=\left(B^{i}\right): \Omega \rightarrow \mathbb{R}^{n}, c: \Omega \rightarrow \mathbb{R}$ be maps such that $|B|, c \in L_{\mathrm{loc}}^{1}(\Omega, \mu)$. Assume that

$$
\begin{equation*}
\int L \varphi(x) \mu(d x)=\int \varphi(x) \nu(d x) \quad \forall \varphi \in C_{0}^{\infty}(\Omega) \tag{1}
\end{equation*}
$$

where

$$
\begin{equation*}
L \varphi(x):=\Delta \varphi(x)+\langle B(x), \nabla \varphi(x)\rangle+c(x) \varphi(x) \tag{2}
\end{equation*}
$$

Then:
(i) $\mu \in H_{\mathrm{loc}}^{p, 1-n(p-1) / p-\varepsilon}(\Omega)$ for any $p \geq 1$ and $\varepsilon>0$. Here $1-n(p-1) / p>0$ if $p \in\left[1, \frac{n}{n-1}\right)$ and, in particular, $\mu$ admits a density $F \in L_{\mathrm{loc}}^{p}(\Omega, d x)$ for any $p \in\left[1, \frac{n}{n-1}\right)$.
(ii) If $|B| \in L_{\mathrm{loc}}^{\gamma}(\Omega, \mu), c \in L_{\mathrm{loc}}^{\gamma / 2}(\Omega, \mu)$ and $v \in L_{\mathrm{loc}}^{n /(n-\gamma+2)}(\Omega, d x)$ where $n \geq$ $\gamma>1$, then $F:=\frac{d \mu}{d x} \in H_{\mathrm{loc}}^{p, 1}(\Omega)$ for any $p \in[1, n /(n-\gamma+1))$. In particular, $F \in L_{\mathrm{loc}}^{p}(\Omega, d x)$ for any $p \in[1, n /(n-\gamma)$ ), where (here and below) $\frac{n}{n-\gamma}:=\infty$, if $\gamma=n$.
(iii) If $\gamma>n$ and either.
(a) $|B| \in L_{\mathrm{loc}}^{\gamma}(\Omega, d x)$ and $c, v \in L_{\mathrm{loc}}^{\gamma n /(n+\gamma)}(\Omega, d x)$,
or
(b) $|B| \in L_{\mathrm{loc}}^{\gamma}(\Omega, \mu), c \in L_{\mathrm{loc}}^{\gamma n /(n+\gamma)}(\Omega, \mu)$, and $v \in L_{\mathrm{loc}}^{\gamma n /(n+\gamma)}(\Omega, d x)$, then $\mu$ admits a density $F \in H_{\mathrm{loc}}^{\gamma, 1}(\Omega)$, and, in particular, $F \in C_{\mathrm{loc}}^{1-n / \gamma}(\Omega)$.

Remark 2 (i) There is a similar regularity result for $n=1$ (whose proof is easier and, in fact, quite elementary). Therefore, Theorem 7 and Corollary 8 below also hold in this case. However, our conditions there for $n=1$ are then obviously equivalent to: $\varphi(=\sqrt{\rho}) \in H_{\mathrm{loc}}^{2,1}(\mathbb{R})$ and (the continuous version of) $\rho$ is strictly positive. But under these conditions in the special case $n=1$ both results are already contained in [25]. So, we state and prove our results only for $n \geq 2$.
(ii) Note that since $c \in L_{\mathrm{loc}}^{1}(\Omega, \mu)$ and $\nu$ is a Radon measure, the assumptions on $c$, $v$ in Theorem 1 (ii) are automatically fulfilled, if $\gamma \leq 2$, provided $v \ll d x$.

To prove Theorem 1 we use the following lemma.
Lemma 3. (i) For any $r \in(-\infty, \infty)$ and $p>1$, if $\Delta u \in H_{\mathrm{loc}}^{p, r}(\Omega)$, then $u \in H_{\mathrm{loc}}^{p, r+2}(\Omega)$; also if $u \in H_{\mathrm{loc}}^{p, r}(\Omega)$, then $u_{x^{i}} \in H_{\mathrm{loc}}^{p, r-1}(\Omega), 1 \leq i \leq n$.
(ii) We have $H_{\mathrm{loc}}^{p, 1}(\Omega) \subset L_{\mathrm{loc}}^{n p /(n-p)}(\Omega, d x)$ and $L_{\mathrm{loc}}^{p}(\Omega, d x) \subset H_{\mathrm{loc}}^{n p /(n-p),-1}(\Omega)$ whenever $1<p<n$, and $H_{\mathrm{loc}}^{p, 1}(\Omega) \subset C_{\mathrm{loc}}^{1-n / p}(\Omega)$ if $p>n$, so that in the latter case elements of $H_{\mathrm{loc}}^{p, 1}(\Omega)$ are locally bounded. Also for $q>p>1, L_{\mathrm{loc}}^{p}(\Omega, d x) \subset$ $H_{\mathrm{loc}}^{q, n / q-n / p}(\Omega)$.
(iii) If $\mu$ is a Radon measure on $\Omega$, then $\mu \in H_{\mathrm{loc}}^{p,-m}(\Omega)$ whenever $p>1$ and $m>n(1-1 / p)$.

Proof. Assertion (i) is well-known. Specifically, its first statement is a well-known elliptic regularity result and the second statement follows from the boundedness of Riesz's transforms. Assertion (ii) is just the Sobolev imbedding theorems (see [23]). Assertion (iii) follows from these imbedding theorems since, for regular sub-domains $U$ of $\Omega, H^{q, m}(U) \subset C(\bar{U})$ if $q m>n$ whence by duality the space $H^{q /(q-1),-m}(U)=\left[H^{q, m}(U)\right]^{*}$ contains all finite measures on $U$.

Proof of Theorem 1. (i): We have that in the sense of distributions

$$
\begin{equation*}
\Delta \mu=\left(B^{i} \mu\right)_{x^{i}}-c \mu+v \tag{3}
\end{equation*}
$$

on $\Omega$. Here by Lemma 3 (iii), the right-hand side belongs to $H_{\text {loc }}^{p,-m-1}(\Omega)$ if $m>n(1-1 / p)$. By Lemma 3 (i) we conclude $\mu \in H_{\mathrm{loc}}^{p,-m+1}(\Omega)$, which leads to the result after substituting $m=n(1-1 / p)+\varepsilon$.

Before we prove (ii), (iii) we need some preparations. Fix a $p_{1}>1$ and assume that $F:=\frac{d \mu}{d x} \in L_{\mathrm{loc}}^{p_{1}}(\Omega, d x)$. (Such $p_{1}$ exists by (i).) Define

$$
\begin{equation*}
r:=r\left(p_{1}\right):=\frac{\gamma p_{1}}{\gamma-1+p_{1}} \tag{4}
\end{equation*}
$$

and observe that owing to the inequalities $1<\gamma$ and $p_{1}>1$, we have $1<r<$ $\gamma$. Next, starting with the formula

$$
|B F|^{r}=\left(|B \| F|^{1 / \gamma}\right)^{r}|F|^{r-r / \gamma}
$$

and using Hölder's inequality (with $s=\frac{\gamma}{r}(>1)$ and $t:=\frac{s}{s-1}=\frac{\gamma}{\gamma-r}$ ) and the assumptions $|B \| F|^{1 / \gamma} \in L_{\text {loc }}^{\gamma}(\Omega, d x)$ and $F \in L_{\text {loc }}^{p_{1}}(\Omega, d x)$, we get that $B^{i} F \in L_{\mathrm{loc}}^{r}(\Omega, d x)$. By Lemma 3 (i)

$$
\begin{equation*}
B^{i} F \in H_{\mathrm{loc}}^{r, 0}(\Omega), \quad\left(B^{i} F\right)_{x^{i}} \in H_{\mathrm{loc}}^{r,-1}(\Omega) \tag{5}
\end{equation*}
$$

(ii): Set

$$
\begin{equation*}
q:=q\left(p_{1}\right):=\frac{\gamma p_{1}}{\gamma-2+2 p_{1}} \vee 1 \tag{6}
\end{equation*}
$$

and note that $q>1 \Leftrightarrow \gamma>2 \Leftrightarrow q<\frac{\gamma}{2}$, in particular, $q<\gamma$ in any case. Hence repeating the above argument with $c, \gamma / 2, q$ replacing $|B|, \gamma, r$, respectively we obtain that

$$
\begin{equation*}
c F \in L_{\mathrm{loc}}^{q}(\Omega, d x) \tag{7}
\end{equation*}
$$

Fix $p_{1}>1$ such that $F:=\frac{d \mu}{d x} \in L_{\mathrm{loc}}^{p_{1}}(\Omega, d x)$ and let $r, q$ be as in (4), (6), correspondingly. Since $\gamma \leq n$ we have that $q<n$, which by (7) and Lemma 3 (ii) resp. (iii) yields $c F \in H_{\mathrm{loc}}^{n q /(n-q),-1}(\Omega)$ if $q>1$ resp. $c F \in H_{\mathrm{loc}}^{s,-1}(\Omega)$ for any $s \in(1, n /(n-1))$ if $q=1$.

It turns out that if $p_{1}<n /(n-\gamma)$, then

$$
\begin{equation*}
c F \in H_{\mathrm{loc}}^{r,-1}(\Omega) \tag{8}
\end{equation*}
$$

Indeed, if $q>1$, then (8) follows from the fact that if $p_{1} \in(1, n /(n-\gamma))$ the inequality $r \leq n q /(n-q)$ holds. If $q=1$, then $\gamma \leq 2$ and (8) follows from the fact that $r<n /(n-\gamma+1) \leq n /(n-1)$ for $p_{1}<n /(n-\gamma)$.

Finally by Lemma 3 (ii) we have $v \in H_{\mathrm{loc}}^{n /(n-\gamma+1),-1}(\Omega)$ if $\gamma>2$ and $v \in H_{\mathrm{loc}}^{s,-1}(\Omega)$ for any $s \in(1, n /(n-1))$ if $\gamma \leq 2$. In the same way as above, $\nu \in H_{\mathrm{loc}}^{r,-1}(\Omega)$ whenever $1<p_{1}<n /(n-\gamma)$. This along with (5) and (8) shows that the right-hand side of (3) is now in $H_{\mathrm{loc}}^{r,-1}(\Omega)$. By Lemma 3 (i) we have

$$
\begin{equation*}
\mu \in H_{\mathrm{loc}}^{r, 1}(\Omega) \tag{9}
\end{equation*}
$$

and by Lemma 3 (ii) $F \in L_{\text {loc }}^{p_{2}}(\Omega, d x)$, where

$$
p_{2}:=\frac{n r}{n-r}=\frac{n \gamma p_{1}}{n \gamma-n+(n-\gamma) p_{1}}=: f\left(p_{1}\right)
$$

Thus we get

$$
p_{1} \in\left(1, \frac{n}{n-\gamma}\right) \text { and } F \in L_{\mathrm{loc}}^{p_{1}}(\Omega, d x) \Longrightarrow F \in L_{\mathrm{loc}}^{f\left(p_{1}\right)}(\Omega, d x) .
$$

One can easily check that $p_{2}=f\left(p_{1}\right)>p_{1}$ if $p_{1}<n /(n-\gamma)$, and that the only positive solution of the equation $q=f(q)$ is $q=n /(n-\gamma)$. Therefore, by taking $p_{1}$ from ( $1, n /(n-1)$ ), which is possible by (i), and by defining $p_{k+1}=f\left(p_{k}\right)$ we get an increasing sequence of $p_{k} \uparrow n /(n-\gamma)$, which implies that $F \in L_{\mathrm{loc}}^{p}(\Omega, d x)$ for any $p<n /(n-\gamma)$.

But as $p_{k} \nearrow n /(n-\gamma), r\left(p_{k}\right)$ (defined according to (4)) increasingly converges to

$$
\frac{\gamma n /(n-\gamma)}{\gamma-1+n /(n-\gamma)}=\frac{n}{n-\gamma+1} .
$$

By (9) this proves (ii).
(iii): First we consider case (b) in which $|B| \in L_{\mathrm{loc}}^{\gamma}(\Omega, \mu), c \in L_{\mathrm{loc}}^{n \gamma /(n+\gamma)}$ $(\Omega, \mu), \nu \in L_{\text {loc }}^{n \gamma /(n+\gamma)}(\Omega, d x)$. By the last assertion in (ii) we have $F \in$ $L_{\text {loc }}^{p_{1}}(\Omega, d x)$ for any (finite) $p_{1}>1$. Let $r:=r\left(p_{1}\right)$ be defined as in (4). Then $1<r<\gamma$ and (5) holds. Set

$$
\begin{equation*}
q:=q\left(p_{1}\right):=\frac{\frac{n \gamma}{n+\gamma} p_{1}}{\frac{n \gamma}{n+\gamma}-1+p_{1}} \tag{10}
\end{equation*}
$$

$2 \leq n<\gamma$, implies $\frac{n \gamma}{n+\gamma}>1$. Therefore, (since $p_{1}>1$ ) it follows that $1<q<$ $\frac{n \gamma}{n+\gamma}$. Hence repeating the arguments that led to (5) with $c, \frac{n \gamma}{n+\gamma}, q$ replacing $|B|$, $\gamma, r$ respectively we obtain $c F \in L_{\mathrm{loc}}^{q}(\Omega, d x)$, thus $c F \in H^{n q /(n-q),-1}(\Omega)$ by Lemma 3 (ii). Observe that when $p_{1} \rightarrow \infty$, we have $r \uparrow \gamma, q \uparrow n \gamma /(n+\gamma)$, and $n q /(n-q) \uparrow \gamma$. Therefore, combining this with our assumption that $v \in L_{\text {loc }}^{n \gamma /(n+\gamma)}(\Omega, d x)$ which by Lemma 3 (ii) is contained in $H_{\mathrm{loc}}^{\gamma,-1}(\Omega)$, by taking $p_{1}$ large enough, we see that the right-hand side in (3) is in $H_{\mathrm{loc}}^{\gamma-\varepsilon,-1}(\Omega)$ for any $\varepsilon \in(0, \gamma-1)$. By Lemma 3 (ii) we conclude $F \in H_{\mathrm{loc}}^{\gamma-\varepsilon, 1}(\Omega)$ and since $\gamma>n$, the function $F$ is locally bounded. Now we see that above we can take $p_{1}=\infty$ and therefore the right-hand side of (3) is in $H_{\mathrm{loc}}^{\gamma,-1}(\Omega)$, which by Lemma 3 (i) gives us the desired result.

In the remaining case (a) we take $p_{1}>\gamma /(\gamma-1)$ and assume that $F \in$ $L_{\mathrm{loc}}^{p_{1}}(\Omega, d x)$. Then instead of (4) and (10) we define

$$
\begin{equation*}
r:=r\left(p_{1}\right):=\frac{\gamma p_{1}}{\gamma+p_{1}}, \quad q:=q\left(p_{1}\right):=\frac{\frac{n \gamma}{n+\gamma} p_{1}}{\frac{n \gamma}{n+\gamma}+p_{1}} \vee 1 \tag{11}
\end{equation*}
$$

and observe that owing to $p_{1}>\gamma /(\gamma-1)$ we have $r>1$, which (because $p_{1}^{-1}+\gamma^{-1}=r^{-1}$ ) allows us to apply Hölder's inequality starting with $|B F|^{r}=$ $|B|^{r}|F|^{r}$ to conclude that (5) holds. Since $c \in L_{\text {loc }}^{1}(\Omega, \mu)$, resp. $\frac{n \gamma}{n+\gamma}>1$ and $\left(\frac{n \gamma}{n+\gamma}\right)^{-1}+p_{1}^{-1}=q^{-1}$, we also have that $c F \in L_{\mathrm{loc}}^{q}(\Omega, d x)$. Obviously, $q<n$. As in part (ii) this yields that $c F \in H_{\mathrm{loc}}^{n q /(n-q),-1}(\Omega)$ if $q>1$ and $c F \in H_{\text {loc }}^{s,-1}(\Omega)$ for any $s \in(1, n /(n-1))$ if $q=1$. We claim that (8) holds (with $r=r\left(p_{1}\right)$ as in (11) for all $p_{1}>\gamma /(\gamma-1), p_{1} \neq n \gamma /(n \gamma-n-\gamma)$.

Indeed, if $q>1$, then $n q /(n-q)=r$. If $q=1$, then $p_{1} \leq n \gamma /(n \gamma-n-\gamma)$. But since $p_{1} \neq n \gamma /(n \gamma-n-\gamma)$, we have $p_{1}<n \gamma /(n \gamma-n-\gamma)$, which is equivalent to the inequality $r<n /(n-1)$.

Thus, since $v \in L_{\mathrm{loc}}^{n \gamma /(n+\gamma)}(\Omega, d x) \subset H_{\mathrm{loc}}^{\gamma,-1}(\Omega) \subset H_{\mathrm{loc}}^{r,-1}(\Omega)$ (because $r<$ $\gamma$ ), it follows by Lemma 2 (i) that:

$$
\begin{equation*}
\binom{p_{1}>\frac{\gamma}{\gamma-1} \text { and } p_{1} \neq \frac{n \gamma}{n \gamma-n-\gamma}}{\text { and } F \in L_{\mathrm{loc}}^{p_{1}}(\Omega, d x)} \Longrightarrow F \in H_{\mathrm{loc}}^{r, 1}(\Omega) \tag{12}
\end{equation*}
$$

Provided $r<n$ the latter in turn by Lemma 3 (ii) implies that $F \in L_{\mathrm{loc}}^{p_{2}}(\Omega, d x)$. Summarizing we have thus shown:

$$
\left(\begin{array}{ll}
p_{1}>\frac{\gamma}{\gamma-1} & \text { and } p_{1} \neq \frac{n \gamma}{n \gamma-n-\gamma}  \tag{13}\\
\text { and } r:=\frac{\gamma p_{1}}{\gamma+p_{1}}<n & \text { and } F \in L_{\mathrm{loc}}^{p_{1}}(\Omega, d x)
\end{array}\right) \Longrightarrow F \in L_{\mathrm{loc}}^{p_{2}}(\Omega, d x)
$$

where

$$
p_{2}:=\frac{n r}{n-r}=\frac{n \gamma p_{1}}{n \gamma-(\gamma-n) p_{1}}>\frac{n \gamma}{n \gamma-(\gamma-n)} p_{1} .
$$

Also notice that $\gamma /(\gamma-1)<n /(n-1)<\frac{n \gamma}{\gamma n-n-\gamma}$ so that by (i) we can take a $p_{1}$ to start with. Then starting with $p_{1}$ close enough to $n /(n-1)$, by iterating (13) we always increase $p$ by a certain factor $>1$. While doing so we can obviously choose the first $p$ so that the iterated $p$ 's will be never equal to $n \gamma /(n \gamma-n-\gamma)$ and the corresponding $r$ 's will not coincide with $n$. Then after several steps we shall come to the situation where $r>n$, and then we conclude from (12) that $F$ is locally bounded (one cannot keep iterating (13) infinitely having the restriction $r<n$ ). As in case (b) one can now easily complete the proof.

Remark 4 (i) For sufficiently regular $F$ with no zeros operators of the type considered above become special cases of operators $L=\sum_{i, j} \partial_{i}\left(a_{i j} \partial_{j}\right)+q$. Additional information (including further references) about the essential selfadjointness of such operators, however, considered on $L^{2}\left(R^{n}, d x\right)$ can be found in [8], [15].
(ii) In a forthcoming paper the parabolic case will be studied. It is, however, immediate from Theorem 1 that if $t \mapsto \mu_{t}$ is differentiable such that $\frac{\partial}{\partial t} \mu_{t}$ is a

Radon measure, then for fixed $t$ the densities $F_{t}$ of $\mu_{t}$ w.r.t. $d x$ exist and all respective assertions in Theorem 1 hold for $F_{t}$.
(iii) Note that the only property of the operator $L_{0}:=\Delta$ used above was the one mentioned in Lemma 3 (i), i.e., that $u \in H_{\mathrm{loc}}^{p, r+2}(\Omega)$ provided $L_{0} u \in H_{\text {loc }}^{p, r}(\Omega)$. It is known (see, e.g., [21, p. 270]) that this holds for arbitrary non-degenerate second order elliptic operators with smooth coefficients. Therefore, Theorem 1 remains valid if we replace $\Delta$ by any non-degenerate second order elliptic operator $L_{0}$ with smooth coefficients. Moreover, as a thorough inspection of the proof of Theorem 4.2.4 in [22] shows, one can relax the assumption about the smoothness of the coefficients of $L_{0}$ here even more. Note, in particular, that Theorem 1 extends to elliptic second order operators on smooth Riemannian manifolds with non-degenerate smooth second order parts.
(iv) It should be noted that the elliptic equations discussed here cannot be reduced to those considered e.g. in [10], [13], [18], [24]. There are two major differences. The first is that the solutions considered there by definition are supposed to be in $H_{\mathrm{loc}}^{\gamma, 1}\left(\mathbb{R}^{n}\right)$. Secondly, our integrability conditions for $B$ are w.r.t. a measure $\mu$ which is a solution of our equation. For this reason, $B$ need not be locally Lebesgue integrable; e.g. if $\mu$ is given by the density $x^{2} \exp \left(-x^{2}\right)$ on $\mathbb{R}^{1}$, then it solves our elliptic equation with $B(x)=\beta(x)=-2 x+2 / x$. Of course, Theorem 1 (iii) shows that under sufficient integrability conditions our solutions become solutions also in the sense of the above mentioned references. However, in general we get a wider class of solutions. Note also that in our setting due to the weak assumptions on $B$ the elliptic regularity does not imply that solutions belong to the second Sobolev class $H_{\mathrm{loc}}^{\gamma, 2}$ (e.g. any $\mu=\rho d x$ with $\rho \in H_{\mathrm{loc}}^{1,1}$ satisfies (1) with $B:=\nabla \rho / \rho, c:=0, v:=0$ ).

The next example shows that assertion (iii) of Theorem 1 fails if $n+\varepsilon$ is replaced by $n-\varepsilon$. (Then $F$ does not even need to be in $H_{\text {loc }}^{2,1}(\Omega)$.)

Example 5. Let $n>3$ and

$$
L^{*} F(x)=\Delta F(x)+\alpha\left(x^{i}|x|^{-2} F\right)_{x^{i}}(x)-F(x)
$$

where $\alpha=n-3$. Then the function $F(x)=\left(e^{r}-e^{-r}\right) r^{-(n-2)}, r=|x|$, is locally $d x$-integrable and $L^{*} F=0$ in the sense of distributions, but $F$ is not in $H_{\text {loc }}^{2,1}\left(\mathbb{R}^{n}\right)$. Here $B(x)=-\alpha x\|x\|^{-2}=\nabla\left(|x|^{-\alpha}\right) /|x|^{-\alpha}$ and $|B| \in L_{\text {loc }}^{n-\varepsilon}\left(\mathbb{R}^{n}, d x\right)$ for all $\varepsilon>0$. In a similar way, if there is no " $-F$ " in the equation above, then the function $F(x)=r^{-(n-3)}$ has the same properties.

Proof. Observe that $F_{x^{i}}, F_{x^{i} x^{j}}$ are locally $d x$-integrable. Therefore, the equation $L^{*} F=0$ follows easily from the equation on $(0, \infty)$

$$
f^{\prime \prime}+\frac{(n-1+\alpha)}{r} f^{\prime}+\alpha \frac{n-2}{r^{2}} f-f=0
$$

which is satisfied for the function $f(r)=\left(e^{r}-e^{-r}\right) r^{-(n-2)}$. It remains to note that $F, \nabla F$ and $\Delta F$ are locally $d x$-integrable, since $f(r) r^{n-1}, f^{\prime}(r) r^{n-1}$,
$f^{\prime \prime}(r) r^{n-1}$ are locally bounded, but $\nabla F$ is not $d x$-square-integrable at the origin. (If $n \geq 6$, then also $F$ is not $d x$-square-integrable at the origin). In the case without " $-F$ " in the equation similar (but even simpler) arguments can be used to show that $F(x)=r^{-(n-3)}$ has the same properties.

Remark 6. Applying the regularity result in Theorem 1 (ii) above to the case $c=0=v$ we get, in particular, the existence of a density in $H_{\mathrm{loc}}^{p, 1}\left(\mathbb{R}^{n}\right)$, for $p \in\left[1, \frac{n}{n-\varepsilon}\right.$ ), for any invariant measure $\mu$ of a diffusion $\xi_{t}$ driven by the stochastic differential equation $d \xi_{t}=d w_{t}+B\left(\xi_{t}\right) d t$, where the drift $B$ is assumed to be in $L_{\mathrm{loc}}^{1+\varepsilon}\left(\mathbb{R}^{n}, \mu\right)$. This is true for any interpretation of a solution which implies (1) for invariant measures. Thus, we get an improvement of a part of a theorem in [3], [4] (see also [2] for the case of a non-constant second order part). In [3], [4] under the a priori assumption that $\mu$ is a probability measure and assuming that $|B|$ is globally in $L^{2}\left(\mathbb{R}^{n}, \mu\right)$, it was shown that $\mu$ admits a density in $H^{1,1}\left(\mathbb{R}^{n}\right)$. (We would like to mention that under these stronger conditions the latter result can also be deduced from [6]).

We say that a measurable function $f$ on $\mathbb{R}^{n}$ is locally uniformly positive if $\operatorname{essinf}_{U} f>0$ for every ball $U \subset \mathbb{R}^{n}$.

Theorem 7. Let $n \geq 2$ and let $\mu$ be a measure on $\mathbb{R}^{n}$ with density $\rho:=\varphi^{2}$, $\varphi \in H_{\mathrm{loc}}^{2,1}\left(\mathbb{R}^{n}\right)$, which is locally uniformly positive. Assume that $|\beta| \in L_{\mathrm{loc}}^{\gamma}\left(\mathbb{R}^{n}, \mu\right)$, where $\beta:=\nabla \rho / \rho$ and $\gamma>n$. Then the operator

$$
L \psi=\Delta \psi+\langle\nabla \psi, \beta\rangle
$$

with domain $C_{0}^{\infty}\left(\mathbb{R}^{n}\right)$ is essentially selfadjoint on $L^{2}\left(\mathbb{R}^{n}, \mu\right)$.
Proof. First we note that since $\mu$ satisfies (1) with $B:=\beta, c: \equiv 0, v: \equiv 0$, it follows by Theorem 1 (iii), part (b), that $\rho$ is continuous, hence locally bounded. Assume that there is a function $g \in L^{2}\left(\mathbb{R}^{n}, \mu\right)$ such that

$$
\begin{equation*}
\int(L-1) \zeta(x) g(x) \mu(d x)=0 \quad \forall \zeta \in C_{0}^{\infty}\left(\mathbb{R}^{n}\right) \tag{14}
\end{equation*}
$$

Recall that by definition $\beta=0$ on the set $\{\rho=0\}$ (which is reasonable since $\nabla \rho=0 d x$-a.e. on $\{\rho=0\}$ ). Clearly, $|\beta| \in L_{\text {loc }}^{\gamma}\left(\mathbb{R}^{n}, d x\right)$. Consequently, by Theorem 1 (iii), Part (a), $F \in H_{\mathrm{loc}}^{\gamma, 1}\left(\mathbb{R}^{n}\right)$. In particular, $F$ is continuous and locally bounded. Then $g=F / \rho \in H_{\mathrm{loc}}^{\gamma, 1}\left(\mathbb{R}^{n}\right) \bigcap L_{\mathrm{loc}}^{\infty}\left(\mathbb{R}^{n}\right), g|\beta| \in L_{\mathrm{loc}}^{\gamma}\left(\mathbb{R}^{n}, d x\right)$. Therefore, we can integrate by parts in equality (14) which yields for every $\zeta \in C_{0}^{\infty}\left(\mathbb{R}^{n}\right)$

$$
\begin{align*}
0 & =-\int\langle\nabla \zeta, \nabla g\rangle d \mu-\int\langle\nabla \zeta, \beta\rangle g d \mu+\int\langle\nabla \zeta, \beta\rangle g d \mu-\int \zeta g d \mu  \tag{15}\\
& =-\int\langle\nabla \zeta, \nabla g\rangle d \mu-\int \zeta g d \mu
\end{align*}
$$

Now let $\psi \in C_{0}^{\infty}\left(\mathbb{R}^{n}\right)$ and $\varphi \in H_{\mathrm{loc}}^{2,1}\left(\mathbb{R}^{n}\right)$. Then by the product rule

$$
\begin{equation*}
\langle\nabla \varphi, \nabla(\psi g)\rangle=\langle\nabla(\psi \varphi), \nabla g\rangle-\varphi\langle\nabla \psi, \nabla g\rangle+g\langle\nabla \varphi, \nabla \psi\rangle \tag{16}
\end{equation*}
$$

Since equality (15) extends to all $\zeta$ in $H^{2,1}\left(\mathbb{R}^{n}\right)$ with compact support, we can apply (15) to $\zeta:=\psi \varphi$ and use (16) to obtain

$$
\begin{aligned}
& \int\langle\nabla \varphi, \nabla(\psi g)\rangle d \mu+\int \varphi \psi g d \mu \\
& \overline{(16)} \int\langle\nabla(\psi \varphi), \nabla g\rangle d \mu-\int \varphi\langle\nabla \psi, \nabla g\rangle d \mu \\
&+\int g\langle\nabla \varphi, \nabla \psi\rangle d \mu+\int \varphi \psi g d \mu \\
&\left(\overline{\overline{(15)}}-\int \varphi\langle\nabla \psi, \nabla g\rangle+\right. \int g\langle\nabla \varphi, \nabla \psi\rangle d \mu
\end{aligned}
$$

Taking $\varphi:=\psi g$, one gets

$$
\begin{aligned}
\int\langle\nabla(\psi g), & \nabla(\psi g)\rangle d \mu+\int(\psi g)^{2} d \mu \\
& =-\int \psi g\langle\nabla \psi, \nabla g\rangle d \mu+\int g\langle\nabla(\psi g), \nabla \psi\rangle d \mu \\
& =\int g^{2}\langle\nabla \psi, \nabla \psi\rangle d \mu
\end{aligned}
$$

Hence, we get

$$
\begin{equation*}
\int(\psi g)^{2} d \mu \leq \int g^{2}|\nabla \psi|^{2} d \mu \tag{17}
\end{equation*}
$$

Taking a sequence $\psi_{k} \in C_{0}^{\infty}\left(\mathbb{R}^{n}\right), k \in \mathbb{N}$, such that $0 \leq \psi_{k} \leq 1, \psi_{k}(x)=1$ if $|x| \leq k, \psi_{k}(x)=0$ if $|x| \geq k+1$, and $\sup _{k}\left|\nabla \psi_{k}\right|=M<\infty$, we get by Lebesgue's dominated convergence theorem that the left hand side of (17) tends to $\|g\|_{2}^{2}$, while the right hand side tends to zero. Thus, $g=0$. By a standard result (see, e.g., [12]) this implies the essential self-adjointness of ( $L, C_{0}^{\infty}\left(\mathbb{R}^{n}\right)$ ) on $L^{2}\left(\mathbb{R}^{n}, \mu\right)$.

Corollary 8. The assertion of the previous theorem holds true if $\mu$ is a measure on $\mathbb{R}^{n}$ with density $\rho:=\varphi^{2}, \varphi \in H_{\mathrm{loc}}^{2,1}\left(\mathbb{R}^{n}\right)$, and $|\beta| \in L_{\mathrm{loc}}^{\gamma}\left(\mathbb{R}^{n}, d x\right)$, where $\beta:=\nabla \rho / \rho$ and $\gamma>n$.

Proof. Note that $\rho$ admits a continuous strictly positive modification. Indeed, if $f_{n}:=\log \left(\rho+\frac{1}{n}\right), n \in \mathbb{N}$, then $f_{n} \xrightarrow[n \rightarrow \infty]{\longrightarrow} \log \rho$ in $L_{\text {loc }}^{1}\left(\mathbb{R}^{n}, d x\right)$, which easily follows from the fact that $\log \rho \in L_{\mathrm{loc}}^{1}\left(\mathbb{R}^{n}, d x\right)$. The latter in turn follows from [3, Lemma 6.4]. Consequently by the Poincaré inequality, the sequence $\left(f_{n}\right)_{n \in \mathbb{N}}$ is bounded in $H^{\gamma, 1}(U)$ for every open ball $U \subset \mathbb{R}^{n}$. By the compactness of the embedding $H^{\gamma, 1}(U) \rightarrow C(U)$, a subsequence of the sequence of the continuous modifications of $\left(f_{n}\right)_{n \in \mathbb{N}}$ converges locally uniformly to $\log \rho$. Whence $\rho$ is continuous and strictly positive. In particular, $|\nabla \rho / \rho| \in L_{\mathrm{loc}}^{\gamma}\left(\mathbb{R}^{n}, \mu\right)$.

Remark 9. If $\mu=\rho d x$ with $\rho=\varphi^{2}$ and $\varphi \in H_{\mathrm{loc}}^{2,1}\left(\mathbb{R}^{n}\right)$, the so-called Markov uniqueness (i.e, the uniqueness of a Markovian semigroup on $L^{2}\left(\mathbb{R}^{n}, \mu\right)$ with generator given by $L f=\Delta f+\langle\nabla f, \beta\rangle$ on $C_{0}^{\infty}\left(\mathbb{R}^{n}\right)$ ) always holds with $\beta:=\nabla \rho / \rho$ (see [16], [17]). However, in general Markov uniqueness is weaker than the essential self-adjointness of $\left(L, C_{0}^{\infty}\left(\mathbb{R}^{n}\right)\right)$ on $L^{2}\left(\mathbb{R}^{n}, \mu\right)$. (see [7]). Optimal (local or global) conditions for the essential self-adjointness remain unknown except for the one-dimensional case investigated in [25] and [7]. In fact, recently in [7] a complete characterization of the essential self-adjointness for Dirichlet operators has been given in the case $n=1$.

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