

# Potentials in the Kato Class of Measures And Other Very Singular Potentials

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## 1. Introduction

Let  $M$  be a compact Riemannian manifold, of dimension  $n \geq 2$ , with volume element  $d\omega$  and Laplace operator  $\Delta$ . Set

$$(1.1) \quad G(x, y) = \gamma_n(d(x, y)),$$

where  $d(x, y)$  is the Riemannian distance from  $x$  to  $y$ , and

$$(1.2) \quad \begin{aligned} \gamma_n(r) &= r^{2-n}, & n \geq 3, \\ &\log \frac{A}{r}, & n = 2, \end{aligned}$$

where  $A > \text{diam } M$ . An essentially equivalent specification of  $G$  is

$$(1.3) \quad (1 - \Delta)^{-1}u(x) = \int_M G(x, y)u(y) d\omega(y).$$

The Kato class  $\mathcal{K}(M)$  of integrable functions, introduced in [K], is defined as follows.

**Definition.** Given  $V \in L^1(M)$ , we say  $V \in \mathcal{K}(M)$  provided

$$(1.4) \quad \lim_{r \rightarrow 0} \sup_{x \in M} \int_{B_r(x)} G(x, y)|V(y)| d\omega(y) = 0.$$

Equivalent conditions are

$$(1.5) \quad \lim_{\lambda \rightarrow +\infty} \sup_{x \in M} (\lambda - \Delta)^{-1}|V|(x) = 0,$$

hence

$$(1.6) \quad \lim_{\lambda \rightarrow +\infty} \sup_{\|u\|_{L^1} \leq 1} |(u, (\lambda - \Delta)^{-1}|V||) = 0,$$

hence

$$(1.7) \quad \lim_{\lambda \rightarrow +\infty} \sup_{\|u\|_{L^1} \leq 1} |(|V|(\lambda - \Delta)^{-1}u, 1)| = 0,$$

hence (since  $(\lambda - \Delta)^{-1}$  is positivity-preserving)

$$(1.8) \quad \lim_{\lambda \rightarrow +\infty} \|V(\lambda - \Delta)^{-1}\|_{\mathcal{L}(L^1(M))} = 0.$$

Here is another characterization (a consequence of Theorem 4.15 of [AS]).

**Proposition 1.1.** *If  $V \in L^1(M)$ , then*

$$(1.9) \quad V \in \mathcal{K}(M) \iff (1 - \Delta)^{-1}|V| \in C(M).$$

There is also the following property of  $M_V$  (defined by  $M_V u = Vu$ ), important for defining  $-\Delta + V$  as a self-adjoint operator.

**Proposition 1.2.** *If  $V \in \mathcal{K}(M)$ , then*

$$(1.10) \quad M_V : H^1(M) \longrightarrow H^{-1}(M),$$

and for each  $\varepsilon > 0$ , there exists  $C(\varepsilon) < \infty$  such that, for  $u \in H^1(M)$ ,

$$(1.11) \quad \int_M |V|u^2 d\omega \leq \varepsilon \|u\|_{H^1}^2 + C(\varepsilon) \|u\|_{L^2}^2.$$

For a generalization involving fractional powers of  $1 - \Delta$ , see Theorem 4.2 of [ZY].

The notion of Kato class for integrable functions given above has the following natural extension (introduced in [BM]).

**Definition.** Let  $\mu$  be a finite, signed Borel measure on  $M$ . We say  $\mu \in \mathcal{K}(M)$  provided

$$(1.12) \quad \lim_{r \rightarrow 0} \sup_{x \in M} \int_{B_r(x)} G(x, y) d|\mu|(y) = 0,$$

where  $|\mu|$  denotes the positive “total variation” measure associated with  $\mu$ , via the Hahn decomposition.

Lots of singular measures satisfy (1.12):

**Proposition 1.3.** *Assume that there exists  $A < \infty$  such that, for all  $x \in M$ ,  $r > 0$ ,*

$$(1.13) \quad |\mu|(B_r(x)) \leq Ar^\alpha, \quad \alpha > n - 2.$$

*Then (1.12) holds.*

*Proof.* We have

$$(1.14) \quad \int_{B_r(x) \setminus B_{r/2}(x)} G(x, y) d|\mu|(y) \leq C\gamma_n(r)r^\alpha.$$

For  $n \geq 3$ , this is dominated by

$$(1.15) \quad Cr^{\alpha-(n-2)},$$

and for  $n = 2$  it is dominated by

$$(1.16) \quad Cr^\alpha \log \frac{A}{r}.$$

In either case, we have (1.12).

REMARK. In case  $n = 2$ , the condition (1.13) is

$$(1.17) \quad |\mu|(B_r(x)) \leq Ar^\alpha, \quad \text{for some } \alpha > 0.$$

We see in Appendix A (which arose in the course of writing [CST]) that (1.17) holds for various Cantor sets in a 2D manifold  $M$ , endowed with Hausdorff measure. It is also seen in Chapter 7 of [T2] that (1.17) holds when  $\mu$  is the maximal entropy measure on a Julia set in  $S^2$ , associated with a holomorphic map on  $S^2$  of degree  $\geq 2$ . In Appendix A we show directly that

$$(1.18) \quad M_\mu : H^1(M) \longrightarrow H^{-1}(M) \text{ is compact,}$$

whenever  $\mu$  is a finite positive measure on an  $n$ -dimensional Riemannian manifold  $M$  satisfying

$$(1.19) \quad \mu(B_r(x)) \leq Cr^\alpha, \quad \alpha > n - 1 - \frac{1}{n-1}.$$

Returning to the description of our analysis of Kato-class measures, in §2 we show that Propositions 1.1 and 1.2 extend from  $L^1$  functions  $V \in \mathcal{K}(M)$  to finite, signed measures  $\mu \in \mathcal{K}(M)$ . We deduce that, for such  $\mu$ , there exists a constant  $C$  such that

$$(1.20) \quad A = -\Delta + M_\mu + C \text{ is a positive self-adjoint operator on } L^2(M),$$

satisfying

$$(1.21) \quad \mathcal{D}(A^{1/2}) = H^1(M).$$

Going further, in §3, we examine self-adjoint extensions of such operators as

$$(1.22) \quad A = -\Delta + M_\mu + FM_\nu + C,$$

with  $\mu$  and  $C$  as in (1.20),  $\nu$  a positive measure on  $M$  satisfying

$$(1.23) \quad M_\nu : H^1(M) \longrightarrow H^{-1}(M),$$

and

$$(1.24) \quad F \geq 1, \quad F \in L^1(M, \nu).$$

In case  $\mu = 0$  and  $\nu = \omega$ , Proposition 3.2 specializes to a classical result of Friedrichs (cf. [CFKS]), Chapter 1.

In §§4–5 we point to further results on Schrödinger operators with singular potentials. We cite various papers that we have not digested.

Appendix A gives a general discussion of multipliers from  $H^1(M)$  to  $H^{-1}(M)$ , taken from [CST].

## 2. Properties of measures in $\mathcal{K}(M)$

Our first goal here is to extend Proposition 1.1. In preparation for this, let us take  $G(x, y)$  as in (1.3) and write (for  $t > 0$ )

$$(2.1) \quad (1 - \Delta)^{-1} e^{t(\Delta-1)} |\mu|(x) = \int_M G_t(x, y) d|\mu|(y),$$

and

$$(2.2) \quad R_t(x, y) = G(x, y) - G_t(x, y),$$

so

$$(2.3) \quad \begin{aligned} \int_M R_t(x, y) d|\mu|(y) &= (1 - \Delta)^{-1} (1 - e^{t(\Delta-1)}) |\mu|(x) \\ &= \int_0^t e^{s(\Delta-1)} |\mu|(x) ds, \end{aligned}$$

the last identity yielding  $R_t(x, y) \geq 0$ . Comparison with (1.12) shows that

$$(2.4) \quad \begin{aligned} \mu \in \mathcal{K}(M) &\iff \lim_{t \searrow 0} \sup_x \int_M R_t(x, y) d|\mu|(y) = 0 \\ &\iff \lim_{t \searrow 0} \|(1 - e^{t(\Delta-1)})(1 - \Delta)^{-1} |\mu|\|_{\text{sup}} = 0. \end{aligned}$$

As an aside, we mention that, if  $\mu$  is a finite signed measure on  $M$ , then, for each  $\varepsilon > 0$ , there exists  $\delta > 0$  such that

$$(2.5) \quad \psi := (1 - \Delta)^{-1} |\mu| \in H^{2-\varepsilon, 1+\delta}(M), \quad \psi \geq 0.$$

Here is our extension of Proposition 1.1.

**Proposition 2.1.** *Given that  $\mu$  is a finite signed measure on  $M$ ,*

$$(2.6) \quad \mu \in \mathcal{K}(M) \iff (1 - \Delta)^{-1} |\mu| \in C(M).$$

*Proof.* Take  $\psi$  as in (2.5). Since  $\{e^{t(\Delta-1)} : t \geq 0\}$  is a strongly continuous semigroup on  $C(M)$  and, for each  $\psi \in \mathcal{D}'(M)$ ,  $e^{t(\Delta-1)} \psi \in C^\infty(M)$  for all  $t > 0$ , we have

$$(2.7) \quad \begin{aligned} \psi \in C(M) &\iff \lim_{t \searrow 0} \|(1 - e^{t(\Delta-1)}) \psi\|_{\text{sup}} = 0 \\ &\iff \mu \in \mathcal{K}(M), \end{aligned}$$

the last equivalence by (2.4).

Our next goal is to extend Proposition 1.2, and show that if  $\mu$  is a finite signed measure,

$$(2.8) \quad \mu \in \mathcal{K}(M) \implies M_\mu : H^1(M) \rightarrow H^{-1}(M),$$

accompanied by estimates of the form (1.11). Our analysis will be adapted from the proof of Theorem 4.2 in [ZY], specialized from  $(-\Delta)^{-\alpha/2}$  to the classical case  $\alpha = 2$ . To start, note that if  $\mu = \mu^+ - \mu^-$  is the Hahn decomposition,  $\mu \in \mathcal{K}(M)$  implies  $\mu^\pm \in \mathcal{K}(M)$ , so we can restrict attention to the case

$$(2.9) \quad \mu \in \mathcal{K}(M) \text{ is a positive measure.}$$

It will be convenient to pass to

$$(2.10) \quad \nu = \mu + \omega,$$

where  $\omega$  is the volume measure on  $M$ . Clearly also  $\nu \in \mathcal{K}(M)$ . Parallel to (1.5), we have, for  $\lambda > 0$ ,

$$(2.11) \quad \sup_{x \in M} (\lambda - \Delta)^{-1} \nu(x) \leq \varepsilon(\lambda), \quad \varepsilon(\lambda) \rightarrow 0 \text{ as } \lambda \rightarrow +\infty.$$

In order to adapt arguments from [ZY], we set

$$(2.12) \quad V_s = e^{s\Delta} \nu, \quad s > 0,$$

so  $V_s \in C^\infty(M)$ , and

$$(2.13) \quad V_s = W_s + 1, \quad W_s = e^{s\Delta} \mu \geq 0.$$

Here, as is natural, we implement the identification  $C^\infty(M) \hookrightarrow \mathcal{D}'(M)$  by  $f \mapsto f\omega$ . Note that

$$(2.14) \quad \sup_{x \in M} (\lambda - \Delta)^{-1} V_s(x) \leq \varepsilon(\lambda),$$

with  $\varepsilon(\lambda)$  exactly as in (2.11), since  $\{e^{s\Delta} : s \geq 0\}$  is a contraction semigroup on  $C(M)$ . In particular, the estimate (2.14) is independent of  $s \in (0, \infty)$ . Hence, parallel to (1.8), we have

$$(2.15) \quad \|V_s(\lambda - \Delta)^{-1}\|_{\mathcal{L}(L^1(M))} \leq \varepsilon(\lambda).$$

By duality,

$$(2.16) \quad \|(\lambda - \Delta)^{-1} V_s\|_{\mathcal{L}(L^\infty(M))} \leq \varepsilon(\lambda).$$

Now, following [ZY], we bring in weighted  $L^p$ -spaces,

$$(2.17) \quad L_s^p(M) = L^p(M, V_s d\omega).$$

We see that  $M_{V_s^{1/p}} : L_s^p(M) \rightarrow L^p(M)$  is an isometric isomorphism. Hence

$$(2.18) \quad \begin{aligned} \|(\lambda - \Delta)^{-1} V_s\|_{\mathcal{L}(L_s^1(M))} &= \|V_s(\lambda - \Delta)^{-1}\|_{\mathcal{L}(L^1(M))} \leq \varepsilon(\lambda), \\ \|(\lambda - \Delta)^{-1} V_s\|_{\mathcal{L}(L_s^\infty(M))} &= \|(\lambda - \Delta)^{-1} V_s\|_{\mathcal{L}(L^\infty(M))} \leq \varepsilon(\lambda). \end{aligned}$$

Therefore, by the Riesz-Thorin interpolation theorem,

$$(2.19) \quad \|(\lambda - \Delta)^{-1} V_s\|_{\mathcal{L}(L_s^2(M))} \leq \varepsilon(\lambda),$$

which translates to

$$(2.20) \quad \|V_s^{1/2}(\lambda - \Delta)^{-1} V_s^{1/2}\|_{\mathcal{L}(L^2(M))} \leq \varepsilon(\lambda).$$

Hence, for  $f \in L^2(M)$ ,

$$(2.21) \quad \begin{aligned} \|(\lambda - \Delta)^{-1/2} V_s^{1/2} f\|_{L^2(M)}^2 &= ((\lambda - \Delta)^{-1/2} V_s^{1/2} f, (\lambda - \Delta)^{-1/2} V_s^{1/2} f) \\ &= (V_s^{1/2}(\lambda - \Delta)^{-1} V_s^{1/2} f, f) \\ &\leq \varepsilon(\lambda) \|f\|_{L^2(M)}^2, \end{aligned}$$

so

$$(2.22) \quad \|(\lambda - \Delta)^{-1/2} V_s^{1/2}\|_{\mathcal{L}(L^2(M))} \leq \varepsilon(\lambda)^{1/2},$$

and, by duality,

$$(2.23) \quad \|V_s^{1/2}(\lambda - \Delta)^{-1/2}\|_{\mathcal{L}(L^2(M))} \leq \varepsilon(\lambda)^{1/2}.$$

Hence, for  $u \in L^2(M)$ ,

$$(2.24) \quad \begin{aligned} &(V_s(\lambda - \Delta)^{-1/2} u, (\lambda - \Delta)^{-1/2} u) \\ &= (V_s^{1/2}(\lambda - \Delta)^{-1/2} u, V_s^{1/2}(\lambda - \Delta)^{-1/2} u) \\ &\leq \varepsilon(\lambda) \|u\|_{L^2(M)}^2. \end{aligned}$$

Let us set

$$(2.25) \quad \varphi = (\lambda - \Delta)^{-1/2} u, \quad \|\varphi\|_{H^1(M)} = \|(1 - \Delta)^{1/2} \varphi\|_{L^2(M)}.$$

Then (2.24) gives (for  $\lambda > 1$ )

$$\begin{aligned}
 (2.26) \quad \int_M V_s \varphi^2 d\omega &\leq \varepsilon(\lambda) \|(\lambda - \Delta)^{1/2} \varphi\|_{L^2}^2 \\
 &= \varepsilon(\lambda) ((\lambda - \Delta)\varphi, \varphi) \\
 &= \varepsilon(\lambda) \left\{ \|\varphi\|_{H^1(M)}^2 + (\lambda - 1) \|\varphi\|_{L^2(M)}^2 \right\}.
 \end{aligned}$$

Taking  $\lambda = 1$  and applying Cauchy's inequality gives, for  $\varphi, \psi \in H^1(M)$ ,

$$(2.27) \quad \left| \int_M V_s \varphi \psi d\omega \right| \leq \varepsilon(1) \|\varphi\|_{H^1(M)} \|\psi\|_{H^1(M)},$$

hence  $M_{V_s} : H^1(M) \rightarrow H^{-1}(M)$  satisfies

$$(2.28) \quad \|M_{V_s} \varphi\|_{H^{-1}(M)} \leq \varepsilon(1) \|\varphi\|_{H^1(M)},$$

an estimate that is independent of  $s > 0$ . Taking  $s \searrow 0$  and applying a little basic functional analysis yields the following.

**Proposition 2.2.** *Let  $\mu \in \mathcal{K}(M)$  be a positive measure, and set  $\nu = \mu + \omega$ . Then*

$$(2.29) \quad M_\nu : H^1(M) \longrightarrow H^{-1}(M),$$

and, for all  $\varphi \in H^1(M)$ ,  $\lambda > 1$ ,

$$(2.30) \quad (M_\nu \varphi, \varphi) \leq \varepsilon(\lambda) \left\{ \|\varphi\|_{H^1(M)}^2 + (\lambda - 1) \|\varphi\|_{L^2(M)}^2 \right\},$$

with  $\varepsilon(\lambda)$  as in (2.11). Consequently,

$$(2.31) \quad M_\mu : H^1(M) \longrightarrow H^{-1}(M),$$

with a similar estimate.

If we denote by  $Q_\lambda(\varphi)$  the quantity in brackets on the right side of (2.30), then Cauchy's inequality yields

$$(2.32) \quad |(M_\nu \varphi, \psi)| \leq \varepsilon(\lambda) Q_\lambda(\varphi)^{1/2} Q_\lambda(\psi)^{1/2},$$

for  $\varphi, \psi \in H^1(M)$ ,  $\lambda \geq 1$ . Recalling the reduction to (2.9) via the Hahn decomposition  $\mu = \mu^+ - \mu^-$ , we have the following.



**Corollary 2.3.** *Let  $\mu \in \mathcal{K}(M)$  be a finite signed measure. Then (2.31) holds, and*

$$(2.33) \quad |(M_\mu \varphi, \psi)| \leq \varepsilon(\lambda) Q_\lambda(\varphi)^{1/2} Q_\lambda(\psi)^{1/2}.$$

We now discuss how to use these estimates to define  $-\Delta + M_\mu$  as a self-adjoint operator, when  $\mu \in \mathcal{K}(M)$  is a finite, signed measure. Let us rewrite the case  $\psi = \varphi$  of (2.33) as

$$(2.34) \quad |(M_\mu \varphi, \varphi)| \leq \varepsilon \|\varphi\|_{H^1(M)}^2 + C(\varepsilon) \|\varphi\|_{L^2(M)}^2.$$

Set  $C^\#(\varepsilon) = C(\varepsilon) + 1$ . It follows that

$$(2.35) \quad [(-\Delta + M_\mu + C^\#(1/2))\varphi, \varphi] \geq \frac{1}{2} \|\varphi\|_{H^1(M)}^2.$$

This together with (2.34) implies (e.g., via the Lax-Milgram theorem) that the self-adjoint map

$$(2.36) \quad -\Delta + M_\mu + C^\#(1/2) : H^1(M) \longrightarrow H^{-1}(M) \text{ is bijective,}$$

with inverse

$$(2.37) \quad [-\Delta + M_\mu + C^\#(1/2)]^{-1} : H^{-1}(M) \longrightarrow H^1(M),$$

which is self adjoint. Restriction to  $L^2(M)$  yields

$$(2.38) \quad [-\Delta + M_\mu + C^\#(1/2)]^{-1} : L^2(M) \longrightarrow L^2(M),$$

self adjoint and compact, and injective. We can hence apply the classical theory of von Neumann/Friedrichs, to deduce the following.

**Proposition 2.4.** *Let  $\mu \in \mathcal{K}(M)$  be a finite signed measure. Then*

$$(2.39) \quad A = -\Delta + M_\mu + C^\#(1/2)$$

*is a positive self-adjoint operator on  $L^2(M)$ , with domain*

$$(2.40) \quad \mathcal{D}(A) = \{u \in H^1(M) : -\Delta u + M_\mu u \in L^2(M)\}.$$

*This self-adjoint operator has compact resolvent, hence discrete spectrum. Furthermore,*

$$(2.41) \quad \mathcal{D}(A^{1/2}) = H^1(M).$$

### 3. More singular potentials

Take  $M$  as before, and take an operator  $L$ , satisfying

$$(3.1) \quad L : H^1(M) \xrightarrow{\approx} H^{-1}(M), \quad \text{positive and self adjoint.}$$

For example, we might have  $L = -\Delta + 1$ , or more generally

$$(3.2) \quad L = -\Delta + M_\mu + C, \quad \mu \in \mathcal{K}(M), \text{ finite signed measure,}$$

and  $C \in (0, \infty)$  sufficiently large. Such  $L$  gives rise to a positive, self-adjoint operator on  $L^2(M)$ , with  $\mathcal{D}(L^{1/2}) = H^1(M)$ . To go further, take a finite, positive measure  $\nu$  on  $M$  such that

$$(3.3) \quad M_\nu : H^1(M) \longrightarrow H^{-1}(M).$$

Next, take a Borel function  $F$  on  $M$ , satisfying

$$(3.4) \quad F \geq 1, \quad F \in L^1(M, \nu).$$

Consider the space

$$(3.5) \quad \mathcal{H}_{F\nu}(M) = \{\varphi \in H^1(M) : \varphi \in L^2(M, F\nu + \omega)\}$$

(where, recall,  $\omega$  is volume measure on  $M$ ), and set

$$(3.6) \quad Q(\varphi, \psi) = (L\varphi, \psi) + \int_M \varphi\psi (F d\nu + d\omega),$$

for  $\varphi, \psi \in \mathcal{H}_{F\nu}(M)$ . We aim to prove the following.

**Proposition 3.1.** *The quadratic form  $Q$  defined by (3.6), with form domain  $\mathcal{H}_{F\nu}(M)$ , is closed.*

This result has the following implication.

**Proposition 3.2.** *Under the hypotheses given in (3.1)–(3.6), we have a self-adjoint operator*

$$(3.7) \quad A = L + M_{F\nu},$$

satisfying

$$(3.8) \quad \mathcal{D}(A^{1/2}) = \mathcal{H}_{F\nu}(M),$$

and

$$(3.9) \quad \mathcal{D}(A) = \{\varphi \in \mathcal{H}_{F\nu}(M) : L\varphi + FM\nu\varphi \in L^2(M)\}.$$

*Proof of Proposition 3.1.* Take  $\varphi_k \in \mathcal{H}_{F\nu}(M)$ , and assume

$$(3.10) \quad (\varphi_k) \text{ is Cauchy in the } Q\text{-norm.}$$

Equivalently,

$$(3.11) \quad (\varphi_k) \text{ is Cauchy in the } H^1(M)\text{-norm,}$$

and

$$(3.12) \quad (\varphi_k) \text{ is Cauchy in the } L^2(M, F\nu + \omega)\text{-norm.}$$

Our task is to show that  $(\varphi_k)$  converges in the  $Q$ -norm to an element of  $\mathcal{H}_{F\nu}(M)$ . Now (3.11)–(3.12) imply that there exist

$$(3.13) \quad \varphi \in H^1(M) \quad \text{and} \quad \psi \in L^2(M, F\nu + \omega)$$

such that

$$(3.14) \quad \begin{aligned} \varphi_k &\rightarrow \varphi, \quad \text{in } H^1(M), \text{ and} \\ \varphi_k &\rightarrow \psi, \quad \text{in } L^2(M, F\nu + \omega). \end{aligned}$$

Now, thanks to (3.3), the first part of (3.14) implies that

$$(3.15) \quad \varphi_k \longrightarrow \varphi \text{ in } L^2(M, \nu + \omega).$$

Meanwhile, the second part of (3.14) implies

$$(3.16) \quad \varphi_k \longrightarrow \psi \text{ in } L^2(M, \nu + \omega).$$

Hence

$$(3.17) \quad \varphi = \psi, \quad (\nu + \omega)\text{-a.e., on } M,$$

so we have

$$(3.18) \quad \varphi_k \longrightarrow \varphi \text{ in } Q\text{-norm,}$$

and Proposition 3.1 is proved.

#### 4. Schrödinger operators defined via Dirichlet forms

A number of papers, cited in “more references,” treat self-adjoint extensions of  $-\Delta + M_\mu$  for substantially more singular potentials than discussed in §§2–3. The approach involves more general “Dirichlet forms.” To introduce this approach in a simple fashion, we start with a quadratic form

$$(4.1) \quad Q(\varphi, \psi) = \int_M \langle \nabla \varphi, \nabla \psi \rangle A^2 d\omega,$$

where

$$(4.2) \quad A \in C^\infty(M), \quad A > 0.$$

We have

$$(4.3) \quad \begin{aligned} Q(\varphi, \psi) &= \int_M \langle \nabla \varphi, A^2 \nabla \psi \rangle d\omega \\ &= \int_M \langle \varphi, \nabla(A^2 \psi) - (\nabla A^2) \psi \rangle d\omega \\ &= - \int_M [\Delta \varphi + 2A^{-1} \nabla A \cdot \nabla \varphi] \psi A^2 d\omega. \end{aligned}$$

Under the hypotheses in (4.2), we see that  $Q$  is a closed quadratic form, with form domain  $H^1(M)$ , determining the operator

$$(4.4) \quad L = \Delta + 2A^{-1} \nabla A \cdot \nabla,$$

as a self-adjoint operator on  $L^2(M, A^2 \omega)$ , with domain  $H^2(M)$ .

Now we have the unitary operator

$$(4.5) \quad M_A : L^2(M, A^2 \omega) \xrightarrow{\approx} L^2(M, \omega),$$

given by  $M_A \varphi = A\varphi$ , giving rise to a self-adjoint operator

$$(4.6) \quad M_A L M_A^{-1} \text{ on } L^2(M, \omega),$$

also with domain  $H^2(M)$ . A calculation gives

$$(4.7) \quad \Delta(A\varphi) = A\Delta\varphi + 2\langle \nabla A, \nabla \varphi \rangle + (\Delta A)\varphi,$$

hence

$$(4.8) \quad A^{-1}\Delta(A\varphi) = \Delta\varphi + 2A^{-1}\nabla A \cdot \nabla\varphi + A^{-1}(\Delta A)\varphi,$$

so

$$(4.9) \quad M_A^{-1}\Delta M_A = L + M_{A^{-1}(\Delta A)},$$

hence

$$(4.10) \quad M_A L M_A^{-1} = \Delta - M_{A^{-1}(\Delta A)}.$$

If

$$(4.11) \quad A = (-\Delta + 1)^{-1}\gamma,$$

with  $\gamma \in C^\infty(M)$ ,  $\gamma \geq 0$  ( $\gamma \neq 0$ ), we are in the setting of (4.2), and (4.10) becomes

$$(4.12) \quad \begin{aligned} M_A L M_A^{-1} &= \Delta - M_{A^{-1}(\Delta A - A)} - 1 \\ &= \Delta + M_{A^{-1}\gamma} - 1. \end{aligned}$$

Now if (4.11) holds for smooth  $\gamma \geq 0$ , then the calculations leading to (4.12) are simply a hugely indirect way of defining  $\Delta + M_{A^{-1}\gamma}$  as a self-adjoint operator on  $L^2(M, \omega)$ . To go further, we want to take  $\gamma$  to be singular.

For example, by Proposition 2.1,

$$(4.13) \quad \gamma \in \mathcal{K}(M), \gamma \geq 0, \gamma \neq 0 \implies A \in C(M), A > 0.$$

Under these conditions, again  $Q$  in (4.1) is a closed quadratic form, with form domain  $H^1(M)$ , yielding a self-adjoint operator  $L$  on  $L^2(M, A^2\omega)$  satisfying

$$(4.14) \quad -(L\varphi, \psi)_{L^2(M, A^2\omega)} = Q(\varphi, \psi),$$

with

$$(4.15) \quad \mathcal{D}((-L)^{1/2}) = H^1(M),$$

and

$$(4.16) \quad \mathcal{D}(L) = \{\varphi \in H^1(M) : |Q(\varphi, \psi)| \leq C(\varphi)\|\psi\|_{L^2(M, A^2\omega)}\}.$$

Now our principal object of interest is not  $L$ , but

$$(4.17) \quad L_A = M_A L M_A^{-1},$$

as a self-adjoint operator on  $L^2(M, \omega)$ , with the idea of extending the identity (4.12), and understanding its significance. In connection with this, we have from (4.13) that

$$(4.18) \quad \gamma \in \mathcal{K}(M), \gamma \geq 0, \gamma \neq 0 \implies A^{-1}\gamma \in \mathcal{K}(M),$$

so the results of §2 apply to the operator on the right side of (4.12), in this setting. For a clear connection from this to (4.14)–(4.15), it would be useful to have

$$(4.19) \quad M_A, M_A^{-1} : H^1(M) \longrightarrow H^1(M).$$

To see this, take  $\varphi \in H^1(M)$ . We have

$$(4.20) \quad \begin{aligned} \nabla(A\varphi) &= A\nabla\varphi + (\nabla A)\varphi, \\ \nabla(A^{-1}\varphi) &= A^{-1}\nabla\varphi - A^{-2}(\nabla A)\varphi, \end{aligned}$$

so (4.19) holds provided

$$(4.21) \quad M_{\nabla A} : H^1(M) \longrightarrow L^2(M).$$

That this holds is a consequence of Theorem 11.1.1 of [MS], which we state.

**Proposition 4.1.** *Given  $\gamma \in \mathcal{D}'(M)$ ,  $A = (-\Delta + 1)^{-1}\gamma$ , we have*

$$(4.22) \quad M_\gamma : H^1(M) \rightarrow H^{-1}(M) \iff M_{\nabla A} : H^1(M) \rightarrow L^2(M).$$

Actually, the result of [MS] is done in the Euclidean space setting. We should look into checking the manifold case. We note the following related result, Theorem 11.2.2 of [MS] (again, formulated here in the manifold setting).

**Proposition 4.2.** *Given  $\gamma \in \mathcal{D}'(M)$ ,  $B = (-\Delta + 1)^{-1/2}\gamma$ , we have*

$$(4.23) \quad M_\gamma : H^1(M) \rightarrow H^{-1}(M) \iff M_B : H^1(M) \rightarrow L^2(M).$$

The proof in [MS] makes use of a result of [MV] that suggests the following.

**Claim 4.3.** *Let  $B \in L^2(M)$ ,  $P \in OPS^0(M)$ . Then*

$$(4.24) \quad M_B : H^1(M) \rightarrow L^2(M) \implies M_{PB} : H^1(M) \rightarrow L^2(M).$$

In fact, the following result is suggested by Lemmas 12.1.3–12.1.4 of [MS] (which cites [MV] for proofs):

**Claim 4.4.** *Assume*

$$(4.25) \quad p \in (1, \infty), \quad 0 < s < \frac{n}{p}.$$

Then, for  $B \in L^p(M)$ ,  $P \in OPS^0(M)$ ,

$$(4.26) \quad M_B : H^{s,p}(M) \rightarrow L^p(M) \implies M_{PB} : H^{s,p}(M) \rightarrow L^p(M).$$

According to Theorem 2.3.2 of [MS], given  $m \in \mathbb{N}$ ,  $p \in (1, \infty)$ , and  $B \in L^p(M)$ , we have  $M_B : H^{m,p}(M) \rightarrow L^p(M)$  if and only if there exists  $a < \infty$  such that, for each compact set  $K$ ,

$$(4.27) \quad \int_K |B|^p dV(x) \leq aC_{p,m}(K),$$

where  $C_{p,m}(K)$  is a certain capacity. One might also want to check whether this characterization (in case  $m = 1, p = 2$ ) is directly applicable to (4.21).

So we see that the approach via Dirichlet forms applies to potentials given by measures in the Kato class. A number of works have gone beyond this, and treated more singular potentials. Papers that do this include [AM], [ABR], [BM], [Br], [BrT], [Her], [KT]. Unfortunately, lots of these papers are tucked away in conference proceedings and are hard to access.

## 5. Point scatterers in 2D and 3D

There seem to be treatments of this topic using Dirichlet forms in papers cited at the end of §4, but I have not been able to verify their claims.

Another approach to point scattering in 2D and 3D can be found in [RU] and [U]. I have not figured out how these papers tie in with other works.



### A. Measures that multiply $H^1(M)$ to $H^{-1}(M)$

Let  $M$  be a compact,  $n$ -dimensional Riemannian manifold, and  $\Omega \subset M$  a connected open subset. Let  $\mu$  be a finite, positive measure on  $\overline{\Omega}$ . We want to give conditions that imply

$$(A.1) \quad M_\mu : H_0^1(\Omega) \longrightarrow H^{-1}(\Omega), \quad \text{compactly, where } M_\mu f = f\mu.$$

It will actually be convenient to let  $\mu$  be a positive finite measure on  $M$ , and ask when

$$(A.2) \quad M_\mu : H^1(M) \longrightarrow H^{-1}(M), \quad \text{compactly.}$$

Since  $H_0^1(\Omega)$  is a closed linear subspace of  $H^1(M)$ , we always get (A.1) from (A.2), by restriction (also restricting  $\mu$  to  $\overline{\Omega}$ ). We emphasize that the support of  $\mu$  can have nonempty intersection with  $\partial\Omega$ .

Note that, if  $\dim M = n \geq 2$ ,

$$(A.3) \quad f, g \in H^1(M) \Rightarrow fg \in H^{1,p}(M), \quad \text{for } p = \frac{n}{n-1} \text{ if } n \geq 3, \\ \text{for all } p < 2 \text{ if } n = 2.$$

Since  $C(M)$  is dense in the duals of such spaces  $H^{1,p}(M)$ , we deduce that (A.2) holds under the following conditions.

$$(A.4) \quad \mu \in H^{1,n/(n-1)}(M)^* = H^{-1,n}(M), \quad \text{if } n \geq 3, \\ \mu \in H^{-1,r}(M), \quad \text{for some } r > 2, \quad \text{if } n = 2.$$

We can use this to obtain the following class of singular measures satisfying (A.2).

**Proposition A.1.** *Let a compact  $S \subset M$  be locally a Lipschitz graph, of dimension  $n - 1$ , equipped with surface measure, i.e.,  $(n - 1)$ -dimensional Hausdorff measure,  $\sigma_S$ . Let*

$$(A.5) \quad \mu = h \sigma_S,$$

with

$$(A.6) \quad h \in L^{n-1}(S, d\sigma_S), \quad \text{if } n \geq 3, \\ h \in L^{1+\delta}(S, d\sigma_S), \quad \text{for some } \delta > 0, \text{ if } n = 2.$$

Then (A.4) holds, hence (A.2) holds.

*Proof.* We can apply the trace theorem, followed by the embedding theorem,

$$(A.7) \quad \text{Tr} : H^{1,p}(M) \longrightarrow B_{p,p}^s(S) \subset L^q(S),$$

with

$$(A.8) \quad \begin{aligned} p &= \frac{n}{n-1}, & s &= 1 - \frac{1}{p}, \\ q &= \frac{(n-1)p}{n-1-sp} = \frac{n-1}{n-2}, & q' &= n-1, \end{aligned}$$

if  $n \geq 3$ . For  $n = 2$ , (A.7) applies for all  $p \in (1, 2)$ , again with  $s = 1 - 1/p$ , and taking  $p \nearrow 2$  yields  $q \nearrow \infty$ , hence  $q' \searrow 1$ . Thus (A.5)–(A.6) imply (A.4).

NOTE. The trace result is perhaps better known when  $S$  is smooth. (Cf. [BL], Theorem 6.6.1.) However, all the function spaces involved are invariant under bi-Lipschitz maps.

Positive measures satisfying (A.2) can have much wilder support than a Lipschitz surface. For example, one can take an infinite sequence of measures  $\mu_k$  satisfying the hypotheses of Proposition A.1, supported on surfaces  $S_k$ , and set  $\mu = \sum_{k=1}^{\infty} a_k \mu_k$ , with positive  $a_k$  decreasing sufficiently fast.

Here is another class of examples. Let  $\mathcal{O} \subset M$  be an open set whose boundary  $\partial\mathcal{O}$  is locally the graph of a continuous function. Then one can take a smooth vector field  $X$  on  $M$ , vanishing nowhere on  $\partial\mathcal{O}$ , whose flow  $\mathcal{F}_X^t$  has the property that, for each  $y \in \partial\mathcal{O}$ ,  $\mathcal{F}_X^t y$  belongs to  $\mathcal{O}$  for small  $t > 0$ . Then

$$(A.9) \quad \mu = X\chi_{\mathcal{O}}$$

is a positive measure, supported by  $\partial\mathcal{O}$ , and it belongs to  $H^{-1,\infty}(M)$ . The positivity of  $\mu$  is a consequence of the fact that

$$(A.10) \quad \chi_{\mathcal{O}} \circ \mathcal{F}_X^t \geq \chi_{\mathcal{O}},$$

for all small  $t > 0$ , since

$$(A.11) \quad t^{-1}(\chi_{\mathcal{O}} \circ \mathcal{F}_X^t - \chi_{\mathcal{O}}) \longrightarrow X\chi_{\mathcal{O}} \text{ in } \mathcal{D}'(M), \text{ as } t \rightarrow 0.$$

In the last class of examples, (A.9), the support of  $\mu$  has topological dimension  $n - 1$ , but its Hausdorff dimension can be  $> n - 1$ . We next produce measures satisfying (A.4) and supported on “fractal” sets of Hausdorff dimension  $< n - 1$ . We make use of the following result, contained in Theorem 4.7.4 of [Zie]. Here,  $B_r(x)$  denotes the ball of radius  $r$  centered at  $x$ .

**Lemma.** *Let  $\mu$  be a positive measure on  $M$  with the property that there exist  $A < \infty$  and  $\varepsilon > 0$  such that*

$$(A.12) \quad \mu(B_r(x)) \leq Ar^{n-q+\varepsilon}, \quad \forall r \in (0, 1], x \in M.$$

*Assume  $q \in (1, n)$ . Then*

$$(A.13) \quad \mu \in H^{-1,p}(M), \quad p = q'.$$

We then see that (A.4) holds whenever

$$(A.14) \quad \mu(B_r(x)) \leq Ar^\alpha, \quad \alpha > n - 1 - \frac{1}{n-1}.$$

In particular, for  $n = 2$ , it suffices to have (A.14) for some  $\alpha > 0$ .

We will give some explicit examples of a compactly supported measure on  $\mathbb{R}^2$  satisfying (A.14). It will be clear that many other examples can be constructed. We start with the Cantor middle third set  $\mathcal{K} \subset [0, 1]$ . Now put  $[0, 1] \subset \mathbb{R} \subset \mathbb{R}^2$ , say as part of the  $x$ -axis, so now  $\mathcal{K} \subset \mathbb{R}^2$ . As is well known (cf. [T], p. 170), there is the  $\alpha$ -dimensional Hausdorff measure computation

$$(A.15) \quad \mathcal{H}^\alpha(\mathcal{K}) = \gamma_\alpha, \quad \text{for } \alpha = \frac{\log 2}{\log 3} \approx 0.6309,$$

with  $0 < \gamma_\alpha < \infty$  (in fact,  $\gamma_\alpha = \pi^{\alpha/2} 2^{-\alpha} / \Gamma(\alpha/2 + 1)$ ). Set

$$(A.16) \quad \mu = \mathcal{H}^\alpha \llcorner \mathcal{K},$$

i.e.,  $\mu(S) = \mathcal{H}^\alpha(\mathcal{K} \cap S)$ , for Borel sets  $S \subset \mathbb{R}^2$ . The self similarity of  $\mathcal{K}$  enables one to show that

$$(A.17) \quad \mu(B_{3^{-k}}(x)) \leq C 2^{-k},$$

which readily leads to (A.14), with  $\alpha$  as in (A.15).

The Cantor middle third set described above is one of a family of Cantor sets  $\mathcal{K}(\vartheta) \subset [0, 1]$ , defined for  $\vartheta \in (0, 1)$  as follows. Remove from  $[0, 1] = I$  the open interval of length  $\vartheta \ell(I)$ , with the same center as  $I$ , and repeat this process with the other closed subintervals that remain. (Thus  $\mathcal{K} = \mathcal{K}(1/3)$ .) This time (cf. [T], p. 171), one has

$$(A.18) \quad \mathcal{H}^\alpha(\mathcal{K}(\vartheta)) = \gamma_\alpha, \quad \alpha = \frac{\log 2}{\log b}, \quad b = \frac{2}{1-\vartheta},$$

and again self-similarity yields

$$(A.19) \quad \mu(B_r(x)) \leq Ar^\alpha,$$

with  $\alpha$  as in (A.18), when

$$(A.20) \quad \mu = \mathcal{H}^\alpha \lfloor \mathcal{K}(\vartheta).$$

Note that

$$(A.21) \quad \vartheta \searrow 0 \Rightarrow \alpha \nearrow 1, \quad \vartheta \nearrow 1 \Rightarrow \alpha \searrow 0.$$

As before, we put  $\mathcal{K}(\vartheta) \subset [0, 1] \subset \mathbb{R} \subset \mathbb{R}^2$ , and regard  $\mu$  in (A.20) as a compactly supported measure on  $\mathbb{R}^2$ . Thus the push-forward of  $\mu$  to a measure on a compact two-dimensional manifold  $M$ , via a locally bi-Lipschitz map, yields a measure on  $M$  satisfying (A.4), hence (A.2), whenever  $0 < \vartheta < 1$ .

One way to get measures on higher dimensional spaces satisfying (A.14) is to take products. Say  $n = n_1 + n_2$  and  $\mu_j$  are compactly supported measures on  $\mathbb{R}^{n_j}$  satisfying

$$(A.22) \quad \mu_j(B_r(x_j)) \leq C_j r^{\alpha_j}, \quad j = 1, 2, \quad x_j \in \mathbb{R}^{n_j}.$$

If  $x = (x_1, x_2) \in \mathbb{R}^n$ , note that  $B_r(x) \subset B_r(x_1) \times B_r(x_2)$ , so if

$$(A.23) \quad \mu = \mu_1 \times \mu_2$$

is the product measure on  $\mathbb{R}^n$ , we have

$$(A.24) \quad \mu(B_r(x)) \leq \mu_1(B_r(x_1))\mu_2(B_r(x_2)) \leq C_1 C_2 r^{\alpha_1 + \alpha_2}, \quad x \in \mathbb{R}^n.$$

If  $\alpha = \alpha_1 + \alpha_2$  satisfies the condition on  $\alpha$  in (A.14), we get a compactly supported measure on  $\mathbb{R}^n$  whose push-forward to a measure on an  $n$ -dimensional compact manifold  $M$ , via a locally bi-Lipschitz map, satisfies (A.14), hence (A.4), hence (A.2).

For example, take  $\vartheta_j \in (0, 1)$ , and set

$$(A.25) \quad \mu_j = \mathcal{H}^{\alpha_j} \lfloor \mathcal{K}(\vartheta_j),$$

with  $\alpha_j$  as in (A.18), i.e.,

$$(A.26) \quad \alpha_j = \frac{\log 2}{\log b_j}, \quad b_j = \frac{2}{1 - \vartheta_j}.$$

We regard  $\mu_1$  as a measure on  $\mathbb{R}^2$  and  $\mu_2$  as a measure on  $\mathbb{R}$ , via  $\mathcal{K}(\vartheta_1) \subset [0, 1] \subset \mathbb{R} \subset \mathbb{R}^2$  and  $\mathcal{K}(\vartheta_2) \subset [0, 1] \subset \mathbb{R}$ . Thus  $\mu = \mu_1 \times \mu_2$  is a compactly supported measure on  $\mathbb{R}^3$  (actually supported on a 2D linear subspace of  $\mathbb{R}^3$ ). In this case, the condition for (A.18) to hold is

$$(A.27) \quad \alpha_1 + \alpha_2 > \frac{3}{2}.$$

Looking at (A.15), we see that (A.27) fails when  $\vartheta_1 = \vartheta_2 = 1/3$ . In case  $\vartheta_1 = \vartheta_2 = \vartheta$ , the condition that (A.27) hold is that

$$(A.28) \quad 2 \frac{\log 2}{\log 2 - \log(1 - \vartheta)} > \frac{3}{2},$$

or equivalently,

$$(A.29) \quad \vartheta < 1 - 2^{-1/3} \approx 0.2063.$$

NOTE. When the measures  $\mu_j$  are given by (A.25),  $\mu = \mu_1 \times \mu_2$  is supported on  $\mathcal{K}(\vartheta_1) \times \mathcal{K}(\vartheta_2)$ , but it is generally *not*  $(\alpha_1 + \alpha_2)$ -dimensional Hausdorff measure on this set. See pp. 70–74 of [Fal] for a discussion of this matter.

Returning to generalities, we mention that while (A.12) is a sufficient condition for (A.13), it is not quite necessary. There is a (somewhat more elaborate) necessary and sufficient condition for (A.13) to hold, provided  $1 < q < n$ , given in terms of an estimate on  $\mu(B_r(x))$ . See Theorem 4.7.5 of [Zie]. The condition is subtly weaker than (A.12).

We also mention conditions for

$$(A.30) \quad M_\mu : H^1(M) \longrightarrow H^{-1}(M),$$

i.e., (A.2) but disregarding compactness. If we take  $f \in H^1(M)$  and pair  $M_\mu f$  with  $f$ , since  $\mu$  is a positive measure, a necessary condition for (A.30) is that

$$(A.31) \quad \int_M |f|^2 d\mu \leq C \|f\|_{H^1(M)}^2, \quad \forall f \in H^1(M).$$

Applying Cauchy's inequality to  $\int_M f \bar{g} d\mu$ , we see that (A.31) is also sufficient for (A.30). Furthermore, it follows from Theorem 1.2.2 of [MS] that the existence of  $A < \infty$  such that

$$(A.32) \quad \mu(S) \leq A \text{Cap}(S), \quad \forall \text{Borel } S \subset M$$

is necessary and sufficient for (A.31) to hold, hence for (A.30) to hold, where  $\text{Cap}(S)$  is a variant of electrostatic capacity, appropriate for this situation.

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