

Non-Equilibrium Statistical Mechanics of Strongly Anharmonic Chains of Oscillators

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Abstract: We study the model of a strongly non-linear chain of particles coupled to two heat baths at different temperatures. Our main result is the existence and uniqueness of a stationary state at all temperatures. This result extends those of Eckmann, Pillet, Rey-Bellet [EPR99a, EPR99b] to potentials with essentially arbitrary growth at infinity. This extension is possible by introducing a stronger version of Hörmander’s theorem for Kolmogorov equations to vector fields with polynomially bounded coefficients on unbounded domains.

1. Introduction

In this paper, we study the statistical mechanics of a highly non-linear chain of oscillators coupled to two heat reservoirs which are at (arbitrary) different temperatures. We show that such systems have, under suitable conditions, a *unique* stationary state, in which heat flows from the hotter reservoir to the cooler one.

These results are an extension of the same statements obtained by Eckmann, Pillet and Rey-Bellet in [EPR99a, EPR99b] where it was assumed that the Hamiltonian is essentially “quadratic at high energies”. Since quadratic Hamiltonians have been discussed much earlier by Lebowitz and Spohn [LS77], there is an issue here of whether the quadratic nature of the forces at infinite energies is an essential ingredient of existence and uniqueness of the stationary state. Our result shows that this is not the case, since we allow for potentials of arbitrary polynomial growth.

Our models, which are described in Sect. 2, treat a Hamiltonian of the form

$$H_S(p, q) = \sum_{i=0}^N \left(\frac{p_i^2}{2} + V_1(q_i) \right) + \sum_{i=1}^N V_2(q_i - q_{i-1}),$$

describing a chain of particles with nearest-neighbor interaction (see Fig. 3.1). This chain is linearly coupled to heat baths B_i represented by free fields at temperatures T_i .

We proceed then, as in [EPR99a], to a reduction to a stochastic differential equation, see (2.2). Associated with it is an “effective energy” G , described in (2.6), which is equal to H_S with some quadratic terms from the heat baths added. The generator corresponding to the stochastic differential equation above, represented in a space weighted with an exponential of G , will be called K and is the main object of study of this paper.

It is for this generator that we show existence and uniqueness of an invariant state. This will be done by first showing that K has compact resolvent (which is really more than needed), and then using this result to derive the properties of the invariant measure. Our conditions on H_S are spelled out in Sect. 3 below. They basically say that the coupling *between* the oscillators must be stronger than the single particle potential. This condition might be physically relevant, since it implies that transport is favored over storing of energy, but we have not found a counterexample when this condition is violated. Furthermore, the interparticle coupling must be convex.

The main technical insight behind our generalization of the results of [EPR99a, EPR99b] is a new, and stronger version of the Hörmander theorem for Kolmogorov equations. We will develop this in more generality in Sect. 5, but here we just indicate how we use this result. The operator K is of the form

$$K = \sum_{i=1}^n X_i^* X_i + X_0, \quad (1.1)$$

where the X_i are smooth vector fields on \mathbf{R}^d . For example, see Eq. (3.13), X_0 contains terms of the form $p_i \partial_{q_i}$ and $(\partial_{q_i} V) \partial_{p_i}$, where V is the interaction. The X_i for $i \neq 0$ are first order operators. Here, ∂V is polynomially bounded, whereas, in [EPR99a], ∂V was assumed to be linearly bounded. Letting g_0 be an adequate inverse power of the effective energy G , one successively considers the *finite* sets of operators

$$\mathcal{A}_{-1} = \{X_1, \dots, X_n\}, \quad \mathcal{A}_0 = \{g_0 X_0, X_1, \dots, X_n\},$$

and then – see Sect. 6 for the detailed definition –

$$\mathcal{A}_\ell = \mathcal{A}_{\ell-1} \cup [g_0 X_0, \mathcal{A}_{\ell-1}].$$

We stop this iteration after at most $2N$ steps, where N is the number of particles in the chain, obtaining the set $\mathcal{A} = \mathcal{A}_{2N+1}$. We now define the operator $\Lambda_{\mathcal{A}}$ as the finite sum

$$\Lambda_{\mathcal{A}}^2 = 1 + \sum_{A \in \mathcal{A}} A^* A.$$

This is a generalization to our case of an elliptic operator of the type $\Lambda^2 = 1 - \sum_i \partial_{x_i}^2$ used in [Hör85] or $\Lambda^2 = 1 - \sum_i \partial_{x_i}^2 + \sum_i x_i^2$ used in [EPR99a].

With these definitions, one then has the bound

Proposition 1.1 (Momentum space bound). *There is a constant C such that for all $f \in C_0^\infty(\mathbf{R}^d)$ one has in L^2 :*

$$\|\Lambda_{\mathcal{A}}^{16-N} f\| \leq C (\|Kf\| + \|f\|). \quad (1.2)$$

We also derive a similar bound in the conjugate variables:

Proposition 1.2 (Position space bound). *There is a constant C such that for all $f \in C_0^\infty(\mathbf{R}^d)$ one has in L^2 :*

$$\|G^\varepsilon f\| \leq C(\|Kf\| + \|f\|), \tag{1.3}$$

where $\varepsilon > 0$ depends on the asymptotic behavior of the potential V and on N .

Combining these two propositions one easily shows that K has compact resolvent. Then one derives from that result the existence of an invariant measure. Its properties are then found adapting the techniques of [EPR99a,EPR99b].

The remainder of this paper is organized as follows. In Sect. 2 we describe the physical model and in Sect. 3 we refine the setting and state the results. Section 4 will be devoted to the proof of the position space bound (Proposition 3.7). In Sect. 5, we present in detail the general scheme for studying operators of the form of (1.1), and show the inequality corresponding to (1.2). This section is as much as possible self-contained as it presents some independent interest. The detailed application of this general scheme to the problem of the chain allows us to prove the momentum space bound (Proposition 3.8) in Sect. 6. In Sect. 7 we combine these two bounds and prove Theorem 3.6 showing that K has compact resolvent and hence discrete spectrum. In Sect. 8 we show existence, uniqueness, and further properties of the invariant measure (Theorem 3.9).

Appendix A contains some technical estimates used in Sect. 5. Appendix B contains the proof of a result concerning the domains of K and K^* . The method used there probably works for more general accretive second-order differential operators. Appendix C finally contains the proof of a technical result used in Sect. 8.

2. The Model

We will study the model of a (small) classical N -particle Hamiltonian system coupled to M stochastic heat baths proposed in [EPR99a]. The small system without the heat baths is governed by a Hamiltonian

$$H_S \in C^\infty(\mathbf{R}^{2N}).$$

(We stay here with $d = 1$ dimensional position space for each particle to simplify notation.) The heat baths are modeled by classical field theories associated to the wave equation. The fields will be called φ_i and their conjugate momenta π_i , where the index i ranges from 1 to M .

The Hamiltonian for one heat bath is given by

$$H_B(\pi, \varphi) = \frac{1}{2} \int_{\mathbf{R}} (|\partial\varphi|^2 + |\pi|^2) dx.$$

The couplings allowed for the model are linear in the field variables. The total Hamiltonian for our model is then given by

$$H(p, q, \pi, \varphi) = \sum_{i=1}^M \left(H_B(\pi_i, \varphi_i) + F_i(p, q) \int_{\mathbf{R}} \partial\varphi_i(x) \varrho_i(x) dx \right) + H_S(p, q). \tag{2.1}$$

We assume the initial conditions describe the heat baths at equilibrium at inverse temperatures β_i , i.e. they are distributed in a sense according to the measure with “weight”

$$e^{-\beta_i H_B(\pi_i, \varphi_i)}.$$

The paper [EPR99a] explains in detail how, and under which conditions on the coupling functions q_i , one can reduce the resulting “big” system to a “small” system, where the heat baths are described by a finite number of variables. The price to pay for that is that we are now dealing with the following system of stochastic differential equations:

$$\begin{aligned} dq_j &= \partial_{p_j} H_S dt - \sum_{i=1}^M (\partial_{p_j} F_i) r_i dt, & j = 1, \dots, N, \\ dp_j &= -\partial_{q_j} H_S dt + \sum_{i=1}^M (\partial_{q_j} F_i) r_i dt, & (2.2) \\ dr_i &= -\gamma_i r_i dt + \lambda_i^2 \gamma_i F_i(p, q) dt - \lambda_i \sqrt{2\gamma_i T_i} dw_i(t), & i = 1, \dots, M, \end{aligned}$$

where the w_i are independent Wiener processes. The various constants appearing in (2.2) have the following meaning. T_i is the temperature of the i^{th} heat bath, λ_i is the strength of the coupling between that heat bath and the small system and $1/\gamma_i$ is the relaxation time of the i^{th} heat bath. The value of γ_i depends on the choice of the coupling function q_i . If we wanted to be more general, we would have to introduce for each bath a family of auxiliary variables $r_{i,m}$ as is done in [EPR99a]. This would only cause notational problems and does not change our argument.

If we consider a generic n -dimensional system of stochastic differential equations with additive noise of the form

$$dx_i(t) = b_i(x(t)) dt + \sum_{j=1}^n \sigma_{ij} dw_j(t), \quad (2.3)$$

we can associate with it the second-order differential operator \mathcal{L} formally defined by

$$\mathcal{L} \equiv \frac{1}{2} \sum_{i,j=1}^n \partial_i (\sigma \sigma^T)_{ij} \partial_j + \sum_{i=1}^n b_i(x) \partial_i. \quad (2.4)$$

It is a classical result that if the solution of such a system of stochastic differential equations exists, the probability density of the solution satisfies the partial differential equation

$$\partial_t p(x, t) = (\mathcal{L}p)(x, t).$$

In our case, the differential operator \mathcal{L} is given by

$$\mathcal{L} = \sum_{i=1}^M \lambda_i^2 \gamma_i T_i \partial_{r_i}^2 - \sum_{i=1}^M \gamma_i (r_i - \lambda_i^2 F_i(p, q)) \partial_{r_i} + X^{H_S} - \sum_{i=1}^M r_i X^{F_i}, \quad (2.5)$$

where the symbol X^F denotes the Hamiltonian vector field associated to the function F . It is convenient to introduce the “effective energy” given by

$$G(p, q, r) = H_S(p, q) + \sum_{i=1}^M \left(\frac{r_i^2}{2\lambda_i^2} - F_i(p, q) r_i \right). \quad (2.6)$$

At this point, we make the following assumption on the asymptotic behavior of G .

Assumption 0. There exist constants $\tilde{d}_i, C > 0$ and $\alpha > 0$, as well as constants $\tilde{c}_i > 2/\lambda_i^2$ such that

$$H_S(p, q) \geq C(1 + \|p\|^\alpha + \|q\|^\alpha), \tag{2.7a}$$

$$F_i^2(p, q) \leq \tilde{c}_i H_S(p, q) + \tilde{d}_i. \tag{2.7b}$$

Remark 2.1. This assumption essentially means that the effective energy G grows at infinity at least like $1 + \|r\|^2 + \|p\|^\alpha + \|q\|^\alpha$. This implies the stability of the system, as follows easily from the inequality

$$|r_i F_i(p, q)| \leq s^2 r_i^2 + \frac{F_i^2(p, q)}{s^2},$$

which holds for every $s > 0$. In particular, this implies that $\exp(-\beta G)$ is integrable for every $\beta > 0$.

We also define

$$W \equiv \sum_{i=1}^M \gamma_i T_i,$$

which is, in some sense that will be clear in a moment, the maximal power the heat baths can pull into the chain. We have the following result.

Proposition 2.2. *Assume Assumption 0 holds. Then the solution $\xi(t; x_0, w)$ of (2.2) exists and is continuous for all $t > 0$ with probability 1. Moreover, the mean energy of the system satisfies for all values of t and x_0 the estimate*

$$\mathbf{E}[G(x(t; x_0, w))] - G(x_0) \leq Wt, \tag{2.8}$$

where $\mathbf{E}[\cdot]$ denotes the expectation with respect to the M -dimensional Wiener process w .

Remark 2.3. The bound (2.8) allows the energy to grow forever, which would cause the system to “explode”. But this is not the case for the systems we consider in this paper. Indeed, we will prove that the process possesses a unique stationary state. This implies among other features that the mean time needed to reach any compact region is finite, and so the energy can not grow forever.

Proof. A classical result (see e.g. [Has80, Thm 4.1]) states the following. Assume that the vector field b of (2.3) is locally Lipschitz and that there exists a confining \mathcal{C}^2 function $G : \mathbf{R}^n \rightarrow \mathbf{R}$ and a constant k such that

$$(\mathcal{L}G)(x) \leq k \quad \text{for all } x \in \mathbf{R}^n.$$

Then there exists a unique stochastic process $\xi(t)$ solving (2.3). The process ξ is regular (i.e. it does not blow up in a finite time) and continuous for all $t > 0$. It satisfies the statistics of a Markovian diffusion process with generator \mathcal{L} . Moreover, we have the estimate

$$\mathbf{E}[G(x(t; x_0, w))] - G(x_0) \leq kt.$$

This result can be applied to our case, if we take for G the effective energy defined in (2.6). An explicit computation yields indeed

$$\mathcal{L}G = W - \sum_{i=1}^M \frac{\gamma_i}{\lambda_i^2} (r_i - \lambda_i^2 F_i(p, q))^2. \tag{2.9}$$

Moreover, G is confining by Assumption 0. This proves the assertion. \square

2.1. Definition and simple properties of the semigroup. In this paper, we will mainly be interested in studying under which assumptions on the chain Hamiltonian H_S it is possible to prove the existence of a *unique invariant measure* for the stochastic process $\xi(t; x_0, w)$ solving (2.2). Throughout, we will use the notation

$$\mathcal{X} = \mathbf{R}^{2N+M}$$

for the extended phase space (p, q, r) . This stochastic process defines a semigroup \mathcal{T}^t on $C_0^\infty(\mathcal{X})$ by

$$\mathcal{T}^t f(x_0) = \mathbf{E}[f(\xi(t; x_0, w))]. \quad (2.10)$$

This semigroup satisfies the following

Proposition 2.4. *Assume Assumption 0 holds. Then \mathcal{T}^t extends to a strongly continuous, quasi-bounded semigroup of positivity preserving operators on $L^2(\mathcal{X})$. Its generator L is the closure of the operator \mathcal{L} with domain $C_0^\infty(\mathcal{X})$. The adjoint L^* is the closure of the formal adjoint \mathcal{L}^T with domain $C_0^\infty(\mathcal{X})$.*

Proof. The proof will be given in Appendix B. \square

This in turn defines a dual semigroup $(\mathcal{T}^t)^*$ by

$$\int (\mathcal{T}^t f)(x) v(dx) = \int f(x) ((\mathcal{T}^t)^* v)(dx).$$

The generator of $(\mathcal{T}^t)^*$ is given by the adjoint of \mathcal{L} in L^2 that will be denoted \mathcal{L}^T . It is possible to check that if the heat baths are all at the same temperature $T = 1/\beta$, we have

$$\mathcal{L}^T \mu_0 = 0, \quad \text{where} \quad \mu_0(p, q, r) = e^{-\beta G(p, q, r)}.$$

Thus, the generalized Gibbs measure

$$d\mu_0 = e^{-\beta G(p, q, r)} dp dq dr = \mu_0(p, q, r) dp dq dr,$$

is an invariant measure for the Markov process described by (2.2). This confirms our definition of G as the effective energy of our system. We want to consider the more interesting case where the temperatures of the heat baths are not the same. The idea is to work in a Hilbert space that is weighted with a Gibbs measure for some reference temperature.

We will therefore study an extension \mathcal{T}_0^t of \mathcal{T}^t acting on an auxiliary weighted Hilbert space \mathcal{H}_0 , given by

$$\mathcal{H}_0 \equiv L^2(\mathcal{X}, Z_0^{-1} e^{-2\beta_0 G(p, q, r)} dp dq dr),$$

where Z_0 is a normalization constant and β_0 is a “reference” inverse temperature that we choose such that

$$1/\beta_0 \equiv T_0 > \max\{T_i \mid i = 1, \dots, M\}. \quad (2.11)$$

We have the following

Proposition 2.5. *Assume Assumption 0 holds. Then the semigroup \mathcal{T}^t given by (2.10) extends to a strongly continuous quasi-bounded semigroup \mathcal{T}_0^t on \mathcal{H}_0 . Moreover, $\mathcal{T}_0^t 1 = 1$ and \mathcal{T}_0^t is positivity preserving, i.e.*

$$\mathcal{T}_0^t f \geq 0 \quad \text{if} \quad f \geq 0.$$

Let L_0 be the generator of \mathcal{T}_0^t . Then L_0 coincides on $C_0^\infty(\mathcal{X})$ with \mathcal{L} of (2.5) and $C_0^\infty(\mathcal{X})$ is a core for both L_0 and L_0^* .

Proof. The statement can be proven by simply retracing the proof of Lemma 3.1 in [EPR99a]. There are only three points that have to be checked. We define the vector fields b and b_0 respectively by

$$b = - \sum_{i=1}^M \gamma_i (r_i - \lambda_i^2 F_i(p, q)) \partial_{r_i} + X^{H_S} - \sum_{i=1}^M r_i X^{F_i},$$

$$b_0 = 2\beta_0 \sum_{i=1}^M \lambda_i^2 \gamma_i T_i (\partial_{r_i} G) \partial_{r_i} = 2\beta_0 \sum_{i=1}^M \gamma_i T_i (r_i - \lambda_i^2 F_i(p, q)) \partial_{r_i}.$$

In order to make the proof of [EPR99a] work, we have to check that

$$\|\operatorname{div} b\|_\infty < \infty, \quad \|\operatorname{div} b_0\|_\infty < \infty, \quad \sup_{x \in \mathcal{X}} (b + \frac{1}{2} b_0) G(x) < \infty,$$

where b and b_0 are considered as first-order differential operators in the last inequality. The divergence of any Hamiltonian vector field vanishes, and so we have

$$\|\operatorname{div} b\|_\infty = - \sum_{i=1}^M \gamma_i < \infty.$$

The term involving the divergence of b_0 can easily be computed to give

$$\|\operatorname{div} b_0\|_\infty = \beta_0 \sum_{i=1}^M \gamma_i T_i < \infty.$$

In order to check the last inequality, we compute the expression

$$(b + \frac{1}{2} b_0) G(p, q, r) = \sum_{i=1}^M \frac{\gamma_i}{\lambda_i^2} (\beta_0 T_i - 1) (r_i - \lambda_i^2 F_i(p, q))^2.$$

We see that condition (2.11) on β_0 obviously implies $\beta_0 T_i - 1 < 0$, and so the desired inequality holds.

The domains of L_0 and L_0^* are controlled by the techniques of Appendix B. \square

We are mainly interested in the case $M = 2$. The Hamiltonian H_S will describe a chain of $N + 1$ strongly anharmonic oscillators coupled to two heat baths at the first and the last particle. In the case in which the Hamiltonian H_S can be written as a quadratic function plus some bounded terms, the existence and uniqueness of a stationary state for every temperature difference has been proved in [EPR99a, EPR99b]. We will extend this result to the case where the potentials grow faster than quadratically at infinity. Besides some weak conditions on the derivatives of the one and two-body potentials, we will only require that they grow algebraically and that the two-body potentials grow asymptotically faster than the one-body potentials, i.e. at large separation the interaction energy between neighboring particles grows *faster* than the one-particle energy.

2.2. *Notations.* Throughout, the domain of an operator A will be denoted by $\mathcal{D}(A)$. Unless specified, the domain of any operator will always be the closure in the graph norm of \mathcal{C}_0^∞ . For example, if we write $[A, B]$, we mean in fact $\overline{(AB - BA) \upharpoonright \mathcal{C}_0^\infty}$, so that the domain of $[A, B]$ can be larger than that of A or B separately.

3. Setting and Results

In order to set up our model, we need to be able to describe precisely the growth rates of the potentials at infinity. This will be achieved with the following function spaces.

Definition 3.1. Choose $\alpha \in \mathbf{R}$. We call \mathcal{F}_α the set of all \mathcal{C}^∞ functions from \mathbf{R}^n to \mathbf{R} such that for every multi-index k there exists a constant C_k for which

$$\|D^k f(x)\| \leq C_k(1 + \|x\|^2)^{\alpha/2}, \quad \text{for all } x \in \mathbf{R}^n.$$

Definition 3.2. Choose $\alpha \in \mathbf{R}$ and $i \in \mathbf{N} \cup \{\infty\}$. We call \mathcal{F}_α^i the set of all \mathcal{C}^∞ functions from \mathbf{R}^n to \mathbf{R} such that for every multi-index k with $|k| \leq i$, we have $D^k f(x) \in \mathcal{F}_{\alpha-|k|}$.

Remark 3.3. For any $\alpha \in \mathbf{R}$, the function

$$\begin{aligned} P^\alpha : \mathbf{R}^n &\rightarrow \mathbf{R} \\ x &\mapsto (1 + \|x\|^2)^{\alpha/2} \end{aligned} \tag{3.1}$$

belongs to $\mathcal{F}_\alpha^\infty$. Moreover, any polynomial of degree n belongs to \mathcal{F}_n^∞ .

3.1. *The chain.* We consider the Hamiltonian

$$H_S(p, q) = \sum_{i=0}^N \left(\frac{p_i^2}{2} + V_1(q_i) \right) + \sum_{i=1}^N V_2(q_i - q_{i-1}), \tag{3.2}$$

describing a chain of particles with nearest-neighbor interaction (Fig. 3.1). We slightly modify the notations used so far. Because there are only two heat baths, we will not use for them the indices $i \in \{1, 2\}$, but rather $i \in \{L, R\}$. Concerning the coupling between the chain and the baths, we assume that we can make a dipole approximation, so we set

$$F_L = q_0 \quad \text{and} \quad F_R = q_N, \tag{3.3}$$

in Eq. (2.1). We will make the Assumptions 1–3 on V_1 and V_2 .

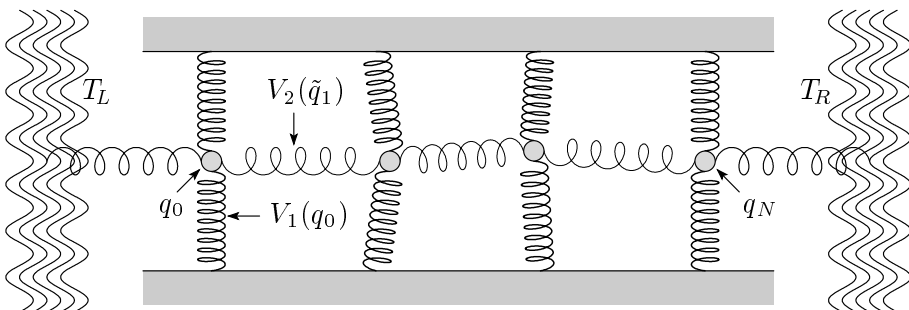


Fig. 3.1. Chain of oscillators

Assumption 1. The potential V_1 is in \mathcal{F}_{2n}^2 for some $n > 1$. Moreover, there are constants $c_i > 0$ such that

$$V_1(x) \geq c_1 P^{2n}(x), \quad (3.4a)$$

$$x V_1'(x) \geq c_2 P^{2n}(x) - c_3, \quad (3.4b)$$

for all $x \in \mathbf{R}$.

Assumption 2. The potential V_2 is in \mathcal{F}_{2m}^2 for some $m > n$. Moreover, there are constants $c'_i > 0$ such that

$$V_2(x) \geq c'_1 P^{2m}(x), \quad (3.5a)$$

$$x V_2'(x) \geq c'_2 P^{2m}(x) - c'_3, \quad (3.5b)$$

for all $x \in \mathbf{R}$.

Assumption 3. The function

$$x \mapsto \frac{1}{V_2''(x)}$$

belongs to \mathcal{F}_ℓ for some ℓ .

Remark 3.4. It is clear that (3.3), together with Assumptions 1 and 2 immediately imply Assumption 0. Notice that the assumptions $V_1 \in \mathcal{F}_{2n}^2$ and $V_2 \in \mathcal{F}_{2m}^2$ give bounds not only on the asymptotic behavior of V_1 and V_2 , but also of their derivatives. The numbers n , m and ℓ need not be integers. The generalization to a Hamiltonian with V_1 , V_2 depending also on the number of the particle only creates notational problems and is left to the reader.

An example of potentials that satisfy Assumptions 1–3 is

$$V_1(x) = x^4 - x^2 + 2 \quad \text{and} \quad V_2(x) = (1 + x^2)^{5/2} - \cos(x).$$

The effective energy of the system chain+baths is given by

$$G(p, q, r) = H_S(p, q) + \frac{r_L^2}{2\lambda_L^2} + \frac{r_R^2}{2\lambda_R^2} - q_0 r_L - q_N r_R + \Gamma, \quad (3.6)$$

where we choose the constant Γ such that $G \geq 1$, which is always possible, because $n > 1$. In fact, it is important that the function $\exp(-\beta G)$ be integrable for any $\beta > 0$. This could also be achieved with for example only one of the one-body potentials non-vanishing, but would cause some unimportant notational difficulties. The case $n = 1$ is marginal, the stability of the system depends on the values of the constants λ_i and was treated in [EPR99a]. We will not treat this case, but it would not cause any difficulties, as long as G remains confining.

In the sequel, we will extensively use the notations

$$\tilde{q}_i \equiv q_i - q_{i-1} \quad \text{and} \quad Q \equiv \sum_{i=0}^N q_i.$$

The system of stochastic differential equations we consider is given by

$$\begin{aligned}
dq_i &= p_i dt, \\
dp_0 &= -V_1'(q_0) dt + V_2'(\tilde{q}_1) dt + r_L dt, \\
dp_j &= -V_1'(q_j) dt - V_2'(\tilde{q}_j) dt + V_2'(\tilde{q}_{j+1}) dt, \\
dp_N &= -V_1'(q_N) dt - V_2'(\tilde{q}_N) dt + r_R dt, \\
dr_L &= -\gamma_L r_L dt + \lambda_L^2 \gamma_L q_0 dt - \lambda_L \sqrt{2\gamma_L T_L} dw_L(t), \\
dr_R &= -\gamma_R r_R dt + \lambda_R^2 \gamma_R q_N dt - \lambda_R \sqrt{2\gamma_R T_R} dw_R(t),
\end{aligned} \tag{3.7}$$

where $i = 1, \dots, N$ and $j = 1, \dots, N - 1$. Since Assumption 0 holds, the results of the preceding section apply. Therefore, there exists for any initial condition x_0 a unique stochastic process $\xi(t; x_0, w)$ solving (3.7). It obeys the statistics of a Markov diffusion process with generator

$$\begin{aligned}
\mathcal{L} &= \lambda_L^2 \gamma_L T_L \partial_{r_L}^2 + \lambda_R^2 \gamma_R T_R \partial_{r_R}^2 - \gamma_L (r_L - \lambda_L^2 q_0) \partial_{r_L} - \gamma_R (r_R - \lambda_R^2 q_N) \partial_{r_R} \\
&+ r_L \partial_{p_0} + r_R \partial_{p_N} + \sum_{i=0}^N (p_i \partial_{q_i} - V_1'(q_i) \partial_{p_i}) - \sum_{i=1}^N V_2'(\tilde{q}_i) (\partial_{p_i} - \partial_{p_{i-1}}).
\end{aligned} \tag{3.8}$$

We want to prove the existence of a smooth invariant measure with density $\mu(p, q, r)$. It is the solution of $(\mathcal{T}^t)^* \mu = 0$, where $(\mathcal{T}^t)^*$ is the dual semigroup of \mathcal{T}^t . To achieve this, we introduce, as above, the Hilbert space

$$\mathcal{H}_0 \equiv L^2(\mathbf{R}^{2N+4}, Z_0^{-1} e^{-2\beta_0 G(p,q,r)} dp dq dr),$$

where Z_0 is a normalization constant and β_0 is a ‘‘reference’’ inverse temperature that we choose such that

$$1/\beta_0 \equiv T_0 > \max\{T_L, T_R\}. \tag{3.9}$$

Proposition 2.4 holds, so the dynamics of our system is described by a semigroup \mathcal{T}_0^t acting in \mathcal{H}_0 with generator L_0 , formally given by \mathcal{L} . The extended phase space of our system will again be denoted by $\mathcal{X} \equiv \mathbf{R}^{2N+4}$.

For convenience, we would like to work in $\mathcal{H} = L^2(\mathcal{X})$, so we define the unitary transformation $U : \mathcal{H} \rightarrow \mathcal{H}_0$ by

$$(Uf)(x) = e^{\beta_0 G(x)} f(x).$$

So L_0 is unitarily equivalent to the operator $L_{\mathcal{H}} : \mathcal{D}(L_{\mathcal{H}}) \rightarrow \mathcal{H}$ defined by

$$L_{\mathcal{H}} = U^{-1} L_0 U = e^{-\beta_0 G} L_0 e^{\beta_0 G}.$$

An explicit computation shows that $L_{\mathcal{H}}$ is given by

$$L_{\mathcal{H}} = \alpha - K,$$

where the formal expression for the differential operator K is

$$\begin{aligned}
 K = & \alpha_K - c_L^2 \partial_{r_L}^2 + a_L^2 (r_L - \lambda_L^2 q_0)^2 - c_R^2 \partial_{r_R}^2 + a_R^2 (r_R - \lambda_R^2 q_N)^2 \\
 & - r_L \partial_{p_0} + b_L (r_L - \lambda_L^2 q_0) \partial_{r_L} - r_R \partial_{p_N} + b_R (r_R - \lambda_R^2 q_N) \partial_{r_R} \\
 & - \sum_{i=0}^N (p_i \partial_{q_i} - V_1'(q_i) \partial_{p_i}) + \sum_{i=1}^N V_2'(\tilde{q}_i) (\partial_{p_i} - \partial_{p_{i-1}}).
 \end{aligned} \tag{3.10}$$

Since $\mathcal{C}_0^\infty(\mathcal{X})$ is invariant under the unitary transformation U , it remains a core for both K and K^* . The various constants appearing in (3.10) are given by

$$\begin{aligned}
 a_i^2 &= \gamma_i (\beta_0 T_i - 1), \\
 b_i &= \frac{\gamma_i \beta_0}{\lambda_i^2} (\beta_0 T_i - 1), \quad i \in \{L, R\}, \\
 c_i &= \lambda_i \sqrt{\gamma_i T_i}, \\
 \alpha_K &= -\frac{b_L}{2} - \frac{b_R}{2}, \\
 \alpha &= \alpha_K + \beta_0 \sum_{i \in \{L, R\}} \gamma_i T_i.
 \end{aligned}$$

We see that condition (3.9) ensures the positivity of the constants a_L^2 and a_R^2 , which in turn implies that the closure of $\text{Re}K = (K + K^*)/2$ is a strictly positive self-adjoint operator.

The first feature we notice about K is that Assumption 3 implies the hypoellipticity of the operators K , K^* , $\partial_t + K$ and $\partial_t + K^*$. We recall that a differential operator L acting on functions in a finite-dimensional differentiable manifold \mathcal{M} is called hypoelliptic if

$$\text{sing supp } f = \text{sing supp } Lf, \quad \text{for all } f \in \mathcal{D}'(\mathcal{M}),$$

where $\mathcal{D}'(\mathcal{M})$ is the space of distributions on $\mathcal{C}_0^\infty(\mathcal{M})$. In particular, the eigenfunctions of a hypoelliptic operator are \mathcal{C}^∞ .

The hypoellipticity of the above operators is a consequence of a theorem by Hörmander [Hör67, Hör85]: given a second-order differential operator

$$L = \sum_{i=1}^n L_i^* L_i + L_0 + c,$$

where $c : \mathcal{M} \rightarrow \mathbf{C}$ is a smooth function and the L_i are smooth vector fields. Then a sufficient condition for L to be hypoelliptic is that the Lie algebra generated by $\{L_i \mid i = 0, \dots, n\}$ has maximal rank everywhere. It is not hard to verify that Assumption 3 ensures that this condition is verified for K , K^* , $\partial_t + K$ and $\partial_t + K^*$.

Proposition 3.5. *If Assumptions 0 and 3 are satisfied, the transition probabilities of the Markov process solving (3.7) have a smooth density*

$$P(t, x, y) \in \mathcal{C}^\infty((0, \infty) \times \mathcal{X} \times \mathcal{X}).$$

Proof. This is an immediate consequence of the Kolmogorov equations which state that

$$\partial_t P = \mathcal{L}P \quad \Rightarrow \quad (\partial_t + K - \alpha)U^{-1}P = 0,$$

so $U^{-1}P$ is an eigenfunction of the operator $\partial_t + K - \alpha$, which is hypoelliptic. \square

3.2. *Main results.* Our main technical result is

Theorem 3.6. *If Assumptions 1–3 are satisfied, then the operator K defined in (3.10) has compact resolvent.*

In order to prepare the proof of Theorem 3.6, we will prove the following two propositions.

Proposition 3.7. *If Assumptions 1 and 2 are satisfied, there exist constants C and $\varepsilon > 0$ such that*

$$\|G^\varepsilon f\| \leq C(\|Kf\| + \|f\|), \quad \text{for all } f \in \mathcal{D}(K), \quad (3.11a)$$

$$\|G^\varepsilon f\| \leq C(\|K^* f\| + \|f\|), \quad \text{for all } f \in \mathcal{D}(K^*). \quad (3.11b)$$

Proposition 3.8. *If Assumptions 1–3 are satisfied, there exist constants $C, \varepsilon > 0$, a positive function $a_0 : \mathcal{X} \rightarrow \mathbf{R}$ and a finite number \tilde{N} of smooth vector fields L_i with bounded coefficients such that, for every function $f \in C_0^\infty(\mathcal{X})$, we have*

$$\|\tilde{\Delta}^\varepsilon f\| \leq C(\|Kf\| + \|f\|), \quad (3.12)$$

where

$$\tilde{\Delta} = \sum_{i=1}^{\tilde{N}} L_i^* L_i + a_0.$$

Moreover, the L_i span the whole of \mathbf{R}^{2N+4} at every point.

Given Theorem 3.6, we can state and prove the main result of this paper, namely the existence and uniqueness of an invariant measure for our Markov process. More precisely, we have the following result.

Theorem 3.9. *If Assumptions 1–3 are satisfied, then the stochastic process $\xi(t)$ solving (2.2) possesses a unique and strictly positive invariant measure μ . Its density h is C^∞ and satisfies for any $\beta_0 < \min\{\beta_L, \beta_R\}$,*

$$h(x) = \tilde{h}(x)e^{-\beta_0 G(x)},$$

where \tilde{h} decays at infinity faster than any polynomial.

The above results say that the spectrum of K looks roughly like the one schematically depicted in Fig. 3.2. We see that it is discrete (compactness of the resolvent) and located in the right half of the complex plane (m -accreitivity). Moreover, it is symmetric along the real axis, because K is a differential operator with real coefficients.

Most of the remainder of this paper is devoted to the proofs of Theorems 3.6 and 3.9. In the sequel, we will always use the notation

$$K = \sum_{i=1}^4 X_i^* X_i + X_0,$$

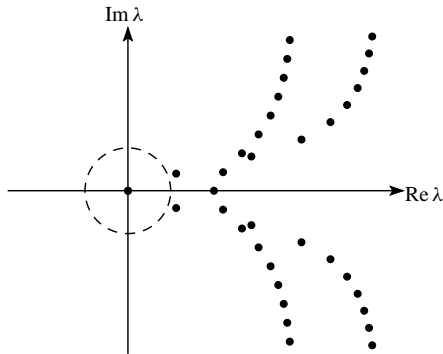


Fig. 3.2. Spectrum of K

where we define

$$X_1 = c_L \partial_{r_L}, \quad X_2 = a_L (r_L - \lambda_L^2 q_0), \tag{3.13a}$$

$$X_3 = c_R \partial_{r_R}, \quad X_4 = a_R (r_R - \lambda_R^2 q_N), \tag{3.13b}$$

$$\begin{aligned} X_0 = & -r_L \partial_{p_0} + b_L (r_L - \lambda_L^2 q_0) \partial_{r_L} - r_R \partial_{p_N} + b_R (r_R - \lambda_R^2 q_N) \partial_{r_R} \\ & - \sum_{i=0}^N (p_i \partial_{q_i} - V_1'(q_i) \partial_{p_i}) + \sum_{i=1}^N V_2'(\tilde{q}_i) (\partial_{p_i} - \partial_{p_{i-1}}) - \alpha_K. \end{aligned} \tag{3.13c}$$

The operator X_0 is antisymmetric, i.e.

$$X_0^* = -X_0. \tag{3.14}$$

This implies that

$$\text{Re} K = \sum_{i=1}^4 X_i^* X_i \quad \text{and} \quad X_0 = K - \text{Re} K, \tag{3.15}$$

and thus $\text{Re} K$ is a positive self-adjoint operator. We have one more estimate that will be extensively used in the sequel. If f is some function in $C_0^\infty(\mathcal{X})$ and $i \in \{1, \dots, 4\}$ we have

$$\begin{aligned} \|X_i f\|^2 &= \langle f, X_i^* X_i f \rangle \leq \langle f, \text{Re} K f \rangle = \text{Re} \langle f, K f \rangle \\ &\leq \|f\| \|K f\| \leq (\|K f\| + \|f\|)^2, \end{aligned} \tag{3.16}$$

and by a similar argument also

$$\|X_i^* f\|^2 \leq (\|K f\| + \|f\|)^2. \tag{3.17}$$

4. Proof of the Bound in Position Space (Proposition 3.7)

First of all, we need a collection of functions belonging to \mathcal{F}_0 , as defined in Definition 3.1. We have the following result.

Proposition 4.1. *Let r , p , q and \tilde{q} designate the vectors*

$$\begin{aligned} r &= (r_L, r_R), & q &= (q_0, \dots, q_N), \\ p &= (p_0, \dots, p_N), & \tilde{q} &= (\tilde{q}_1, \dots, \tilde{q}_N). \end{aligned}$$

Choose $\alpha \geq 0$ and let $h_k : \mathbf{R}^k \rightarrow \mathbf{R}$ be functions in \mathcal{F}_α . Then the functions

$$G^{-\alpha/2}h_2(r), \quad G^{-\alpha/2}h_{N+1}(p), \quad G^{-\alpha/(2n)}h_{N+1}(q), \quad \text{and} \quad G^{-\alpha/(2m)}h_N(\tilde{q})$$

belong to \mathcal{F}_0 .

Proof. We will only sketch the proof of the statement for $G^{-\alpha/(2m)}h_N(\tilde{q})$. The other expressions can easily be treated in a similar way.

We first notice that $G^{-1}(D^k G)$ is bounded for every multi-index k . This is a straightforward consequence of two observations. The first one is that because of the lower bounds (3.4a) and (3.5a) of Assumptions 1 and 2 and the expression (3.6) of G , there exists a constant $C > 0$ for which

$$G(p, q, r) \geq C(r^2 + p^2 + P^{2n}(q) + P^{2m}(\tilde{q})), \quad (4.1)$$

where P^k was defined in (3.1). The second observation is that, because $V_1 \in \mathcal{F}_{2n}$ and $V_2 \in \mathcal{F}_{2m}$, we have for every multi-index k some constant C_k for which

$$|D^k G(p, q, r)| \leq C_k(r^2 + p^2 + P^{2n}(q) + P^{2m}(\tilde{q})). \quad (4.2)$$

Notice that $G^{-\alpha/(2m)}D^k h_N(\tilde{q})$ is bounded by a similar argument, in particular because $h_N \in \mathcal{F}_\alpha$.

We set $\alpha = -\alpha/(2m)$ and write

$$\partial_i(G^\alpha h_N(\tilde{q})) = \alpha(G^{-1}\partial_i G)G^\alpha h_N(\tilde{q}) + G^\alpha \partial_i h_N(\tilde{q}).$$

Both terms are bounded by (4.1), (4.2) and the fact that $h_N \in \mathcal{F}_\alpha$. It is easy to see that all the derivatives can be bounded similarly. The proof of Proposition 4.1 is complete. \square

Let us define

$$\Lambda_1 \equiv G^{1/2}.$$

The symbol Λ_1 was chosen in order to emphasize the similarity between the proof of Proposition 3.7 and the proof of the main result of Sect. 5, Theorem 5.5.

Before we start the proof of Proposition 3.7, we notice two more facts. Let us choose $\alpha, \beta \in \mathbf{R}$ with $0 \leq \beta \leq 1$, and let A, B be two operators of multiplication by positive functions $A \leq B$. We then have

$$\langle \Lambda_1^\alpha A f, f \rangle \leq \langle \Lambda_1^\alpha B f, f \rangle, \quad (4.3)$$

as well as the implication

$$\|\Lambda_1^\alpha A f\| \leq C(\|K f\| + \|f\|) \quad \Rightarrow \quad \|\Lambda_1^{\alpha\beta} A^\beta f\| \leq C(\|K f\| + \|f\|). \quad (4.4)$$

Both inequalities are trivial consequences of the fact that Λ_1 is an operator of multiplication by a positive function and the estimate $x^s \leq 1 + x$ if $x \geq 0$ and $s \leq 1$.

4.1. *The main tool of the proof.* The main tool in the proof of Proposition 3.7 is the following lemma.

Lemma 4.2. *Let Λ_1 and K be defined as above. Let A and B be multiplication operators represented by functions of the form*

$$h(p, q, r) = c_L r_L + c_R r_R + \tilde{h}(p, q), \quad \tilde{h} \in C^\infty(\mathbf{R}^{2N+2}).$$

Assume moreover that there are exponents α_i and β_i and positive constants C_i such that the following estimates are true for every $f \in C_0^\infty(\mathcal{X})$:

$$\begin{aligned} \|\Lambda_1^{-\alpha_1} A f\| &\leq C_1(\|K f\| + \|f\|), & \|\Lambda_1^{-\beta_1} B f\| &\leq C_2(\|K f\| + \|f\|), \\ \|\Lambda_1^{-\alpha_2} A f\| &\leq C_3\|f\|, & \|\Lambda_1^{-\beta_2} B f\| &\leq C_4\|f\|, \\ \|\Lambda_1^{-\alpha_3} [X_0, A] f\| &\leq C_5(\|K f\| + \|f\|), & \|\Lambda_1^{-\beta_3} [X_0, B] f\| &\leq C_6(\|K f\| + \|f\|). \end{aligned}$$

If γ satisfies the conditions

$$\gamma \geq \alpha_3 + \beta_1, \tag{4.5}$$

$$\gamma \geq \alpha_2 + \frac{\beta_1 + \max\{\beta_2, \beta_3\}}{2}, \tag{4.6}$$

$$\gamma \geq \min\{\alpha_1 + \beta_2, \alpha_2 + \beta_1\}, \tag{4.7}$$

then there exists a constant C such that

$$|\langle [X_0, B] f, \Lambda_1^{-\gamma} A f \rangle| \leq C(\|K f\| + \|f\|)^2, \quad \text{for all } f \in C_0^\infty(\mathcal{X}). \tag{4.8}$$

Proof. The proof of this lemma involves some of the commutation techniques developed by Hörmander [Hör85], but it uses the fact that most operators involved are multiplication operators, i.e. they commute. An explicit computation, using (3.6) and (3.13) yields

$$[X_0, G] = \sum_{j \in \{R, L\}} \frac{b_j}{\lambda_j^2} (r_j - \lambda_j^2 F_j)^2, \quad [X_1, G] = c_L (r_L / \lambda_L^2 - q_0), \tag{4.9a}$$

$$[X_2, G] = [X_4, G] = 0, \quad [X_3, G] = c_R (r_R / \lambda_R^2 - q_N). \tag{4.9b}$$

We therefore see that, by Proposition 4.1, we have for $i = 0, \dots, 4$,

$$G^{-1} [X_i, G] \in \mathcal{F}_0. \tag{4.10}$$

Since the X_i are either differentiation operators or multiplicative operators, we have, for any $\alpha \in \mathbf{R}$, the relation

$$G^{-\alpha} [X_i, G^\alpha] = \alpha G^{-1} [X_i, G] \in \mathcal{F}_0,$$

and so, since $\Lambda_1^2 = G$,

$$\|\Lambda_1^\alpha [X_i, \Lambda_1^{-\alpha}]\| < \infty. \tag{4.11}$$

We can now start to bound (4.8). Since $[X_0, B] = -X_0^* B - B X_0$, we can write (4.8) as

$$\begin{aligned} |\langle [X_0, B] f, \Lambda_1^{-\gamma} A f \rangle| &\leq |\langle B X_0 f, \Lambda_1^{-\gamma} A f \rangle| + |\langle B f, X_0 \Lambda_1^{-\gamma} A f \rangle| \\ &\equiv T_1 + T_2. \end{aligned}$$

Both terms will be estimated separately.

Term T_1 . Since we know by (3.15) that $X_0 = K - \operatorname{Re}K$, we can write it as

$$T_1 \leq |\langle B(\operatorname{Re}K)f, \Lambda_1^{-\gamma}Af \rangle| + |\langle BKf, \Lambda_1^{-\gamma}Af \rangle| \equiv T_{11} + T_{12}.$$

The term T_{12} can be estimated by using (4.7). We indeed have either $\gamma \geq \alpha_1 + \beta_2$, or $\gamma \geq \alpha_2 + \beta_1$. In the former case, we write

$$T_{12} \leq \|\Lambda_1^{-\beta_2}B\| \|Kf\| \|\Lambda_1^{-\gamma+\beta_2}Af\| \leq C(\|Kf\| + \|f\|)^2.$$

In the latter case, we use the fact that A , B and Λ_1 commute and are self-adjoint to write similarly

$$T_{12} = |\langle AKf, \Lambda_1^{-\gamma}Bf \rangle| \leq \|\Lambda_1^{-\alpha_2}A\| \|Kf\| \|\Lambda_1^{-\gamma+\alpha_2}Bf\| \leq C(\|Kf\| + \|f\|)^2.$$

Let us now focus on the term T_{11} . Using the positivity of $\operatorname{Re}K$, it can be written as

$$\begin{aligned} T_{11} &= \langle (\operatorname{Re}K)^{1/2} \Lambda_1^{-\gamma_1} Bf, (\operatorname{Re}K)^{1/2} \Lambda_1^{-\gamma_2} Af \rangle + \langle [\Lambda_1^{-\gamma_1} B, \operatorname{Re}K]f, \Lambda_1^{-\gamma_2} Af \rangle \\ &\equiv T_{13} + T_{14}, \end{aligned}$$

where

$$\gamma_1, \gamma_2 > 0, \quad \gamma_1 + \gamma_2 = \gamma,$$

are to be chosen later. We estimate both terms separately. The commutator in T_{14} can be expanded to give

$$T_{14} = \langle \Lambda_1^{-\gamma_1} [B, \operatorname{Re}K]f, \Lambda_1^{-\gamma_2} Af \rangle + \langle [\Lambda_1^{-\gamma_1}, \operatorname{Re}K]Bf, \Lambda_1^{-\gamma_2} Af \rangle.$$

In order to estimate these terms, we recall that $\operatorname{Re}K = \sum_{i=1}^4 X_i^* X_i$. We therefore have

$$\begin{aligned} T_{14} &= \sum_{i=1}^4 \left(\langle \Lambda_1^{-\gamma_1} [B, X_i^*] X_i f, \Lambda_1^{-\gamma_2} Af \rangle + \langle \Lambda_1^{-\gamma_1} X_i^* [B, X_i] f, \Lambda_1^{-\gamma_2} Af \rangle \right. \\ &\quad \left. + \langle [\Lambda_1^{-\gamma_1}, X_i^*] X_i Bf, \Lambda_1^{-\gamma_2} Af \rangle + \langle X_i^* [\Lambda_1^{-\gamma_1}, X_i] Bf, \Lambda_1^{-\gamma_2} Af \rangle \right) \\ &\equiv \sum_{i=1}^4 (T_i^{(1)} + T_i^{(2)} + T_i^{(3)} + T_i^{(4)}). \end{aligned}$$

Noticing that $[B, X_i^*]$ is a multiple of the identity operator and that Λ_1 is self-adjoint, we have

$$|T_i^{(1)}| \leq C \langle X_i f, \Lambda_1^{-\gamma} Af \rangle \leq \|X_i f\| \|\Lambda_1^{-\gamma} Af\| \leq C(\|Kf\| + \|f\|)^2,$$

where we used (3.16) and the fact that $\gamma > \alpha_2$ to get the last inequality. The term $T_i^{(2)}$ is bounded by $C(\|Kf\| + \|f\|)^2$ in a similar way. The term $T_i^{(3)}$ is written as

$$\begin{aligned} |T_i^{(3)}| &= |\langle \Lambda_1^{\gamma_1} [\Lambda_1^{-\gamma_1}, X_i^*] X_i f, \Lambda_1^{-\gamma} ABf \rangle + \langle \Lambda_1^{\gamma_1} [\Lambda_1^{-\gamma_1}, X_i^*] [X_i, B] f, \Lambda_1^{-\gamma} Af \rangle| \\ &\leq C \|X_i f\| \|\Lambda_1^{-\gamma} ABf\| + C \|f\| \|\Lambda_1^{-\gamma} Af\|, \end{aligned}$$

where we used (4.11) and the fact that $[X_i, B]$ is bounded. Now we can bound $T_i^{(3)}$ by $C(\|Kf\| + \|f\|)^2$, using (3.16) to estimate $\|X_i f\|$ and (4.7) to estimate $\|\Lambda_1^{-\gamma} A B f\|$ and $\|\Lambda_1^{-\gamma} A f\|$. The term $T_i^{(4)}$ can be estimated in a similar way.

Let us now focus on the term T_{13} . We can write

$$|T_{13}| \leq |\operatorname{Re}\langle K \Lambda_1^{-\gamma_1} B f, \Lambda_1^{-\gamma_1} B f \rangle|^{1/2} \sqrt{\sum_{i=1}^4 \|X_i \Lambda_1^{-\gamma_2} A f\|}.$$

If we choose

$$\gamma_2 = \alpha_2, \quad (4.12)$$

the terms under the square root are easily estimated by writing them as

$$\begin{aligned} \|X_i \Lambda_1^{-\gamma_2} A f\| &\leq \|\Lambda_1^{-\gamma_2} A\| \|X_i f\| + \|[X_i, \Lambda_1^{-\gamma_2}] \Lambda_1^{\gamma_2}\| \|\Lambda_1^{-\gamma_2} A f\| \\ &\quad + \|\Lambda_1^{-\gamma_2} [X_i, A] f\|, \end{aligned}$$

and estimating the two commutators by (4.11) and (4.9) respectively.

The term preceding the square root can be written as

$$\begin{aligned} \langle K \Lambda_1^{-\gamma_1} B f, \Lambda_1^{-\gamma_1} B f \rangle &= \langle \Lambda_1^{-\gamma_1} B K f, \Lambda_1^{-\gamma_1} B f \rangle + \langle [K, \Lambda_1^{-\gamma_1} B] f, \Lambda_1^{-\gamma_1} B f \rangle \\ &\equiv T_{15} + T_{16}. \end{aligned}$$

The term T_{15} can be bounded if we choose

$$2\gamma_1 \geq \beta_1 + \beta_2, \quad (4.13)$$

because we have then

$$T_{15} \leq \|K f\| \|\Lambda_1^{-\beta_2} B\| \|\Lambda_1^{-\beta_1} B f\| \leq C(\|K f\| + \|f\|)^2.$$

In order to estimate the term T_{16} , we use $K = \operatorname{Re}K + X_0$ to write

$$\begin{aligned} T_{16} &= \langle [X_0, \Lambda_1^{-\gamma_1} B] f, \Lambda_1^{-\gamma_1} B f \rangle + \langle [\operatorname{Re}K, \Lambda_1^{-\gamma_1} B] f, \Lambda_1^{-\gamma_1} B f \rangle \\ &\equiv T_{16}^{(1)} + T_{16}^{(2)}. \end{aligned}$$

The term $T_{16}^{(1)}$ can be estimated by writing it as

$$T_{16}^{(1)} = \langle \Lambda_1^{-\gamma_1} [X_0, B] f, \Lambda_1^{-\gamma_1} B f \rangle + \langle [X_0, \Lambda_1^{-\gamma_1}] \Lambda_1^{\gamma_1} \Lambda_1^{-\gamma_1} B f, \Lambda_1^{-\gamma_1} B f \rangle.$$

The first term can be bounded by $C(\|K f\| + \|f\|)^2$ if we choose

$$2\gamma_1 \geq \beta_1 + \beta_3. \quad (4.14)$$

In order to bound the second term, it suffices to have $\gamma_1 \geq \beta_1$, which is the case because of (4.13) and the fact that $\beta_2 \geq \beta_1$.

The term $T_{16}^{(2)}$ can be bounded by $C(\|K f\| + \|f\|)^2$, by treating it in a similar way as the term T_{14} . We leave to the reader the verification that no additional conditions on γ_1 have to be made. This completes the estimate of T_1 , because (4.12), (4.13) and (4.14) can be satisfied simultaneously by (4.6).

Term T_2 . We decompose this term as

$$\begin{aligned} T_2 &\leq |\langle Bf, \Lambda_1^{-\gamma} A X_0 f \rangle| + |\langle Bf, \Lambda_1^{-\gamma} [X_0, A] f \rangle| + |\langle Bf, [X_0, \Lambda_1^{-\gamma}] A f \rangle| \\ &\equiv T_{21} + T_{22} + T_{23}. \end{aligned}$$

Since $\gamma \geq \alpha_3 + \beta_1$ the term T_{22} is easily estimated by

$$T_{22} \leq \|\Lambda_1^{-\beta_1} Bf\| \|\Lambda_1^{-\alpha_3} [X_0, A] f\| \leq C(\|Kf\| + \|f\|)^2.$$

Noticing that we can assume $\alpha_1 \leq \alpha_2$ and $\beta_1 \leq \beta_2$, condition (4.7) implies $\gamma \geq \alpha_1 + \beta_1$. Since $[X_0, \Lambda_1^{-\gamma}]$ is a function, it commutes with Λ_1 , and so T_{23} can be estimated writing

$$\begin{aligned} T_{23} &\leq |\langle \Lambda_1^{-\beta_1} Bf, \Lambda_1^\gamma [X_0, \Lambda_1^{-\gamma}] \Lambda_1^{-\alpha_1} A f \rangle| \\ &\leq \|\Lambda_1^{-\beta_1} Bf\| \| [X_0, \Lambda_1^{-\gamma}] \Lambda_1^\gamma \| \| \Lambda_1^{-\alpha_1} A f \| \leq C(\|Kf\| + \|f\|)^2, \end{aligned}$$

where we used (4.11) to get the last bound.

We finally bound T_{21} . Since $X_0 = K - \operatorname{Re}K$, it can be expanded as

$$T_{21} \leq |\langle Bf, \Lambda_1^{-\gamma} A K f \rangle| + |\langle Bf, \Lambda_1^{-\gamma} A (\operatorname{Re}K) f \rangle| \equiv T_{21}^{(1)} + T_{21}^{(2)}.$$

The term $T_{21}^{(1)}$ can be estimated by writing

$$T_{21}^{(1)} \leq \|Kf\| \|\Lambda_1^{-\gamma} A Bf\|,$$

and using (4.7). The term $T_{21}^{(2)}$ can be written as

$$T_{21}^{(2)} = \langle Bf, \Lambda_1^{-\gamma} A (\operatorname{Re}K) f \rangle = T_{13} + \langle \Lambda_1^{-\gamma_1} Bf, [\Lambda_1^{-\gamma_2} A, \operatorname{Re}K] f \rangle.$$

The term T_{13} has already been estimated. The other term can be treated like the term T_{14} . We leave to the reader the verification that one can indeed bound it by $C(\|Kf\| + \|f\|)^2$ without any further restriction on γ_1 and γ_2 .

This completes the proof of the lemma. \square

4.2. The main step of the proof of Proposition 3.7. By an elementary approximation argument, it is sufficient to prove the inequalities (3.11) for $f \in C_0^\infty(\mathcal{X})$, since this is a core for both K and K^* . Moreover, we will prove only (3.11a). The interested reader may verify that the same arguments also apply for (3.11b).

We want to show that we can find constants ε and C such that

$$\|\Lambda_1^\varepsilon f\| \leq C(\|Kf\| + \|f\|), \quad \text{for all } f \in C_0^\infty(\mathcal{X}).$$

In order to show this, we notice that there is a constant C such that

$$\begin{aligned} \Lambda_1^2 &\leq C \left(1 + (r_L - \lambda_L^2 q_0)^2 + (r_R - \lambda_r^2 q_N)^2 + \sum_{i=0}^N p_i^2 + P^{2n}(Q) + \sum_{i=1}^N P^{2m}(\tilde{q}_i) \right) \\ &\equiv \tilde{G}. \end{aligned}$$

The immediate consequence is that

$$\|\Lambda_1^\varepsilon f\|^2 = \langle f, \Lambda_1^{2\varepsilon} f \rangle \leq \langle f, \Lambda_1^{2\varepsilon-2} \tilde{G} f \rangle.$$

It is therefore enough to show that there exists a (small) constant ε such that the terms

$$\|\Lambda_1^{\varepsilon-1} P^n(Q)f\|, \|\Lambda_1^{\varepsilon-1} p_i f\|, \|\Lambda_1^{\varepsilon-1} P^m(\tilde{q}_i)f\|, \dots$$

are bounded by $C(\|Kf\| + \|f\|)$.

We are first going to bound the terms involving variables near the boundary of the chain. Then, we will proceed by induction towards the middle of the chain.

The term $\|\Lambda_1^{\varepsilon-1}(r_L - \lambda_L^2 q_0)f\|$. We have

$$\begin{aligned} \|(r_L - \lambda_L^2 q_0)f\|^2 &= |\langle (r_L - \lambda_L^2 q_0)^2 f, f \rangle| \leq C|\langle (\operatorname{Re}K)f, f \rangle| \\ &= C|\operatorname{Re}\langle Kf, f \rangle| \leq C\|Kf\|\|f\| \leq C(\|Kf\| + \|f\|)^2, \end{aligned} \quad (4.15)$$

where we used the fact that $a_L \neq 0$ to obtain the first inequality. Since $\Lambda_1 \geq 1$, we thus have the estimate

$$\|\Lambda_1^{\varepsilon-1}(r_L - \lambda_L^2 q_0)f\| \leq C(\|Kf\| + \|f\|)^2$$

if we take $\varepsilon \leq 1$.

The term $\|\Lambda_1^{\varepsilon-1} p_0 f\|$. We will prove the estimate

$$\|\Lambda_1^{\varepsilon_0-1} p_0 f\| \leq C(\|Kf\| + \|f\|), \quad (4.16)$$

for $\varepsilon_0 \leq 1/(2m)$. An explicit computation yields the relation

$$[X_0, r_L - \lambda_L^2 q_0] = b_L(r_L - \lambda_L^2 q_0) - \lambda_L^2 p_0. \quad (4.17)$$

Solving (4.17) for p_0 , we get

$$\begin{aligned} \|\Lambda_1^{\varepsilon_0-1} p_0 f\|^2 &= \langle \lambda_L^{-2}(b_L(r_L - \lambda_L^2 q_0) - \lambda_L^{-2}[X_0, r_L - \lambda_L^2 q_0])f, \Lambda_1^{2\varepsilon_0-2} p_0 f \rangle \\ &\equiv X_0^{(1)} - X_0^{(2)}. \end{aligned}$$

The term $X_0^{(1)}$ can be estimated as

$$|X_0^{(1)}| \leq \lambda_L^{-2} \|b_L(r_L - \lambda_L^2 q_0)f\| \|\Lambda_1^{2\varepsilon_0-2} p_0 f\| \leq C(\|Kf\| + \|f\|)^2,$$

where the last inequality holds because $\varepsilon_0 \leq 1/2$.

In order to estimate $X_0^{(2)}$, we apply Lemma 4.2 with $A = p_0$ and $B = r_L - \lambda_L^2 q_0$. An explicit computation yields $[X_0, A] = V_1'(q_0) + V_2'(\tilde{q}_1) - r_L$. The term $[X_0, B]$ has already been computed in (4.17). Because of Proposition 4.1 and of (4.15), we can choose

$$\begin{aligned} \alpha_1 &= 1, & \beta_1 &= 0, \\ \alpha_2 &= 1, & \beta_2 &= 1, \\ \alpha_3 &= 2 - 1/m, & \beta_3 &= 1. \end{aligned}$$

The hypotheses of Lemma 4.2 are thus fulfilled if we choose $\gamma = 2 - 1/m$. We therefore have the estimate (4.16) with $\varepsilon_0 = 1/(2m)$. We have a similar estimate for the symmetric term at the other end of the chain.

The term $\|\Lambda_1^{\varepsilon_0^{-1}} P^m(\tilde{q}_1) f\|$. We will prove the estimate

$$\|\Lambda_1^{\varepsilon_0^{-1}} P^m(\tilde{q}_1) f\| \leq C(\|Kf\| + \|f\|),$$

for some $\varepsilon'_0 < \varepsilon_0$. Because of the bound (3.5b) of Assumption 2, we can find some constants c_1 and c_2 such that

$$\langle \Lambda_1^{2\varepsilon_0^{-2}} P^{2m}(\tilde{q}_1) f, f \rangle \leq c_1 |\langle \Lambda_1^{2\varepsilon_0^{-2}} V_2'(\tilde{q}_1) f, \tilde{q}_1 f \rangle| + c_2 |\langle \Lambda_1^{2\varepsilon_0^{-2}} f, f \rangle|, \quad (4.18)$$

where we also used (4.3). The second term is easily estimated because $\Lambda_1^{2\varepsilon_0^{-2}}$ is bounded if $\varepsilon'_0 \leq 1$. We once again use the fact that $[X_0, p_0] = V_1'(q_0) + V_2'(\tilde{q}_1) - r_L$ to write the first term as

$$\begin{aligned} |\langle \Lambda_1^{2\varepsilon_0^{-2}} V_2'(\tilde{q}_1) f, \tilde{q}_1 f \rangle| &= |\langle \Lambda_1^{2\varepsilon_0^{-2}} ([X_0, p_0] - V_1'(q_0) + r_L) f, \tilde{q}_1 f \rangle| \\ &\equiv |Y_1^{(1)} + Y_1^{(2)} + Y_1^{(3)}|. \end{aligned}$$

The term $Y_1^{(2)}$ can be written as

$$\begin{aligned} |Y_1^{(2)}| &= |\langle \Lambda_1^{2\varepsilon_0^{-2}+1/m} V_1'(q_0) f, \Lambda_1^{-1/m} \tilde{q}_1 f \rangle| \\ &\leq \|\Lambda_1^{2\varepsilon_0^{-2}+1/m} V_1'(q_0) f\| \|\Lambda_1^{-1/m} \tilde{q}_1 f\|. \end{aligned}$$

By Proposition 4.1 and the fact that $V_1' \in \mathcal{F}_{2n-1}$, this term is bounded by $C\|f\|^2$ if we take ε'_0 so small that

$$2\varepsilon'_0 \leq 1/n - 1/m. \quad (4.19)$$

The term $Y_1^{(3)}$ is bounded similarly by writing

$$|Y_1^{(3)}| \leq \|\Lambda_1^{2\varepsilon_0^{-2}+1/m} r_L f\| \|\Lambda_1^{-1/m} \tilde{q}_1 f\|,$$

if we impose

$$2\varepsilon'_0 \leq 1 - 1/m. \quad (4.20)$$

Both conditions can be satisfied because we assumed that $1 < n < m$. In order to estimate $Y_1^{(1)}$, we apply once again Lemma 4.2. This time we have $A = \tilde{q}_1$ and $B = p_0$. Using (4.16) and Proposition 4.1, we see that we can choose

$$\begin{aligned} \alpha_1 &= 1/m, & \beta_1 &= 1 - \varepsilon_0, \\ \alpha_2 &= 1/m, & \beta_2 &= 1, \\ \alpha_3 &= 1, & \beta_3 &= 2 - 1/m. \end{aligned}$$

By using $m > 1$, we see that the hypotheses of Lemma 4.2 are fulfilled if (4.19) and (4.20) hold, together with $\varepsilon'_0 < \varepsilon_0/2$. Once again, we have the same estimate at the other end of the chain.

We can now go along the chain by induction. At each step, we go one particle closer towards the middle of the chain. We present here only the terms arising when we go from the left to the right of the chain.

The term $\|\Lambda_1^{\varepsilon-1} p_i f\|$. We already treated the case $i = 1$. Let us therefore assume $i > 1$. We moreover assume that there exist constants $\varepsilon_{i-1}, \varepsilon'_{i-1} > 0$ such that we have the estimates

$$\begin{aligned} \|\Lambda_1^{\varepsilon_{i-1}-1} p_{i-1} f\| &\leq C(\|Kf\| + \|f\|), \\ \|\Lambda_1^{\varepsilon'_{i-1}-1} P^m(\tilde{q}_i) f\| &\leq C(\|Kf\| + \|f\|). \end{aligned} \quad (4.21)$$

We will show that this implies the existence of a constant $\varepsilon_i > 0$ such that

$$\|\Lambda_1^{\varepsilon_i-1} p_i f\| \leq C(\|Kf\| + \|f\|). \quad (4.22)$$

We use $p_i = p_{i-1} + [X_0, \tilde{q}_i]$ to write

$$\|\Lambda_1^{\varepsilon_i-1} p_i f\|^2 = \langle \Lambda_1^{2\varepsilon_i-1} p_{i-1} f, \Lambda_1^{-1} p_i f \rangle + \langle [X_0, \tilde{q}_i] f, \Lambda_1^{2\varepsilon_i-2} p_i f \rangle \equiv X_i^{(1)} + X_i^{(2)}.$$

The term $X_i^{(1)}$ is easily bounded if we write

$$|X_i^{(1)}| \leq \|\Lambda_1^{2\varepsilon_i-1} p_{i-1} f\| \|\Lambda_1^{-1} p_i f\| \leq C(\|Kf\| + \|f\|)^2,$$

where the last inequality is obtained by using Proposition 4.1 and (4.21). We only have to make the assumption $2\varepsilon_i \leq \varepsilon_{i-1}$.

In order to estimate the term $X_i^{(2)}$, we apply Lemma 4.2 with $A = p_i$ and $B = \tilde{q}_i$. Explicit computation yields $[X_0, p_i] = V'_1(q_i) - V'_2(\tilde{q}_{i+1}) - V'_2(\tilde{q}_i)$. Using the induction hypothesis (4.21) and Proposition 4.1, we see that we can choose

$$\begin{aligned} \alpha_1 &= 1, & \beta_1 &= (1 - \varepsilon'_{i-1})/m, \\ \alpha_2 &= 1, & \beta_2 &= 1/m, \\ \alpha_3 &= 2 - 1/m, & \beta_3 &= 1. \end{aligned}$$

If we take $\varepsilon_i \leq \varepsilon'_{i-1}/(2m)$, we see that the hypotheses of Lemma 4.2 are satisfied. We thus have the desired bound (4.22).

The term $\|\Lambda_1^{\varepsilon-1} P^m(\tilde{q}_{i+1}) f\|$. We assume that there exist strictly positive constants ε_i and ε'_{i-1} such that

$$\|\Lambda_1^{\varepsilon_i-1} p_i f\| \leq C(\|Kf\| + \|f\|), \quad \|\Lambda_1^{\varepsilon'_{i-1}-1} P^m(\tilde{q}_i) f\| \leq C(\|Kf\| + \|f\|).$$

We will show that this implies the existence of a constant $\varepsilon'_i > 0$ for which

$$\|\Lambda_1^{\varepsilon'_i-1} P^m(\tilde{q}_{i+1}) f\| \leq C(\|Kf\| + \|f\|). \quad (4.23)$$

Expression (4.18) with \tilde{q}_1 replaced by \tilde{q}_{i+1} holds. In order to prove (4.23), it suffices therefore to show that

$$|\langle \Lambda_1^{2\varepsilon'_i-2} V'_2(\tilde{q}_{i+1}) f, \tilde{q}_{i+1} f \rangle| \leq C(\|Kf\| + \|f\|)^2.$$

Since, for $i > 1$ we have $[X_0, p_i] = V'_1(q_i) - V'_2(\tilde{q}_{i+1}) - V'_2(\tilde{q}_i)$, the preceding term can be written as

$$|\langle \Lambda_1^{2\varepsilon'_i-2} ([X_0, p_i] + V'_1(q_i) + V'_2(\tilde{q}_i)) f, \tilde{q}_{i+1} f \rangle| \equiv |Y_i^{(1)} + Y_i^{(2)} + Y_i^{(3)}|.$$

We impose $2\varepsilon'_i \leq 1/n - 1/m$. The term $Y_i^{(2)}$ is then estimated as

$$|Y_i^{(2)}| \leq \|\Lambda_1^{-1/m} \tilde{q}_{i+1} f\| \|\Lambda_1^{2\varepsilon'_i - 2 + 1/m} V'_1(q_i) f\| \leq C(\|Kf\| + \|f\|)^2,$$

where the last step uses Proposition 4.1 and $V'_1 \in \mathcal{F}_{2n-1}$. In order to estimate the term $Y_i^{(3)}$, we notice that by the Cauchy-Schwarz inequality and Assumption 2, we have

$$\begin{aligned} |Y_i^{(3)}| &\leq C \|\Lambda_1^{-1/m} \tilde{q}_{i+1} f\| \|\Lambda_1^{2\varepsilon'_i - 2 + 1/m} P^{2m-1}(\tilde{q}_i) f\| \\ &\leq C \|f\| \|\Lambda_1^{1/m-1} P^{m-1}(\tilde{q}_i) \Lambda_1^{2\varepsilon'_i-1} P^m(\tilde{q}_i) f\| \\ &\leq C \|f\| \|\Lambda_1^{2\varepsilon'_i-1} P^m(\tilde{q}_i) f\|. \end{aligned}$$

We can choose $2\varepsilon'_i < \varepsilon'_{i-1}$, so this term can be estimated by the induction hypothesis.

The term $Y_i^{(1)}$ is once again estimated by using Lemma 4.2, this time with $A = \tilde{q}_{i+1}$ and $B = p_i$. Using Proposition 4.1, it is easy to verify that one can take

$$\begin{aligned} \alpha_1 &= 1/m, & \beta_1 &= 1 - \varepsilon_i, \\ \alpha_2 &= 1/m, & \beta_2 &= 1, \\ \alpha_3 &= 1, & \beta_3 &= 2 - 1/m. \end{aligned}$$

It suffices then to choose $2\varepsilon'_i < \varepsilon_i$ to satisfy the assumptions of Lemma 4.2 and get the desired estimate.

It is obvious that this induction also works in the other direction, starting from the other end of the chain. It also accommodates to a little bit more complicated topologies, as long as the chain does not contain any closed loop. In order to complete the proof of the lemma, we have to estimate the last term corresponding to the motion of the center of mass.

The term $\|\Lambda_1^{\varepsilon-1} P^n(Q) f\|$. Finally, we want to show the estimate

$$\|\Lambda_1^{\varepsilon-1} P^n(Q) f\| \leq C(\|Kf\| + \|f\|), \quad (4.24)$$

for some ε . We start with a little computation. We write

$$(N+1)q_0 = Q + (q_{N-1} - q_N) + 2(q_{N-2} - q_{N-1}) + \dots + N(q_0 - q_1).$$

Moreover, we have $q_i = q_0 + (q_1 - q_0) + \dots + (q_i - q_{i-1})$. We can thus write

$$\frac{Q}{N+1} - q_i = \sum_{j=1}^N b_{ij} \tilde{q}_j, \quad \text{with } b_{ij} \in \mathbf{R}.$$

This, together with the mean-value theorem, implies the useful relation

$$\begin{aligned} (N+1)QV'_1(Q/(N+1)) &= Q \sum_{i=0}^N V'_1(q_i) + Q \sum_{i=0}^N \left(V'_1(Q/(N+1)) - V'_1(q_i) \right) \\ &= Q \sum_{i=0}^N V'_1(q_i) + Q \sum_{i=0}^N V''_1(\xi_i) \sum_{j=1}^N b_{ij} \tilde{q}_j, \end{aligned} \quad (4.25)$$

where ξ_i is located somewhere on the $Q/(N+1)$ and q_i .

In the case of d -dimensional particles, the expression corresponding to (4.25) is

$$\begin{aligned} & |(N+1)QV_1'(Q/(N+1))| \\ & \leq (N+1)|Q|\|\nabla V_1(q_i)\| + |Q| \sum_{i=0}^N \sup_{t \in (0,1)} \left| \nabla^2 V_1 \left(\frac{tQ}{N+1} + (1-t)q_i \right) \right| \left| \sum_{j=1}^N b_{ij} |\tilde{q}_j| \right|. \end{aligned}$$

The subsequent expressions can be rewritten accordingly.

We use Assumption 1 and (4.25) to write the left-hand side of (4.24) as

$$\begin{aligned} \|\Lambda_1^{\varepsilon-1} P^n(Q)f\|^2 & = |\langle \Lambda_1^{2\varepsilon-2} P^{2n}(Q)f, f \rangle| \\ & \leq C(N+1) \left| \langle \Lambda_1^{2\varepsilon-2} V_1'(Q/(N+1))f, Qf \rangle \right| + C\|f\|^2 \\ & \leq C \left| \langle \Lambda_1^{2\varepsilon-2} \left(\sum_{i=0}^N V_1'(q_i) \right) f, Qf \rangle \right| \\ & \quad + C \sum_{i,j=1}^N b_{ij} |\langle \Lambda_1^{2\varepsilon-2} \tilde{q}_j V_1''(\xi_i) f, Qf \rangle| + C\|f\|^2 \\ & \equiv Y^{(1)} + Y^{(2)} + C\|f\|^2. \end{aligned}$$

The term $Y^{(2)}$ can be bounded because $V_1'' \in \mathcal{F}_{2n-2}$, and so

$$|V_1''(\xi_i)| \leq C(1 + \xi_i^2)^{n-1} \leq CP^{2n-2}(Q) + CP^{2n-2}(q_i) \leq C \sum_{k=0}^N P^{2n-2}(q_k).$$

Thus, $Y^{(2)}$ can be split in terms of the form

$$|\langle \Lambda_1^{2\varepsilon-2} \tilde{q}_j P^{2n-2}(q_k) f, Qf \rangle| \leq \|\Lambda_1^{1/n-2} P^{2n-2}(q_k) Qf\| \|\Lambda_1^{2\varepsilon-1/n} \tilde{q}_j f\|.$$

The first factor clearly can be bounded by $C\|f\|$ if we notice that $q \mapsto P^{2n-2}(q_k)Q$ belongs to \mathcal{F}_{2n-1} and then apply Proposition 4.1. The second factor can also be bounded by $C\|f\|$ if we impose

$$0 < \varepsilon \leq \frac{1}{2n} - \frac{1}{2m},$$

which can be done because we assumed $n < m$.

It remains to estimate $Y^{(1)}$. We define $P = \sum_{i=0}^N p_i$. Since it may easily be verified that $[X_0, P] = \sum_{i=0}^N V_1'(q_i) - r_L - r_R$, we can write Y_1 as

$$Y_1 = \langle \Lambda_1^{2\varepsilon-2} ([X_0, P] + r_L + r_R) f, Qf \rangle \equiv Y^{(3)} + Y^{(4)} + Y^{(5)}.$$

We leave to the reader the verification that the terms $Y^{(4)}$ and $Y^{(5)}$ can be bounded by $C\|f\|^2$ without introducing any stronger condition on ε . The term $Y^{(3)}$ can be estimated by using Lemma 4.2 with $A = Q$ and $B = P$. We have already verified that (4.22) holds for every i , so we can define

$$\varepsilon_P \equiv \min\{\varepsilon_i \mid i = 0, \dots, N\}.$$

This, together with Proposition 4.1, allows us to choose,

$$\begin{aligned} \alpha_1 &= 1/n, & \beta_1 &= 1 - \varepsilon_P, \\ \alpha_2 &= 1/n, & \beta_2 &= 1, \\ \alpha_3 &= 1, & \beta_3 &= 2 - 1/n, \end{aligned}$$

and thus (4.24) is fulfilled if we choose $2\varepsilon \leq \varepsilon_P$. This completes the proof of the lemma. \square

5. Generalization of Hörmander’s Theorem

In a celebrated paper [Hör67], Hörmander studied second-order differential operators of the form

$$P = \sum_{j=1}^r L_j^* L_j + L_0, \tag{5.1}$$

where the L_j are some smooth vector fields acting in \mathbf{R}^d . He showed that a sufficient condition for the operator P to be hypoelliptic is that the Lie algebra generated by $\{L_0, \dots, L_r\}$ has maximal rank everywhere. The main step in his proof is to show that there exists a constant $\varepsilon > 0$ and, for every compact domain $\mathcal{K} \subset \mathbf{R}^d$, a constant $C_{\mathcal{K}}$ such that

$$\|u\|_{(\varepsilon)} \leq C_{\mathcal{K}}(\|Pu\| + \|u\|), \quad \forall u \in C_0^\infty(\mathcal{K}). \tag{5.2}$$

In this expression, the norm $\|\cdot\|_{(\varepsilon)}$ is the natural norm associated to the Sobolev space $H^\varepsilon(\mathbf{R}^d)$, i.e.

$$\|u\|_{(\varepsilon)}^2 = \int_{\mathbf{R}^d} |\hat{u}(k)|^2 (1 + k^2)^\varepsilon d^d k \equiv \|(1 + \Delta)^{\varepsilon/2} u\|.$$

We base our discussion on the proof presented in [Hör85]. Hörmander first defines \mathcal{Q}_1 as the set of all properly supported symmetric first-order differential operators q such that for every compact domain \mathcal{K} , there exist constants $C'_{\mathcal{K}}$ and $C''_{\mathcal{K}}$ with

$$\|qu\|^2 \leq C'_{\mathcal{K}} \operatorname{Re}\langle Pu, u \rangle + C''_{\mathcal{K}} \|u\|^2, \quad u \in C_0^\infty(\mathcal{K}). \tag{5.3}$$

In particular, if we write $L_j^* = -L_j + c_j$, where c_j is some function, \mathcal{Q}_1 contains all the operators of the form

$$(L_j - c_j/2)/i, \quad j \geq 1,$$

as well as their linear combinations. It also contains every operator of order 0. Hörmander then defines \mathcal{Q}_2 as consisting of the operator $(P - P^*)/i$, as well as all the commutators of the form $[q, q']/i$ with $q, q' \in \mathcal{Q}_1$. For $k > 2$, he defines \mathcal{Q}_k as the set of all commutators $[q, q']/i$ with $q \in \mathcal{Q}_{k-1}$ and $q' \in \mathcal{Q}_{k-2}$. One feature of this construction is that a finite number of steps suffices to catch every symmetric first-order differential operator. This is a consequence of the maximal rank hypothesis. The main point of Hörmander’s proof is then the following result.

Lemma 5.1 (Hörmander). *If $q_k \in \mathcal{Q}_k$ and $\varepsilon \leq 2^{1-k}$, we have for every $\mathcal{K} \subset \mathbf{R}^d$,*

$$\|q_k u\|_{(\varepsilon-1)} \leq C(\|Pu\| + \|u\|), \quad u \in C_0^\infty(\mathcal{K}). \tag{5.4}$$

The proof can be found in [Hör85, p. 355]. The result (5.2) then follows almost immediately, because the operators $i\partial_j$ all belong to some Q_k . Thus there exists some $\varepsilon > 0$ such that

$$\sum_{j=1}^d \|\partial_j u\|_{(\varepsilon^{-1})}^2 \leq C_{\mathcal{K}}(\|Pu\| + \|u\|)^2, \quad u \in \mathcal{C}_0^\infty(\mathcal{K}),$$

which implies (5.2). One of the major problems encountered in this paper is to find a *global* estimate analogous to (5.2), i.e. to find constants C and ε such that

$$\|\tilde{\Delta}^\varepsilon u\| \leq C(\|Pu\| + \|u\|), \quad \text{for all } u \in \mathcal{C}_0^\infty(\mathbf{R}^d),$$

where $\tilde{\Delta}$ is some modified Laplacean. There are two major difficulties:

- If we were to construct the sets Q_k as in [Hör85], they would not necessarily “close” in the sense that the successive commutators could blow up, and the whole proof would break down. To avoid this we do not necessarily put $(P - P^*)/i$ into Q_2 , but rather $g_0(P - P^*)/i$, where g_0 is some bounded function. This allows to get decreasing bounds on the successive commutators. This problem does *not* appear in [EPR99a], where the successive commutators are all first-order differential operators with constant (or bounded) coefficients. On the other hand, the commutator technique is essentially the same as in [EPR99a].
- The above construction does not allow to deal with arbitrary symmetric first-order differential operators. The reason is that if we want a global equivalent of (5.3), the set Q_1 is no longer allowed to contain products of the L_j and unbounded functions. We thus work with fewer operators, which means that we track much more closely the expressions which appear in the constructions.

5.1. General setting. Let us consider the Hilbert space $\mathcal{H} = L^2(\mathbf{R}^d, dx)$ for some integer $d \geq 1$. We define the set $\mathcal{C}(\mathcal{H})$ as the set of closed operators on \mathcal{H} and the algebra $\mathfrak{B}(\mathcal{H})$ as the everywhere defined bounded operators on \mathcal{H} .

We define $\mathcal{D} \equiv \mathcal{C}_0^\infty(\mathbf{R}^d)$, which is dense in \mathcal{H} . Let us fix some sub-algebra $\mathfrak{F} \subset \mathfrak{B}(\mathcal{H})$ that is closed under conjugation and such that $F\mathcal{D} \subset \mathcal{D}$ for all $F \in \mathfrak{F}$ (typically \mathfrak{F} is some algebra of bounded functions). The advantage of considering $\mathcal{C}_0^\infty(\mathbf{R}^d)$ is that *every* differential operator with sufficiently smooth coefficients is closable on it (see [Yos80] for a justification). Moreover, every differential operator with smooth coefficients maps \mathcal{D} into itself. This allows us to make a formal calculus, i.e. every relationship between operators appearing in this section is supposed to hold on \mathcal{D} . The actual operators are then the closures of the operators defined on \mathcal{D} .

We define \mathcal{L} as the set of all formal expressions of the form

$$\sum_{|\ell| \leq k} a_\ell(x) D^\ell, \quad k \geq 0, \quad a \in \mathcal{C}^\infty(\mathbf{R}^d),$$

where D^ℓ denotes the $|\ell|^{\text{th}}$ derivative with respect to the multi-index ℓ . By the above remark, any element of \mathcal{L} can naturally be identified with a differential operator in $\mathcal{C}(\mathcal{H})$.

Consider a differential operator K that can be written as

$$K = \sum_{i=1}^n X_i^* X_i + X_0, \quad X_j \in \mathcal{L}, \quad j = 1, \dots, n, \quad (5.5)$$

where X_0 is such that

$$X_0^* = -X_0 + g, \quad g \in \mathfrak{F}. \quad (5.6)$$

We introduce now a definition that will be very useful in the sequel.

Definition 5.2. Let $\mathcal{S} \subset \mathfrak{L}$ be a finite set of differential operators and $i \geq 0$ a natural number. We define the set $\mathcal{Y}_{\mathfrak{F}}^i(\mathcal{S})$ as the module on \mathfrak{F} generated by the terms

$$S_1 S_2 \cdots S_i, \quad S_k \in \mathcal{S} \cup \{1\}, \quad k = 1, \dots, i.$$

The elements of $\mathcal{Y}_{\mathfrak{F}}^i(\mathcal{S})$ are naturally identified with densely defined closed operators on \mathcal{H} . If $i = 0$, we use the convention $\mathcal{Y}_{\mathfrak{F}}^0(\mathcal{S}) \equiv \mathfrak{F}$.

The subscript \mathfrak{F} will be dropped in the sequel when the algebra \mathfrak{F} is clear from the context. We construct the sets

$$\mathcal{A}_{-1} = \{X_1, \dots, X_n\}, \quad \mathcal{A}_0 = \{g_0 X_0, X_1, \dots, X_n\}, \quad g_0 \in \mathfrak{F}, \quad (5.7)$$

where the operator g_0 is assumed to be self-adjoint, positive and such that

$$[g_0, X_0] \in \mathfrak{F}. \quad (5.8)$$

Let us now construct recursively up to a level $R < \infty$ some finite sets $\mathcal{B}_i, \mathcal{A}_i \subset \mathfrak{L}$ by the following procedure. Assume \mathcal{A}_{i-1} is known. Consider next the set $\mathcal{B}_i^{(0)}$ of all A of the form

$$A = \sum_{B \in \mathcal{A}_{i-1}} \left(f_B B + \sum_{X \in \mathcal{A}_0} f_{XB} [X, B] \right), \quad f_B, f_{XB} \in \mathfrak{F}. \quad (5.9)$$

We then select a finite subset $\mathcal{B}_i \subset \mathcal{B}_i^{(0)}$. The set \mathcal{A}_i is then defined as

$$\mathcal{A}_i \equiv \mathcal{A}_{i-1} \cup \mathcal{B}_i.$$

Remark 5.3. It is here that our construction differs from similar ones where *all* elements of $\mathcal{B}_i^{(0)}$ would have been selected. This makes the set of operators which we study much smaller, but then we of course have to verify that the operators of interest are really covered by our construction.

We will make some working hypotheses on the sets \mathcal{A}_i .

Hypothesis 1. The pair $(\mathcal{A}_R, \mathfrak{F})$ satisfies the following. If $A, B \in \mathcal{A}_R$ and $f \in \mathfrak{F}$, then

$$[A, B] \in \mathcal{Y}^1(\mathcal{A}_R), \quad A^* \in \mathcal{Y}^1(\mathcal{A}_R), \quad [A, f] \in \mathfrak{F}.$$

Hypothesis 2. If $A \in \mathcal{A}_i$ with $i \geq -1$, we have $A^* \in \mathcal{Y}^1(\mathcal{A}_i)$.

Remark 5.4. Hypothesis 1 implies that if $X \in \mathcal{Y}^j(\mathcal{A}_R)$ and $Y \in \mathcal{Y}^k(\mathcal{A}_R)$, then $[X, Y] \in \mathcal{Y}^{k+j-1}(\mathcal{A}_R)$. This will be very useful in the sequel. Hypothesis 2 implies that the classes $\mathcal{Y}^k(\mathcal{A}_i)$ are closed under conjugation.

We define now the operator Λ^2 by

$$\Lambda^2 = 1 + \sum_{A \in \mathcal{A}_R} A^* A. \tag{5.10}$$

This is, in some sense that will immediately be clear from Lemma 5.6, the “biggest” operator contained in $\mathcal{Y}^2(\mathcal{A}_R)$. The operator Λ^2 is symmetric, densely defined and positive. We will moreover assume that

Hypothesis 3. Λ^2 is essentially self-adjoint on \mathcal{D} .

The powers Λ^α thus exist and are also essentially self-adjoint on \mathcal{D} for $\alpha \leq 2$.

5.2. Results and a preliminary lemma. The following theorem is the main result of this section.

Theorem 5.5. *Let K and Λ be defined as above and assume Hypotheses 1–3 are satisfied for some R . Then there exist some constants $C, \varepsilon > 0$ such that for every $f \in \mathcal{D}$, we have*

$$\|\Lambda^\varepsilon f\| \leq C(\|Kf\| + \|f\|). \tag{5.11}$$

In the sequel, we will write \mathcal{A} instead of \mathcal{A}_R to simplify the notation.

In order to prove Theorem 5.5, we need the following lemma, which will be extensively used in the sequel.

Lemma 5.6. *Let Λ, \mathfrak{F} and \mathcal{A} be as above and assume Hypotheses 1 and 3 hold. If $X \in \mathcal{Y}_{\mathfrak{F}}^j(\mathcal{A})$, then the operators*

$$\Lambda^\beta X \Lambda^\gamma \quad \text{with} \quad \beta + \gamma \leq -j$$

are bounded.

If $Y \in \mathfrak{L}$ is such that $[Y, \Lambda^2] \in \mathcal{Y}_{\mathfrak{F}}^j(\mathcal{A})$, then the operators

$$\Lambda^\beta [\Lambda^\alpha, Y] \Lambda^\gamma \quad \text{with} \quad \alpha + \beta + \gamma \leq 2 - j$$

are bounded.

If $X, Y \in \mathfrak{L}$ are such that

$$[X, \Lambda^2] \in \mathcal{Y}_{\mathfrak{F}}^j(\mathcal{A}), \quad [Y, \Lambda^2] \in \mathcal{Y}_{\mathfrak{F}}^k(\mathcal{A}) \quad \text{and} \quad [[\Lambda^2, X], Y] \in \mathcal{Y}_{\mathfrak{F}}^{j+k-2}(\mathcal{A}),$$

then the operators

$$\Lambda^\beta [[\Lambda^\alpha, X], Y] \Lambda^\gamma \quad \text{with} \quad \alpha + \beta + \gamma \leq 4 - j - k$$

are bounded.

Proof. The proof of this lemma is postponed to Appendix A. \square

Remark 5.7. Lemma 5.6 allows us to count powers in the following sense. Each time we see an operator that is a monomial containing fractional powers of Λ and some operators of $\mathcal{Y}^j(\mathcal{A})$, we know that the operator is bounded if its “degree” is less or equal to 0. The rule is that if $Y \in \mathcal{Y}^j(\mathcal{A})$, its degree is j and the degree of Λ^α is α . Moreover, every time we encounter a commutator, we can lower the degree by one unit.

Lemma 5.6 also shows that if $f \in \mathcal{D}$, $A \in \mathcal{A}$ and $\alpha \leq 2$, expressions such as $A\Lambda^\alpha f$ can be well defined by

$$A\Lambda^\alpha f \equiv \Lambda^\alpha Af + [A, \Lambda^\alpha]\Lambda^{-2}\Lambda^2 f,$$

where $[A, \Lambda^\alpha]\Lambda^{-2}$ is bounded and can therefore be defined on all of \mathcal{H} . Similar expressions hold to show that any expression of this section can be well defined.

We are now ready to prove the theorem.

5.3. Proof of Theorem 5.5. The proof uses the commutation techniques developed by Hörmander [Hör85] and improved by Eckmann, Pillet, Rey-Bellet [EPR99a]. Large parts of this proof are inspired from this latter work.

Before we start the proof itself, let us make a few computations, the results of which will be used repeatedly in the sequel. We first show that we can assume $\operatorname{Re}K$ positive. An explicit computation, using (5.5) and (5.6), shows that

$$\operatorname{Re}K = \sum_{i=1}^n X_i^* X_i + \frac{g}{2}, \quad \text{and thus also} \quad X_0 = K - \operatorname{Re}K + g/2. \quad (5.12)$$

Because $g \in \mathfrak{F}$, we can add a sufficiently big constant to X_0 to make $\operatorname{Re}K$ positive. This will change neither the commutation relations, nor the estimate (5.11).

Another useful equality is

$$g_0 \operatorname{Re}K = \operatorname{Re}(g_0 K + K_1) + K_2 \quad K_1, K_2 \in \mathcal{Y}^1(\mathcal{A}_{-1}), \quad (5.13)$$

where K_1 is a self-adjoint operator such that $\operatorname{Re}(g_0 K + K_1)$ is a positive self-adjoint operator. This is a consequence of the following two equalities, which are easily verified by inspection

$$\begin{aligned} g_0 \operatorname{Re}K &= \sum_{i=1}^n X_i^* g_0 X_i + K_2 \quad K_2 \in \mathcal{Y}^1(\mathcal{A}_{-1}), \\ \operatorname{Re}(g_0 K) &= \sum_{i=1}^n X_i^* g_0 X_i - K_1 \quad K_1 \in \mathcal{Y}^1(\mathcal{A}_{-1}). \end{aligned}$$

We therefore have

$$\operatorname{Re}(g_0 K + K_1) = \sum_{i=1}^n X_i^* g_0 X_i.$$

This proves (5.13).

Another useful identity will be

$$\begin{aligned} (g_0 X_0)^* &= -X_0 g_0 + g g_0 = -g_0 X_0 + [g_0, X_0] + g g_0 \\ &= -g_0 X_0 + g'_0, \quad g'_0 \in \mathfrak{F}, \end{aligned} \quad (5.14)$$

where the last equality is a consequence of (5.8).

We will now verify the estimate (5.11) for some vector $f \in \mathcal{D}$. In the sequel, the symbol C will be used to denote some constant depending only on the operator K . This

constant can change from one line to the other. We will first prove that $A \in \mathcal{Y}^1(\mathcal{A}_i)$ with $0 \leq i \leq R$ implies

$$\|\Lambda^{1/4i+1-1}Af\| \leq C(\|Kf\| + \|f\|). \quad (5.15)$$

In fact, an immediate consequence of the first part of Lemma 5.6 is that we only have to prove this assertion for $A \in \mathcal{A}_i$. The proof will proceed by induction on i .

5.3.1. Verification for $i = 0$. We want to verify the estimate

$$\|\Lambda^{-3/4}Af\| \leq C(\|Kf\| + \|f\|), \quad \text{for all } A \in \mathcal{A}_0.$$

The cases $A = g_0X_0$ and $A = X_j$ with $j \neq 0$ will be treated separately.

The case $A = X_j$. We write

$$\begin{aligned} \|\Lambda^{-3/4}X_jf\|^2 &\leq C\|X_jf\|^2 \leq C\langle f, X_j^*X_jf \rangle \leq C\langle f, (K + K^* - g)f \rangle \\ &\leq C\operatorname{Re}\langle f, Kf \rangle + C\|f\|^2 \leq C\|f\|(\|Kf\| + \|f\|). \end{aligned}$$

This implies the desired estimate. Because $X_j^* \in \mathcal{Y}^1(\mathcal{A}_{-1})$ by hypothesis, this computation immediately implies the estimates

$$\|X_jf\| \leq C(\|Kf\| + \|f\|), \quad (5.16a)$$

$$\|X_j^*f\| \leq C(\|Kf\| + \|f\|), \quad (5.16b)$$

which hold for every $j \geq 1$.

The case $A = g_0X_0$. We write, using expression (5.12),

$$\begin{aligned} \|\Lambda^{-3/4}Af\|^2 &= \langle g_0X_0f, \Lambda^{-3/2}Af \rangle \\ &= \langle Kf, g_0\Lambda^{-3/2}Af \rangle + \langle g_0gf, \Lambda^{-3/2}Af \rangle / 2 - \langle (\operatorname{Re}K)f, g_0\Lambda^{-3/2}Af \rangle \\ &\equiv S_1 + S_2 - S_3. \end{aligned}$$

The terms S_1 and S_2 are easily bounded by $C(\|Kf\| + \|f\|)^2$, using the Cauchy-Schwarz inequality and the first part of Lemma 5.6. Using the positivity of $\operatorname{Re}K$ and the explicit form of K , the term S_3 can be bounded as

$$\begin{aligned} |S_3| &= \langle (\operatorname{Re}K)^{1/2}f, (\operatorname{Re}K)^{1/2}g_0\Lambda^{-3/2}Af \rangle \\ &\leq |\operatorname{Re}\langle Kf, f \rangle|^{1/2} |\langle (\operatorname{Re}K)g_0\Lambda^{-3/2}Af, g_0\Lambda^{-3/2}Af \rangle|^{1/2} \\ &\leq \sqrt{\|Kf\|\|f\|} \left| \langle g_0\Lambda^{-3/2}Af, g_0\Lambda^{-3/2}Af \rangle / 2 + \sum_{i=1}^n \|X_i g_0\Lambda^{-3/2}Af\|^2 \right|^{1/2} \\ &\equiv \sqrt{\|Kf\|\|f\|} \sqrt{S_0 + \sum_{i=1}^n S_{0,i}^2}. \end{aligned}$$

The term S_0 is estimated by simple power counting (the Λ 's contribute for -3 and the A 's for 2 in the total degree of the expression, hence $|S_0| \leq C\|f\|^2$). The terms $S_{0,i}$ are estimated by writing

$$|S_{0,i}| \leq \|g_0\Lambda^{-3/2}AX_i f\| + \|[X_i, g_0\Lambda^{-3/2}A]f\|.$$

The first term is estimated by using (5.16) and power counting. The second term is estimated by expanding the commutator as

$$[X_i, g_0 \Lambda^{-3/2} A] = [X_i, g_0] \Lambda^{-3/2} A + g_0 [X_i, \Lambda^{-3/2}] A + g_0 \Lambda^{-3/2} [X_i, A],$$

and estimating separately the resulting terms.

5.3.2. The induction hypothesis. We shall proceed by induction. Let us fix $j > 0$, take $A \in \mathcal{A}_j$ and assume (5.15) holds for $i < j$. Let us moreover define $\varepsilon \equiv 1/4^{j+1}$ in order to simplify the notation. Our assumption is therefore that

$$\|\Lambda^{4\varepsilon-1} Bf\| \leq C(\|Kf\| + \|f\|) \quad \forall B \in \mathcal{Y}^1(\mathcal{A}_{j-1}). \quad (5.17)$$

We will now prove that this assumption implies the desired estimate, i.e.

$$\|\Lambda^{\varepsilon-1} Af\| \leq C(\|Kf\| + \|f\|) \quad \forall A \in \mathcal{Y}^1(\mathcal{A}_j). \quad (5.18)$$

This, together with the preceding paragraph, will imply the estimate (5.15).

5.3.3. Proof of the main estimate. Because of the induction hypothesis, we only have to check (5.18) for $A \in \mathcal{A}_j \setminus \mathcal{A}_{j-1}$. By (5.9), we can write

$$A = \sum_{B \in \mathcal{A}_{j-1}} \left(f_B B + f_B^0 [g_0 X_0, B] + \sum_{i=1}^n f_B^i [X_i, B] \right),$$

with all the f belonging to \mathfrak{F} . We have

$$\begin{aligned} \|\Lambda^{\varepsilon-1} Af\|^2 &= \sum_{B \in \mathcal{A}_{j-1}} \left\langle \left(f_B B + f_B^0 [g_0 X_0, B] + \sum_{i=1}^n f_B^i [X_i, B] \right) f, \Lambda^{2\varepsilon-2} Af \right\rangle \\ &\equiv \sum_{B \in \mathcal{A}_{j-1}} \left(T_B + T_B^0 + \sum_{i=1}^n T_B^i \right). \end{aligned}$$

We are going to bound each term of this sum separately by $C(\|Kf\| + \|f\|)^2$.

Term T_B . We have

$$|T_B| = |\langle \Lambda^{2\varepsilon-1} f_B \Lambda^{1-2\varepsilon} \Lambda^{2\varepsilon-1} Bf, \Lambda^{-1} Af \rangle|.$$

The operators $\Lambda^{2\varepsilon-1} f_B \Lambda^{1-2\varepsilon}$ and $\Lambda^{-1} A$ are bounded by Lemma 5.6. Using the induction hypothesis (5.17), we thus get the bound $|T_B| \leq C(\|Kf\| + \|f\|)^2$.

Term T_B^i with $i \neq 0$. We define $h \equiv f_B^i$. The term T_B^i is then written as

$$T_B^i = \langle Bf, X_i^* h^* \Lambda^{2\varepsilon-2} Af \rangle - \langle X_i f, h^* B^* \Lambda^{2\varepsilon-2} Af \rangle \equiv Q_1 - Q_2.$$

Term Q_1 . It can be estimated by writing

$$Q_1 = \langle Bf, h^* \Lambda^{2\varepsilon-2} A X_i^* f \rangle + \langle Bf, [X_i^*, h^* \Lambda^{2\varepsilon-2} A] f \rangle.$$

The first term is estimated by rewriting it as

$$\begin{aligned} |\langle Bf, h^* \Lambda^{2\varepsilon-2} A X_i^* f \rangle| &= |\langle \Lambda^{2\varepsilon-1} Bf, \Lambda^{1-2\varepsilon} h^* \Lambda^{2\varepsilon-2} A X_i^* f \rangle| \\ &\leq \|\Lambda^{2\varepsilon-1} Bf\| \|\Lambda^{1-2\varepsilon} h^* \Lambda^{2\varepsilon-2} A X_i^* f\| \\ &\leq C(\|Kf\| + \|f\|)^2. \end{aligned}$$

The last inequality has been obtained by using the induction hypothesis (5.17), the estimate (5.16b) and the fact that the operator $\Lambda^{1-2\varepsilon} h^* \Lambda^{2\varepsilon-2} A$ is bounded by Lemma 5.6.

The second term is estimated as

$$\begin{aligned} |\langle Bf, [X_i^*, h^* \Lambda^{2\varepsilon-2} A] f \rangle| &= |\langle \Lambda^{2\varepsilon-1} Bf, \Lambda^{1-2\varepsilon} [X_i^*, h^* \Lambda^{2\varepsilon-2} A] f \rangle| \\ &\leq \|\Lambda^{2\varepsilon-1} Bf\| \|\Lambda^{1-2\varepsilon} [X_i^*, h^* \Lambda^{2\varepsilon-2} A] f\|. \end{aligned}$$

The term $\|\Lambda^{2\varepsilon-1} Bf\|$ is bounded by the induction hypothesis (5.17). The other term can be estimated by writing the commutator as

$$[X_i^*, h^* \Lambda^{2\varepsilon-2} A] = [X_i^*, h^*] \Lambda^{2\varepsilon-2} A + h^* [X_i^*, \Lambda^{2\varepsilon-2}] A + h^* \Lambda^{2\varepsilon-2} [X_i^*, A].$$

The resulting terms are estimated by power counting, using the fact that $X_i^* \in \mathcal{Y}^1(\mathcal{A})$.

Term Q_2 . We bound this term as

$$\begin{aligned} |Q_2| &= |\langle X_i f, h^* \Lambda^{2\varepsilon-2} A B^* f \rangle + \langle X_i f, h^* [B^*, \Lambda^{2\varepsilon-2} A] f \rangle| \\ &\leq \|X_i f\| (\|h^* \Lambda^{2\varepsilon-2} A B^* f\| + \|h^* [B^*, \Lambda^{2\varepsilon-2} A] f\|) \\ &\leq \|X_i f\| (\|h^* \Lambda^{2\varepsilon-2} A \Lambda^{1-2\varepsilon}\| \|\Lambda^{2\varepsilon-1} B^* f\| + \|h^* [B^*, \Lambda^{2\varepsilon-2} A] f\|). \end{aligned}$$

We leave to the reader the not too hard task to verify that it is indeed possible to get the bound $|Q_2| \leq C(\|Kf\| + \|f\|)^2$ by similar estimates as for the term Q_1 .

Term T_B^0 . We define $h \equiv f_B^0$. The term T_B^0 is thus equal to

$$\begin{aligned} T_B^0 &= \langle [g_0 X_0, B] f, h^* \Lambda^{2\varepsilon-2} A f \rangle \\ &= \langle g_0 X_0 B f, h^* \Lambda^{2\varepsilon-2} A f \rangle - \langle B g_0 X_0 f, h^* \Lambda^{2\varepsilon-2} A f \rangle. \end{aligned}$$

We use (5.14) to write this as

$$\begin{aligned} T_B^0 &= -\langle Bf, h^* \Lambda^{2\varepsilon-2} A g_0 X_0 f \rangle + \langle Bf, g_0' h^* \Lambda^{2\varepsilon-2} A f \rangle \\ &\quad - \langle Bf, [g_0 X_0, h^* \Lambda^{2\varepsilon-2} A] f \rangle - \langle B g_0 X_0 f, h^* \Lambda^{2\varepsilon-2} A f \rangle \\ &\equiv -U_1 + U_2 - U_3 - U_4, \end{aligned}$$

where $g_0' \in \mathfrak{F}$. The term U_2 can easily be estimated by

$$\begin{aligned} |U_2| &= |\langle \Lambda^{2\varepsilon-1} Bf, \Lambda^{1-2\varepsilon} g_0' h^* \Lambda^{2\varepsilon-2} A f \rangle| \leq \|\Lambda^{2\varepsilon-1} Bf\| \|\Lambda^{1-2\varepsilon} g_0' h^* \Lambda^{2\varepsilon-2} A f\| \\ &\leq C(\|Kf\| + \|f\|) \|f\|, \end{aligned}$$

using the induction hypothesis. In order to estimate the term U_3 , we notice that $g_0 X_0 \in \mathcal{A}$, and thus $[g_0 X_0, \Lambda^2] \in \mathcal{Y}^2(\mathcal{A})$. We can therefore write

$$\begin{aligned} |U_3| &= |\langle \Lambda^{2\varepsilon-1} Bf, \Lambda^{1-2\varepsilon} [g_0 X_0, h^* \Lambda^{2\varepsilon-2} A] f \rangle| \\ &\leq \|\Lambda^{2\varepsilon-1} Bf\| \|\Lambda^{1-2\varepsilon} [g_0 X_0, h^* \Lambda^{2\varepsilon-2} A] f\|, \end{aligned}$$

expand the commutator and estimate the resulting terms separately by power counting. We use the equality

$$X_0 = K - \operatorname{Re}K + g/2,$$

to write the terms U_1 and U_4 as

$$\begin{aligned} U_1 &= \langle Bf, h^* \Lambda^{2\varepsilon-2} A g_0 (K - (\operatorname{Re}K) + g/2) f \rangle \equiv T_{B,1} - T_{B,2} + T_{B,3}, \\ U_4 &= \langle B g_0 (K - (\operatorname{Re}K) + g/2) f, h^* \Lambda^{2\varepsilon-2} A f \rangle \equiv T_{B,4} - T_{B,5} + T_{B,6}. \end{aligned}$$

Each of these terms will now be estimated separately.

Terms $T_{B,3}$ and $T_{B,6}$. They are easily bounded like the term U_2 by power counting and using the induction hypothesis to bound $\|\Lambda^{2\varepsilon-1} Bf\|$. In the case of $T_{B,6}$, we first have to commute B with $g_0 g/2$, but this does not cause any problem.

Term $T_{B,1}$. This term can be estimated by

$$|T_{B,1}| \leq \|Kf\| \|g_0^* A^* \Lambda^{2\varepsilon-2} h Bf\| \leq \|Kf\| \|g_0^* A^* \Lambda^{2\varepsilon-2} h \Lambda^{2-2\varepsilon}\| \|\Lambda^{2\varepsilon-2} Bf\|.$$

The norm of $g_0^* A^* \Lambda^{2\varepsilon-2} h \Lambda^{2-2\varepsilon}$ is bounded by power counting. Using the induction hypothesis (5.17), we thus have $|T_{B,1}| \leq C(\|Kf\| + \|f\|)^2$.

Term $T_{B,4}$. We have the estimate

$$|T_{B,4}| = |\langle Kf, g_0^* B^* h^* \Lambda^{2\varepsilon-2} A f \rangle| \leq \|Kf\| \|g_0^* B^* h^* \Lambda^{2\varepsilon-2} A f\|.$$

The second norm can be estimated by writing

$$\begin{aligned} \|g_0^* B^* h^* \Lambda^{2\varepsilon-2} A f\| &\leq \|g_0^* h^* \Lambda^{2\varepsilon-2} A \Lambda^{1-2\varepsilon}\| \|\Lambda^{2\varepsilon-1} B^* f\| \\ &\quad + \|g_0^* [B^*, h^* \Lambda^{2\varepsilon-2} A] f\|. \end{aligned}$$

Here, the first term can be bounded by $C(\|Kf\| + \|f\|)$ because, by Hypothesis 2, we have $B^* \in \mathcal{Y}^1(\mathcal{A}_{j-1})$ and so we can use the induction hypothesis. The commutator can be expanded and bounded by power counting.

Term $T_{B,2}$. We can write this term as

$$\begin{aligned} T_{B,2} &= \langle \Lambda^{2\varepsilon-1} h Bf, (g_0 \operatorname{Re}K) \Lambda^{-1} A f \rangle + \langle \Lambda^{2\varepsilon-1} h Bf, [\Lambda^{-1} A, g_0 \operatorname{Re}K] f \rangle \\ &= \langle \Lambda^{2\varepsilon-1} h Bf, K_2 \Lambda^{-1} A f \rangle + \langle \Lambda^{2\varepsilon-1} h Bf, [\Lambda^{-1} A, g_0 \operatorname{Re}K] f \rangle \\ &\quad + \langle \Lambda^{2\varepsilon-1} h Bf, \operatorname{Re}(g_0 K + K_1) \Lambda^{-1} A f \rangle \\ &\equiv M_1 + M_2 + M_3, \end{aligned}$$

where the second equality has been obtained using (5.13). These terms can now be estimated separately.

Term M_1 . We write this term as

$$M_1 = \langle \Lambda^{2\varepsilon-1} h Bf, \Lambda^{-1} A K_2 f \rangle + \langle \Lambda^{2\varepsilon-1} h Bf, [K_2, \Lambda^{-1} A] f \rangle.$$

The first term is estimated by using

$$K_2 \in \mathcal{Y}^1(\mathcal{A}_{-1}) \quad \Rightarrow \quad \|K_2 f\| \leq C(\|Kf\| + \|f\|),$$

where the implication is a straightforward consequence of (5.16a). The second term can be estimated by power counting and the induction hypothesis, using the fact that $K_2 \in \mathcal{Y}^1(\mathcal{A}_{-1})$, so that $[K_2, \Lambda^{-1}A]$ is bounded.

Term M_2 . We use the explicit form of $\text{Re}K$ to write this term as

$$\begin{aligned} M_2 &= \langle \Lambda^{2\varepsilon-1} hBf, [\Lambda^{-1}A, g_0](\text{Re}K)f \rangle \\ &\quad + \sum_{i=1}^n \left(\langle \Lambda^{2\varepsilon-1} hBf, g_0 X_i^* [\Lambda^{-1}A, X_i] f \rangle \right. \\ &\quad \quad \left. + \langle \Lambda^{2\varepsilon-1} hBf, g_0 [\Lambda^{-1}A, X_i^*] X_i f \rangle \right) \\ &\quad + \langle \Lambda^{2\varepsilon-1} hBf, g_0 [\Lambda^{-1}A, g] f \rangle / 2 \\ &\equiv M_{20} + \sum_{i=1}^n (M_{i1} + M_{i2}) + M_{21}. \end{aligned}$$

The term M_{20} is estimated by using the explicit form of $\text{Re}K$ to decompose it in terms of the form

$$|\langle \Lambda^{2\varepsilon-1} hBf, [\Lambda^{-1}A, g_0] X_i^* X_i f \rangle| \leq \|[\Lambda^{-1}A, g_0] X_i^*\| \|\Lambda^{2\varepsilon-1} hBf\| \|X_i f\|.$$

The norm $\|[\Lambda^{-1}A, g_0] X_i^*\|$ is finite by Lemma 5.6. The terms $\|\Lambda^{2\varepsilon-1} hBf\|$ and $\|X_i f\|$ are bounded by $C(\|Kf\| + \|f\|)$, using the induction hypothesis (5.17) and the estimate (5.16a) respectively. The terms M_{21} and M_{i2} are estimated by power counting and the induction hypothesis. In order to estimate the term M_{i1} , we have to commute once more to find

$$M_{i1} = \langle \Lambda^{2\varepsilon-1} hBf, g_0 [\Lambda^{-1}A, X_i] X_i^* f \rangle + \langle \Lambda^{2\varepsilon-1} hBf, g_0 [X_i^*, [\Lambda^{-1}A, X_i]] f \rangle.$$

The first term is estimated by using (5.16b). The second term is estimated by expanding the double commutator and power counting.

Term M_3 . We use the positivity of $\text{Re}(g_0K + K_1)$ to write

$$\begin{aligned} |M_3| &= \left((\text{Re}(g_0K + K_1))^{1/2} \Lambda^{2\varepsilon-1} hBf, (\text{Re}(g_0K + K_1))^{1/2} \Lambda^{-1} Af \right) \\ &\leq |\text{Re} \langle (g_0K + K_1) \Lambda^{2\varepsilon-1} hBf, \Lambda^{2\varepsilon-1} hBf \rangle|^{1/2} \\ &\quad \times |\langle \text{Re}(g_0K + K_1) \Lambda^{-1} Af, \Lambda^{-1} Af \rangle|^{1/2} \\ &\leq \sqrt{|\text{Re}M_4| + |\text{Re}M_5|} \sqrt{|M_6|}. \end{aligned}$$

We will now estimate M_4 , M_5 and M_6 separately.

Term M_4 . We want to put the operator g_0K to the left of f . So we write

$$\begin{aligned} M_4 &= \langle \Lambda^{-1} hB g_0 K f, \Lambda^{4\varepsilon-1} hBf \rangle + \langle [g_0K, \Lambda^{2\varepsilon-1} hB] f, \Lambda^{2\varepsilon-1} hBf \rangle \\ &\equiv M_{41} + M_{42}. \end{aligned}$$

The term M_{41} is estimated easily by using the induction hypothesis and the fact that $\Lambda^{-1}hBg_0$ is bounded. In order to estimate M_{42} , we use the explicit form of K to write

$$\begin{aligned} M_{42} &= \langle \Lambda^{-2\varepsilon} [g_0 X_0, \Lambda^{2\varepsilon-1} hB] f, \Lambda^{4\varepsilon-1} hBf \rangle \\ &\quad + \sum_{i=1}^n \left(\langle g_0 X_i^* [X_i, \Lambda^{2\varepsilon-1} hB] f, \Lambda^{2\varepsilon-1} hBf \rangle \right. \\ &\quad \quad \left. + \langle g_0 [X_i^*, \Lambda^{2\varepsilon-1} hB] X_i f, \Lambda^{2\varepsilon-1} hBf \rangle \right) \\ &\quad + \langle \Lambda^{-2\varepsilon} [g_0, \Lambda^{2\varepsilon-1} hB] Kf, \Lambda^{4\varepsilon-1} hBf \rangle \\ &\equiv M_{40} + \sum_{i=1}^n (M_{i3} + M_{i4}) + M_{4K}. \end{aligned}$$

The terms M_{40} and M_{4K} are estimated by expanding the commutator and power counting. The term M_{i4} can be written as

$$\begin{aligned} |M_{i4}| &= |\langle \Lambda^{-2\varepsilon} g_0 [X_i^*, \Lambda^{2\varepsilon-1} hB] X_i f, \Lambda^{4\varepsilon-1} hBf \rangle| \\ &\leq \|\Lambda^{1-4\varepsilon} h^* \Lambda^{2\varepsilon-1} g_0 [X_i^*, \Lambda^{2\varepsilon-1} hB]\| \|X_i f\| \|\Lambda^{4\varepsilon-1} Bf\|. \end{aligned}$$

It is then estimated by power counting, using moreover the induction hypothesis and the estimate (5.16). In order to estimate the term M_{i3} , we have to commute once more to write

$$\begin{aligned} M_{i3} &= \langle \Lambda^{-2\varepsilon} g_0 [X_i, \Lambda^{2\varepsilon-1} hB] X_i^* f, \Lambda^{4\varepsilon-1} hBf \rangle \\ &\quad + \langle \Lambda^{-2\varepsilon} g_0 [X_i^* [X_i, \Lambda^{2\varepsilon-1} hB]] f, \Lambda^{4\varepsilon-1} hBf \rangle. \end{aligned}$$

The first term is estimated exactly like M_{i4} . The second term can then be estimated by expanding the double commutator and power counting.

Term M_5 . We write this term as

$$M_5 = \langle \Lambda^{-1} hBK_1 f, \Lambda^{4\varepsilon-1} hBf \rangle + \langle \Lambda^{-2\varepsilon} [K_1, \Lambda^{2\varepsilon-1} hB] f, \Lambda^{4\varepsilon-1} hBf \rangle.$$

The first term is estimated using the induction hypothesis and the fact that (5.13) and (5.16) imply

$$K_1 \in \mathcal{Y}^1(\mathcal{A}_{-1}) \quad \Rightarrow \quad \|K_1 f\| \leq C(\|Kf\| + \|f\|). \quad (5.19)$$

The other term is estimated by using the fact that $[K_1, \Lambda^2] \in \mathcal{Y}^2(\mathcal{A})$ and $[K_1, hB] \in \mathcal{Y}^1(\mathcal{A})$, which follows from $\mathcal{A}_{-1} \subset \mathcal{A}$ and thus $K_1 \in \mathcal{Y}^1(\mathcal{A})$.

Term M_6 . We use the explicit expression for $\text{Re}(g_0 K + K_1)$ to write this term as

$$M_6 = \sum_{i=1}^n \|g_0^{1/2} X_i \Lambda^{-1} Af\|^2 \leq C \sum_{i=1}^n \|X_i \Lambda^{-1} Af\|^2.$$

These terms are easily estimated by putting the X_i to the left of f , using (5.16) and estimating the commutators.

Term $T_{B,5}$. This is the last term we have to estimate. Using the expression (5.13) and the positivity of $\text{Re}(g_0K + K_1)$, it can be written in the form

$$\begin{aligned} T_{B,5} &= \left(\langle \text{Re}(g_0K + K_1) \rangle^{1/2} f, \left(\text{Re}(g_0K + K_1) \right)^{1/2} B^* h^* \Lambda^{2\varepsilon-2} A f \right) \\ &\quad + \langle K_2 f, B^* h^* \Lambda^{2\varepsilon-2} A f \rangle \\ &\equiv N_1 + N_2. \end{aligned}$$

These terms are now estimated separately.

Term N_2 . We use the Cauchy-Schwarz inequality to write

$$|N_2| \leq \|K_2 f\| \|B^* h^* \Lambda^{2\varepsilon-2} A f\| \equiv \|K_2 f\| \|N_3\|.$$

We can estimate N_3 by writing

$$B^* h^* \Lambda^{2\varepsilon-2} A = h^* \Lambda^{2\varepsilon-2} A \Lambda^{1-2\varepsilon} \Lambda^{2\varepsilon-1} B^* + [B^*, h^* \Lambda^{2\varepsilon-2} A],$$

and estimating the resulting terms using the induction hypothesis. We already noticed that we have the desired estimate for $\|K_2 f\|$.

Term N_1 . Using the Cauchy-Schwarz inequality, we write it as

$$\begin{aligned} N_1 &\leq \langle f, \text{Re}(g_0K + K_1) f \rangle^{1/2} \langle \text{Re}(g_0K + K_1) B^* h^* \Lambda^{2\varepsilon-2} A f, B^* h^* \Lambda^{2\varepsilon-2} A f \rangle^{1/2} \\ &\leq C(\|Kf\| + \|f\|) |\langle \Lambda^{-2\varepsilon} (g_0K + K_1) B^* h^* \Lambda^{2\varepsilon-2} A f, \Lambda^{2\varepsilon} B^* h^* \Lambda^{2\varepsilon-2} A f \rangle|^{1/2} \\ &\equiv C(\|Kf\| + \|f\|) \sqrt{|\langle f_1 + f_2, f_3 \rangle|} \leq C(\|Kf\| + \|f\|) \sqrt{(\|f_1\| + \|f_2\|) \|f_3\|}. \end{aligned}$$

Estimate of $\|f_3\|$. We write it as

$$f_3 = \Lambda^{2\varepsilon} h^* \Lambda^{2\varepsilon-2} A \Lambda^{1-4\varepsilon} \Lambda^{4\varepsilon-1} B^* f + \Lambda^{2\varepsilon} [B^*, h^* \Lambda^{2\varepsilon-2} A] f.$$

The first term is estimated by using the recurrence hypothesis and the fact that Hypothesis 2 implies $B^* \in \mathcal{Y}^1(\mathcal{A}_{j-1})$. The second term is estimated by power counting and by using the fact that $\varepsilon < 1/4$.

Estimate of $\|f_2\|$. We write it as

$$f_2 = \Lambda^{-2\varepsilon} B^* h^* \Lambda^{2\varepsilon-2} A K_1 f + \Lambda^{-2\varepsilon} [K_1, B^* h^* \Lambda^{2\varepsilon-2} A] f.$$

The first term is estimated using the fact that $\|K_1 f\| \leq C(\|Kf\| + \|f\|)$ and power counting. The second term is simply estimated by power counting, and the fact that $K_1 \in \mathcal{Y}^1(\mathcal{A})$.

Estimate of $\|f_1\|$. We use the explicit form of K to write f_1 as

$$\begin{aligned} f_1 &= \Lambda^{-2\varepsilon} B^* h^* \Lambda^{2\varepsilon-2} A g_0 K f + \Lambda^{-2\varepsilon} [g_0 X_0, B^* h^* \Lambda^{2\varepsilon-2} A] f \\ &\quad + \sum_{i=1}^n (\Lambda^{-2\varepsilon} g_0 X_i^* [X_i, B^* h^* \Lambda^{2\varepsilon-2} A] f + \Lambda^{-2\varepsilon} [g_0 X_i^*, B^* h^* \Lambda^{2\varepsilon-2} A] X_i f) \\ &\equiv Q_K + Q_0 + \sum_{i=1}^n (Q_{i,1} + Q_{i,2}). \end{aligned}$$

These terms will now be estimated separately.

Term Q_K . We notice that the operator

$$\Lambda^{-2\varepsilon} B^* h^* \Lambda^{2\varepsilon-2} A g_0$$

is bounded by power counting. This yields the desired estimate.

Term Q_0 . This term is bounded by $C\|f\|$ by power counting, noticing that $g_0 X_0 \in \mathcal{A}$.

Term $Q_{i,2}$. This term can be estimated by power counting if we expand the commutator and use the estimate (5.16).

Term $Q_{i,1}$. We use once more the trick that consists of putting the X_i^* to the left of f . We write therefore

$$Q_{i,1} = \Lambda^{-2\varepsilon} g_0 [X_i, B^* h^* \Lambda^{2\varepsilon-2} A] X_i^* f + \Lambda^{-2\varepsilon} g_0 [X_i^* [X_i, B^* h^* \Lambda^{2\varepsilon-2} A]] f.$$

The first term is estimated by using (5.16b) and expanding the commutator. The second term is estimated in a similar way by expanding the double commutator. We don't write the resulting terms here, because there are too much of them. They are all bounded by simple power counting and by using Lemma 5.6. This completes the proof of estimate (5.18).

It is now straightforward to prove the theorem. Recall that R is the level up to which the \mathcal{A}_i are defined. We put $\varepsilon = 1/4^{R+1}$, and we write:

$$\begin{aligned} \|\Lambda^\varepsilon f\| &= \langle f, \Lambda^{2\varepsilon-2} \Lambda^2 f \rangle = \sum_{A \in \mathcal{A}} \langle f, \Lambda^{2\varepsilon-2} A^* A f \rangle \\ &= \sum_{A \in \mathcal{A}} (\|\Lambda^{\varepsilon-1} A f\|^2 + \langle f, [A^*, \Lambda^{2\varepsilon-2}] A f \rangle). \end{aligned}$$

The first term in the sum is bounded by using (5.15), the second term by simple power counting. This finally completes the proof of Theorem 5.5. \square

We next note a consequence of this theorem, namely a simple criterion to see if a quadratic differential operator has compact resolvent. It is an easy illustration of the technique that will be used in the sequel to show that K has compact resolvent.

5.4. Quadratic differential operators.

Definition 5.8. An operator $A : \mathcal{D}(A) \rightarrow \mathcal{H}$ is called *accretive* if it satisfies

$$\operatorname{Re}\langle f, A f \rangle \geq 0, \quad \text{for all } f \in \mathcal{D}(A).$$

An operator A is called *quasi accretive* if there exists $\lambda \in \mathbf{R}$ such that $A + \lambda$ is accretive. It is called *strictly accretive* if there exists $\lambda > 0$ such that $A - \lambda$ is still accretive.

If $-A$ is accretive, A is called *dissipative*. An operator A is called *m-accretive* if it is accretive and if $(A + \lambda)^{-1}$ exists for all $\lambda > 0$ and satisfies $\|(A + \lambda)^{-1}\| \leq \lambda^{-1}$. The expressions *m-dissipative*, *quasi dissipative*, etc. are defined similarly in an obvious way. An equivalent characterization of *m-accretive* operators is that they are accretive with no proper accretive extension.

It is a classical result (see e.g. [Dav80]) that the quasi *m-dissipative* operators are precisely the generators of quasi-bounded semigroups. An immediate consequence is that if an operator A is (quasi) *m-accretive* (*m-dissipative*), its adjoint A^* is also (quasi) *m-accretive* (*m-dissipative*).

Proposition 5.9. *Let \mathcal{H} be a Hilbert space and \mathcal{C} be a dense subset of \mathcal{H} . Let $K : \mathcal{D}(K) \rightarrow \mathcal{H}$ be a quasi m -accretive (or quasi m -dissipative) operator and let $\Lambda^2 : \mathcal{D}(\Lambda^2) \rightarrow \mathcal{H}$ be a self-adjoint positive operator such that $\mathcal{C} \subset \mathcal{D}(\Lambda^2)$. Assume moreover that \mathcal{C} is a core for K , that Λ^2 has compact resolvent and that there are constants $C > 0$ and $0 < \varepsilon < 2$ such that*

$$\|\Lambda^\varepsilon f\| \leq C(\|Kf\| + \|f\|), \quad \text{for all } f \in \mathcal{C}. \tag{5.20}$$

Then K has compact resolvent too.

Proof. By assumption, there exists a constant $\lambda > 0$ such that $K + \lambda$ is strictly m -accretive. Moreover, (5.20) with K replaced by $K + \lambda$ holds if we change the constant C . Since \mathcal{C} is a core for K , a simple approximation argument shows that $\mathcal{D}(K) \subset \mathcal{D}(\Lambda^\varepsilon)$ and that (5.20) holds for every $f \in \mathcal{D}(K)$.

This immediately implies that $(K + \lambda)^*(K + \lambda)$ has compact resolvent. Since $(K + \lambda)$ is strictly m -accretive, it is invertible and the operator

$$((K + \lambda)^*(K + \lambda))^{-1} = (K + \lambda)^{-1}((K + \lambda)^{-1})^*,$$

is compact. Moreover, we know that $(K + \lambda)^{-1}$ is closed, so we can make the polar decomposition

$$(K + \lambda)^{-1} = PJ,$$

with P self-adjoint and J unitary. Thus P^2 is compact. By the spectral theorem and the characterization of compact operators, this immediately implies P compact, and thus also PJ compact. Thus K has compact resolvent. \square

We now consider $\mathcal{H} = L^2(\mathbf{R}^d)$ and $\mathfrak{F} = \{\lambda I \mid \lambda \in \mathbf{R}\}$, where I is the identity operator in \mathcal{H} . We define the formal expressions

$$\begin{aligned} x^T &= (x_1, \dots, x_d), \\ \partial_x^T &= (\partial_{x_1}, \dots, \partial_{x_d}). \end{aligned}$$

Let $A : \mathbf{R}^d \rightarrow \mathbf{R}^d$ be a linear map and

$$\mathcal{B} = \{b_i \in \mathbf{R}^d \mid i = 1, \dots, s\}, \quad \mathcal{C} = \{c_i \in \mathbf{R}^d \mid i = 1, \dots, t\},$$

two vector families. Let us consider the differential operator K defined as the closure on $C_0^\infty(\mathbf{R}^d)$ of

$$K = - \sum_{i=1}^s \partial_x^T b_i b_i^T \partial_x + \sum_{j=1}^t x^T c_j c_j^T x + x^T A \partial_x. \tag{5.21}$$

We are interested in giving a geometrical condition on A , \mathcal{B} and \mathcal{C} that implies the compactness of the resolvent of K , and therefore the discreteness of its spectrum. It is possible to prove that K is quasi m -accretive. Just follow the proof of Proposition B.3, replacing $G(x)$ by $x^T x$.

We have the following result.

Proposition 5.10. *A sufficient condition for the resolvent of the operator K defined in (5.21) to be compact is that the vector families*

$$\bigcup_{N \geq 0} (A^T)^N \mathcal{B} \quad \text{and} \quad \bigcup_{N \geq 0} A^N \mathcal{C} \tag{5.22}$$

span the whole space \mathbf{R}^n .

Remark 5.11. The intuitive meaning of this theorem is that we can apply Hörmander’s criterion in both direct and Fourier space to obtain an estimate of the form

$$\|H^\varepsilon f\| \leq C(\|Kf\| + \|f\|), \quad H = -\partial_x^T \partial_x + x^T x. \tag{5.23}$$

It is well known that H has compact resolvent. By Proposition 5.9, (5.23) implies that K has compact resolvent.

Proof. We have the following relations:

$$\begin{aligned} [x^T A \partial_x, b^T \partial_x] &\equiv \sum_{i,j,k} [x_i a_{ij} \partial_{x_j}, b_k \partial_{x_k}] = \sum_{i,j,k} b_k [x_i, \partial_{x_k}] a_{ij} \partial_{x_j} \\ &= - \sum_{i,j,k} b_k \delta_{ki} a_{ij} \partial_{x_j} = -b^T A \partial_x, \\ [x^T A \partial_x, c^T x] &\equiv \sum_{i,j,k} [x_i a_{ij} \partial_{x_j}, c_k x_k] = \sum_{i,j,k} x_i a_{ij} [\partial_{x_j}, x_k] c_k \\ &= \sum_{i,j,k} x_i a_{ij} \delta_{jk} c_k = c^T A^T x. \end{aligned}$$

We take $g_0 = 1$, so we have $\mathcal{A}_0 = \mathcal{A}_{-1} \cup \{x^T A \partial_x\}$. We construct the remaining \mathcal{A}_i by

$$\mathcal{B}_i \equiv [x^T A \partial_x, \mathcal{B}_{i-1}].$$

It is very easy to verify Hypotheses 1 and 2, because the assumptions we made on A , \mathcal{B} and \mathcal{C} imply that $\mathcal{Y}^1(\mathcal{A})$ contains every operator of the form $b^T \partial_x$ or $c^T x$. We have moreover

$$(b^T \partial_x)^* = -b^T \partial_x \quad \text{and} \quad (c^T x)^* = c^T x.$$

It is well-known that Hypothesis 3 concerning the essential self-adjointness of the Λ^2 constructed in Theorem 5.5 holds. Finally, it is straightforward that Λ^2 satisfies $\Lambda^2 \geq CH$, where H is the “harmonic oscillator” defined in (5.23).

This proves the validity of (5.23), and hence of the assertion. \square

The interested reader may verify that Proposition 5.10 is quite stable under perturbations. A similar result indeed still holds when the coefficients b_i and c_i are not constants, but functions in \mathcal{F}_0 . This is precisely what was proved in [EPR99a].

6. Proof of the Bound in Momentum Space (Proposition 3.8)

This proposition is an application of Theorem 5.5. It is just a little bit cumbersome to verify the hypotheses of the theorem. In this section, the symbol K will again denote the operator defined in (3.10).

We choose

$$\mathfrak{F} \equiv \mathcal{F}_0,$$

which is simply the set of bounded smooth functions with all their derivatives bounded. It is trivial to check that \mathfrak{F} is an algebra of closed operators. Moreover, they are all self-adjoint. We also define $\mathcal{D} \equiv C_0^\infty(\mathcal{X})$.

In this section, we will first construct a set \mathcal{A} according to the rules explained in Sect. 5. Then we will check that Hypotheses 1–3 are indeed satisfied, so we will be able to apply Theorem 5.5. This will prove Proposition 3.8 almost immediately.

Before we start this program, we write down once again the definition of X_0 , as it will be used repeatedly throughout this section:

$$\begin{aligned} X_0 = & -r_L \partial_{p_0} + b_L(r_L - \lambda_L^2 q_0) \partial_{r_L} - r_R \partial_{p_N} + b_R(r_R - \lambda_R^2 q_N) \partial_{r_R} \\ & - \sum_{i=0}^N (p_i \partial_{q_i} - V_1'(q_i) \partial_{p_i}) + \sum_{i=1}^N V_2'(\tilde{q}_i) (\partial_{p_i} - \partial_{p_{i-1}}) - \alpha_K. \end{aligned}$$

6.1. Definition of \mathcal{A} . We choose an exponent $\alpha < -3/2 - \ell/(2m)$ and we let g_0 be the operator of multiplication by G^α . It is clear that g_0 is self-adjoint and positive. Moreover, we recall that

$$[X_0, G^\alpha] = \alpha G^\alpha G^{-1} [X_0, G] \in \mathcal{F}_0,$$

and so we have $[X_0, g_0] \in \mathfrak{F}$. The set \mathcal{A}_0 is defined as

$$\mathcal{A}_0 = \{c_L \partial_{r_L}, c_R \partial_{r_R}, G^\alpha X_0\} \cup \bar{\mathcal{A}},$$

with

$$\bar{\mathcal{A}} = \{a_L(r_L - \lambda_L^2 q_0), a_R(r_R - \lambda_R^2 q_N)\}.$$

Before we define the sets \mathcal{A}_i , we need a few functions. Let $i > 0$ be a natural number. The functions $V_L^{(i)}$ and $V_R^{(i)}$ are defined respectively by

$$\begin{aligned} V_L^{(i)}(\tilde{q}) &= V_2''(\tilde{q}_i) V_2''(\tilde{q}_{i-1}) \cdots V_2''(\tilde{q}_1), \\ V_R^{(i)}(\tilde{q}) &= V_2''(\tilde{q}_{N+1-i}) \cdots V_2''(\tilde{q}_{N-1}) V_2''(\tilde{q}_N). \end{aligned}$$

It is useful to notice that

$$\partial_{q_j} V_L^{(i)}(\tilde{q}) = \begin{cases} 0, & \text{if } j > i \text{ ,} \\ (V_2'''(\tilde{q}_1) V_2''(\tilde{q}_1)^{-1}) V_L^{(i)}(\tilde{q}), & \text{if } j = 0 \text{ ,} \\ (V_2'''(\tilde{q}_i) V_2''(\tilde{q}_i)^{-1}) V_L^{(i)}(\tilde{q}), & \text{if } j = i \text{ ,} \\ (V_2'''(\tilde{q}_i) V_2''(\tilde{q}_i)^{-1} - V_2'''(\tilde{q}_{i+1}) V_2''(\tilde{q}_{i+1})^{-1}) V_L^{(i)}(\tilde{q}), & \text{otherwise.} \end{cases} \tag{6.1}$$

There are symmetric relations for the derivatives of $V_R^{(i)}$. At this point, we use Assumption 3 to write

$$\partial_{q_j} V_L^{(i)}(\tilde{q}) = f_{ij}(\tilde{q}) V_L^{(i)}(\tilde{q}), \quad f_{ij} \in \mathcal{F}_{2m-2+\ell}. \quad (6.2)$$

This implies

$$[G^\alpha X_0, V_L^{(i)}(\tilde{q})] = G^\alpha \sum_{j=0}^N p_j f_{ij} V_L^{(i)}(\tilde{q}) = f_i V_L^{(i)}(\tilde{q}), \quad f_i \in \mathfrak{F}, \quad (6.3)$$

because of Proposition 4.1 and by the choice $\alpha < -3/2 - \ell/(2m)$. Moreover, we notice that

$$G^{2i\alpha} V_R^{(i)} \in \mathfrak{F},$$

still because of Proposition 4.1. One more thing we have to remember is (4.10), which implies for example that there exists a function $f_0 \in \mathfrak{F}$ such that

$$[G^\alpha X_0, G^\beta] = \beta f_0 G^\beta, \quad \text{for any } \beta \in \mathbf{R}.$$

We are now ready to complete the construction of \mathcal{A} .

6.1.1. Definition of \mathcal{A}_1 and \mathcal{A}_2 . We verify that in the case of our model, we can find functions f_B and f_{X_B} in (5.9) such that

$$\begin{aligned} \mathcal{A}_1 \setminus \mathcal{A}_0 &= \{G^\alpha \partial_{p_0}, G^\alpha \partial_{p_N}\}, \\ \mathcal{A}_2 \setminus \mathcal{A}_1 &= \{G^{2\alpha} \partial_{q_0}, G^{2\alpha} \partial_{q_N}\}. \end{aligned}$$

Considering the elements of \mathcal{A}_1 , we see that it is indeed possible to write

$$G^\alpha \partial_{p_0} = c_L^{-1} [G^\alpha X_0, c_L \partial_{r_L}] - G^{-1} (\partial_{r_L} G) G^\alpha X_0 + G^\alpha b_L \partial_{r_L},$$

and a similar relation concerning $G^\alpha \partial_{p_N}$. The operators $G^{-1} (\partial_{r_L} G)$ and $G^\alpha b_L c_L^{-1}$ belong to \mathfrak{F} , so we succeeded to construct \mathcal{A}_1 according to (5.9).

Let us now focus on the elements of \mathcal{A}_2 . We can write

$$G^{2\alpha} \partial_{q_0} = [G^\alpha X_0, G^\alpha \partial_{p_0}] - G^{\alpha-1} (\partial_{p_0} G) G^\alpha X_0 - \alpha f_0 G^\alpha \partial_{p_0},$$

and an equivalent expression at the other end of the chain. Since $G^{\alpha-1} (\partial_{p_0} G) \in \mathfrak{F}$ and $f_0 \in \mathfrak{F}$, we succeeded to construct \mathcal{A}_2 according to (5.9).

6.1.2. Definition of \mathcal{A}_{2i-1} and \mathcal{A}_{2i} . For $i \geq 1$, these sets are defined by

$$\begin{aligned} \mathcal{A}_{2i-1} \setminus \mathcal{A}_{2i-2} &= \{G^{(2i+1)\alpha} V_L^{(i)} \partial_{p_i}, G^{(2i+1)\alpha} V_R^{(i)} \partial_{p_{N-i}}\}, \\ \mathcal{A}_{2i} \setminus \mathcal{A}_{2i-1} &= \{G^{(2i+2)\alpha} V_L^{(i)} \partial_{q_i}, G^{(2i+2)\alpha} V_R^{(i)} \partial_{q_{N-i}}\}. \end{aligned}$$

We repeat this construction until $i = N - 1$, i.e. we do not stop at the middle of the chain, but we go on until we reach the other end. We want to check that these sets were constructed according to (5.9). In fact, we will see that any element A of $\mathcal{A}_j \setminus \mathcal{A}_{j-1}$ with $j \geq 2$ can be written as

$$A = [G^\alpha X_0, B] + D, \quad B \in \mathcal{A}_{j-1}, \quad D \in \mathcal{Y}^1(\mathcal{A}_{j-1}). \quad (6.4)$$

We will verify this only for $2 \leq i \leq N - 2$. We let the reader verify that (6.4) is also valid for the remaining sets.

Let us first take $j = 2i - 1$ and $A = G^{(2i+1)\alpha} V_L^{(i)} \partial_{p_i}$. We choose

$$B = G^{2i\alpha} V_L^{(i-1)} \partial_{q_{i-1}} \in \mathcal{A}_{j-1}$$

and write

$$\begin{aligned} [G^\alpha X_0, B] &= f_{i-1} G^{2i\alpha} V_L^{(i-1)} \partial_{q_{i-1}} - G^{2i\alpha} V_L^{(i-1)} G^{-1} (\partial_{q_{i-1}} G) G^\alpha X_0 \\ &\quad + 2i\alpha f_0 B + G^{(2i+1)\alpha} V_L^{(i-1)} [X_0, \partial_{q_{i-1}}]. \end{aligned}$$

The first three terms belong to $\mathcal{Y}^1(\mathcal{A}_{2i-2})$ and can thus be absorbed into D . The last term can be written as

$$\begin{aligned} G^{(2i+1)\alpha} V_L^{(i-1)} [X_0, \partial_{q_{i-1}}] &= \\ &G^{2\alpha} (V_1''(q_{i-1}) + V_2''(\tilde{q}_i) + V_2''(\tilde{q}_{i-1})) G^{(2i-1)\alpha} V_L^{(i-1)} \partial_{p_{i-1}} \\ &\quad + G^{4\alpha} V_2''(\tilde{q}) V_2''(\tilde{q}) G^{(2i-3)\alpha} V_L^{(i-2)} \partial_{p_{i-2}} + G^{(2i+1)\alpha} V_L^{(i)} \partial_{p_i}. \end{aligned}$$

The first two terms also belong to $\mathcal{Y}^1(\mathcal{A}_{2i-2})$, so they can be absorbed into D as well. The remaining term is

$$G^{(2i+1)\alpha} V_L^{(i)} \partial_{p_i} = A,$$

thus we have verified that A can be written as in (6.4). The procedure to get the symmetric term from the other end of the chain is similar.

We take now $j = 2i$ and $A = G^{(2i+2)\alpha} V_L^{(i)} \partial_{q_i}$. We choose

$$B = G^{(2i+1)\alpha} V_L^{(i)} \partial_{p_i} \in \mathcal{Y}^1(\mathcal{A}_{j-1})$$

and write

$$\begin{aligned} [G^\alpha X_0, B] &= f_i G^{(2i+1)\alpha} V_L^{(i)} \partial_{p_i} + G^{(2i+1)\alpha} V_L^{(i)} G^{-1} (\partial_{p_i} G) G^\alpha X_0 \\ &\quad + (2i + 1)\alpha f_0 B + G^{(2i+2)\alpha} V_L^{(i)} \partial_{q_i}. \end{aligned}$$

The first three terms belong to $\mathcal{Y}^1(\mathcal{A}_{2i-1})$ and can be absorbed into D , so we verified that every element of \mathcal{A} can indeed be written as in (5.9).

6.2. Verification of the hypotheses and proof. In order to be able to apply Theorem 5.5, we verify the hypotheses 1–3.

Verification of Hypothesis 2. We want to check that $A \in \mathcal{A}_j$ implies $A^* \in \mathcal{Y}^1(\mathcal{A}_j)$. By Proposition 4.1, we can easily verify that $\mathcal{A} \setminus \bar{\mathcal{A}} \subset \mathcal{L}_0$. But we know that

$$A \in \mathcal{L}_0 \Rightarrow A^* = -A + g,$$

and so Hypothesis 2 holds for $\mathcal{A} \setminus \bar{\mathcal{A}}$. The elements of $\bar{\mathcal{A}}$ being self-adjoint, Hypothesis 2 holds trivially.

Verification of Hypothesis 3. The operator Λ^2 can be written as

$$\Lambda^2 = - \sum_{i,j} \partial_i a_{ij}(x) \partial_j + V(x).$$

It is well-known that if a_{ij} and V are sufficiently nice, such operators are essentially self-adjoint on $C_0^\infty(\mathcal{X})$ (see e.g. [Agm82, Thm. 3.2]).

Verification of Hypothesis 1. Let us define $\mathfrak{L}_0 \subset \mathfrak{L}$ as the set of first-order differential operators with coefficients in \mathcal{F}_0 . We first verify that

$$A \in \mathcal{A}, f \in \mathfrak{F} \quad \Rightarrow \quad [A, f] \in \mathfrak{F}.$$

This is trivial, noticing that $\mathcal{A} \subset \mathfrak{L}_0 \cup \bar{\mathcal{A}}$ and $[\mathfrak{L}_0, \mathfrak{F}] = [\bar{\mathcal{A}}, \mathfrak{F}] = \{0\}$.

We now verify that

$$A \in \mathcal{A} \quad \Rightarrow \quad A^* \in \mathcal{Y}^1(\mathcal{A}).$$

This is also trivial, because $A \in \mathfrak{L}_0 \Rightarrow A^* = -A + g$, with $g \in \mathcal{F}_0$. Moreover, the elements of $\bar{\mathcal{A}}$ are self-adjoint.

Finally, we want to verify that

$$A, B \in \mathcal{A} \quad \Rightarrow \quad [A, B] \in \mathcal{Y}^1(\mathcal{A}).$$

This is a little bit longer to verify.

Concerning the commutators of the elements of $\bar{\mathcal{A}}$ with the other elements of \mathcal{A} , the statement follows easily from the fact that if $F : \mathbf{R}^n \rightarrow R$ is linear and $A \in \mathfrak{L}_0$, then $[A, F] \in \mathcal{F}_0 \equiv \mathfrak{F}$. Moreover, the commutator between two multiplication operators vanishes.

Concerning the commutators between the ∂_r and the other elements, we notice that they commute with the functions $V_L^{(i)}(\tilde{q})$ and $V_R^{(i)}(\tilde{q})$. Moreover, we have for example

$$[\partial_{r_L}, G^\gamma] = \gamma(G^{-1}[\partial_{r_L}, G])G^\gamma, \quad \text{if } \gamma \in \mathbf{R},$$

and $G^{-1}[\partial_{r_L}, G]$ belongs to \mathfrak{F} . It is straightforward to verify that this implies the desired statement.

Concerning the commutators of $G^\alpha X_0$ with the other elements of \mathcal{A} , the statement has already been verified by the construction of \mathcal{A} for every operator, but those in $\mathcal{A}_{2N-2} \setminus \mathcal{A}_{2N-3}$. These operators are of the form

$$A = G^{2N\alpha} V_R^{(N-1)} \partial_{q_1},$$

and a similar term at the other end of the chain. We can make a computation very similar to the one we made when we constructed \mathcal{A}_{2i-1} , to show that

$$[G^\alpha X_0, A] = G^{(2N+1)\alpha} V_R^{(N)} \partial_{p_0} + C, \quad C \in \mathcal{Y}^1(\mathcal{A}).$$

But $G^{2N\alpha} V_R^{(N)} \in \mathfrak{F}$, so $[G^\alpha X_0, A] \in \mathcal{Y}^1(\mathcal{A})$. It remains therefore only to verify the statement for commutators between elements of $\mathcal{A} \setminus \mathcal{A}_0$. We can divide these commutators in three classes.

Both operators contain a ∂_p . We notice that these operators can all be written in the form $G^{\alpha_i} W_i(q) \partial_{p_i}$. The commutator between two such elements is given by

$$\begin{aligned} [G^{\alpha_i} W_i(q) \partial_{p_i}, G^{\alpha_j} W_j(q) \partial_{p_j}] &= G^{-1}(\partial_{p_i} G) G^{\alpha_i} W_i(q) G^{\alpha_j} W_j(q) \partial_{p_j} \\ &\quad - G^{-1}(\partial_{p_j} G) G^{\alpha_j} W_j(q) G^{\alpha_i} W_i(q) \partial_{p_i}. \end{aligned}$$

Both terms belong to $\mathcal{Y}^1(\mathcal{A})$, because $G^{-1}(\partial_p G) \in \mathfrak{F}$.

One operator contains a ∂_p , one contains a ∂_q . Let us compute the commutator between $G^{(2i+2)\alpha} V_L^{(i)} \partial_{q_i}$ and $G^{(2j+1)\alpha} V_L^{(j)} \partial_{p_j}$. We have

$$\begin{aligned}
 [G^{(2i+2)\alpha} V_L^{(i)} \partial_{q_i}, G^{(2j+1)\alpha} V_L^{(j)} \partial_{p_j}] &= G^{(2i+2)\alpha} V_L^{(i)} (\partial_{q_i} G) G^{-1} G^{(2j+1)\alpha} V_L^{(j)} \partial_{p_j} \\
 &\quad + G^{(2i+1)\alpha} V_L^{(i)} (\partial_{p_i} G) G^{-1} G^{(2i+2)\alpha} V_L^{(i)} \partial_{q_i} \\
 &\quad + G^{2i\alpha} V_L^{(i)} G^{2\alpha} f_{ij} G^{(2j+1)\alpha} V_L^{(j)} \partial_{p_j}.
 \end{aligned}$$

All those terms belong to $\mathcal{Y}^1(\mathcal{A})$. The computation is similar if we take for example

$$G^{(2j+1)\alpha} V_R^{(j)} \partial_{p_{N-j}}$$

instead of $G^{(2j+1)\alpha} V_L^{(j)} \partial_{p_j}$.

Both operators contain a ∂_q . The computation is similar to the preceding case and is left to the reader.

It is now easy to give the

Proof of Proposition 3.8. We have just verified that the hypotheses of Theorem 5.5 are satisfied. We apply it, so we have the estimate

$$\|\tilde{\Delta}^\varepsilon f\| \leq C(\|Kf\| + \|f\|),$$

where $\tilde{\Delta}$ is given by

$$\tilde{\Delta} = 1 + \sum_{A \in \mathcal{A}} A^* A.$$

It is easy to see that $\tilde{\Delta}$ has exactly the form (3.12). This completes the proof of Proposition 3.8. \square

7. Proof of Theorem 3.6

It is now possible to prove that the operator K has compact resolvent, which is one of the main results of this paper. Before we start the proof itself, we need two preliminary results. The first one states

Lemma 7.1. *Let $\tilde{\Delta}$ be the closure in $L^2(\mathbf{R}^n)$ of the operator acting on $C_0^\infty(\mathbf{R}^n)$ as*

$$\tilde{\Delta} = \sum_{i=1}^{\tilde{N}} L_i^* L_i + a_0,$$

where the L_i are smooth vector fields with bounded coefficients spanning \mathbf{R}^n at every point and a_0 is a smooth positive function.

Let $V : \mathbf{R}^n \rightarrow \mathbf{R}^n$ be a continuous function such that for every constant $C > 0$, there exists a compact $K_C \subset \mathbf{R}^n$ with the property that $V(x) > C$ for every $x \in \mathbf{R}^n \setminus K_C$. We moreover assume that $V(x) \geq 1$. Define the operator H as the closure in $L^2(\mathbf{R}^n)$ of the operator acting on $f \in C_0^\infty(\mathbf{R}^n)$ as

$$(Hf)(x) = (\tilde{\Delta}f)(x) + V(x)f(x).$$

Then the operator H is self-adjoint.

Suppose V and the L_i are such that the function

$$2a_0V + \sum_{i=1}^{\tilde{N}} \left((L_i^* + L_i)[L_i, V] - [L_i, [L_i, V]] \right) \tag{7.1}$$

is bounded. We then have the estimate

$$\langle f, H^\varepsilon f \rangle \leq \langle f, \tilde{\Delta}^\varepsilon f \rangle + \langle f, V^\varepsilon f \rangle + C \langle f, H^{\varepsilon-1} f \rangle, \quad 0 < \varepsilon < 1, \tag{7.2}$$

which holds for any $f \in C_0^\infty(\mathbf{R}^n)$.

Proof. The result concerning the self-adjointness of H and of $\tilde{\Delta}$ is classical, we will not prove it here. The interested reader can find a proof in [Agm82, Thm. 3.2].

We use the fact that if T is a strictly positive self-adjoint operator and $\alpha = 1 - \varepsilon \in (0, 1)$, we can write

$$T^{-\alpha} = C_\alpha \int_0^\infty z^{-\alpha} (z + T)^{-1} dz, \quad C_\alpha = \frac{\sin(\pi\alpha)}{\pi},$$

and thus

$$T^\varepsilon = C_\alpha \int_0^\infty z^{\varepsilon-1} \frac{T}{z + T} dz.$$

Moreover, a core of T is again a core of T^ε , so (7.2) makes sense. For a proof of these statements, see [Kat80, §V.3]. This allows us to write inequality (7.2) as

$$\begin{aligned} \int_0^\infty z^{\varepsilon-1} \left\langle f, \frac{H}{z + H} f \right\rangle dz &\leq \\ &\int_0^\infty z^{\varepsilon-1} \left\langle f, \frac{\tilde{\Delta}}{z + \tilde{\Delta}} f \right\rangle dz + \int_0^\infty z^{\varepsilon-1} \left\langle f, \frac{V}{z + V} f \right\rangle dz \\ &+ C \int_0^\infty z^{\varepsilon-1} \left\langle f, \frac{1}{z + H} f \right\rangle dz. \end{aligned} \tag{7.3}$$

In order to prove (7.3), let us first show that the operator $\tilde{\Delta}V + V\tilde{\Delta}$ is lower bounded. This is an immediate consequence of (7.1) and the equality

$$L_i^* L_i V + V L_i^* L_i = 2L_i^* V L_i + (L_i + L_i^*)[L_i, V] - [L_i, [L_i, V]],$$

which is easily verified, using the fact that $L_i + L_i^*$ is simply a function. Therefore, there exists a constant $C > 0$ such that

$$\langle g, (\tilde{\Delta}V + V\tilde{\Delta})g \rangle + C \langle g, g \rangle \geq 0, \quad \forall g \in C_0^\infty(\mathbf{R}^n).$$

Since $H \geq 1$ in the sense of quadratic forms, we find

$$\langle g, (\tilde{\Delta}V + V\tilde{\Delta})g \rangle + C \langle g, (z + H)g \rangle \geq 0,$$

which holds for every $z \geq 0$. Since $\tilde{\Delta}$ and V are positive self-adjoint operators, this immediately implies

$$\left\langle g, V \frac{\tilde{\Delta}}{z + \tilde{\Delta}} V g \right\rangle + \left\langle g, \tilde{\Delta} \frac{V}{z + V} \tilde{\Delta} g \right\rangle + \langle g, (\tilde{\Delta}V + V\tilde{\Delta})g \rangle + C \langle g, (z + H)g \rangle \geq 0. \tag{7.4}$$

We can easily check the following identities:

$$\begin{aligned} V\tilde{\Delta}(z + \tilde{\Delta})^{-1}V + \tilde{\Delta}V &= (z + H)\tilde{\Delta}(z + \tilde{\Delta})^{-1}V, \\ \tilde{\Delta}V(z + V)^{-1}\tilde{\Delta} + V\tilde{\Delta} &= (z + H)\tilde{\Delta}(z + V)^{-1}\tilde{\Delta}. \end{aligned}$$

Inserting this in (7.4), we get

$$\langle g, (z + H)(\tilde{\Delta}(z + \tilde{\Delta})^{-1}V + V(z + V)^{-1}\tilde{\Delta})g \rangle + C\langle g, (z + H)g \rangle \geq 0,$$

and thus

$$\begin{aligned} \langle g, (z + H)Hg \rangle &\leq \langle g, (z + H)(H + \tilde{\Delta}(z + \tilde{\Delta})^{-1}V \\ &\quad + V(z + V)^{-1}\tilde{\Delta})g \rangle + C\langle g, (z + H)g \rangle. \end{aligned}$$

We can check the equalities

$$\begin{aligned} \tilde{\Delta}(z + \tilde{\Delta})^{-1}V + \tilde{\Delta} &= \tilde{\Delta}(z + \tilde{\Delta})^{-1}(z + H), \\ V(z + V)^{-1}\tilde{\Delta} + V &= V(z + V)^{-1}(z + H), \end{aligned}$$

which allow us to write

$$\begin{aligned} \langle g, (z + H)Hg \rangle &\leq \langle (z + H)g, (\tilde{\Delta}(z + \tilde{\Delta})^{-1} + V(z + V)^{-1})(z + H)g \rangle \\ &\quad + C\langle g, (z + H)g \rangle. \end{aligned}$$

Let us define $f \equiv (z + H)g$. This immediately yields the estimate

$$\left\langle f, \frac{H}{z + H}f \right\rangle \leq \left\langle f, \frac{\tilde{\Delta}}{z + \tilde{\Delta}}f \right\rangle + \left\langle f, \frac{V}{z + V}f \right\rangle + C\left\langle f, \frac{1}{z + H}f \right\rangle, \quad (7.5)$$

which holds for any f in $\mathcal{W} \equiv (z + H)\mathcal{C}_0^\infty(\mathbf{R}^n)$. But we know that $\mathcal{C}_0^\infty(\mathbf{R}^n)$ is a core for H , therefore \mathcal{W} is dense in $L^2(\mathbf{R}^n)$. Since the operators appearing in (7.5) are all bounded, the inequality (7.5) holds for every $f \in L^2(\mathbf{R}^n)$ and thus in particular also for $f \in \mathcal{C}_0^\infty(\mathbf{R}^n)$. This implies the wanted estimate (7.3). \square

The second result we want to use is

Proposition 7.2. *Let $\tilde{\Delta}$, V and H be as in Lemma 7.1. Then H has compact resolvent.*

Proof. We know that $\tilde{\Delta}$ is a positive self-adjoint operator, so

$$T = (\tilde{\Delta} + 1)^{-1}$$

exists and $\|T\| \leq 1$. The proof of compactness is a modification of the standard proof of the same theorem with $\tilde{\Delta}$ replaced by the true Laplacian Δ , which can be found e.g. in [Agm82]. It is based on the fact that if χ is a function with compact support, then the multiplication operator χ is relatively compact with respect to $\tilde{\Delta}$. We want to prove that χT is a compact operator, i.e. that the closure of

$$Y = \{\chi T f \mid f \in \mathcal{C}_0^\infty(\mathbf{R}^n) \text{ and } \|f\| \leq 1\}$$

is compact.

Let us define $\mathcal{K} = \text{supp } \chi$. By hypothesis, \mathcal{K} is compact. Moreover, we have $Y \subset \mathcal{C}_0^\infty(\mathcal{K})$. It is well-known that if \mathcal{K} is a compact domain of \mathbf{R}^n , then the set

$$\{u \in \mathcal{C}_0^\infty(\mathcal{K}) \mid \|u\| \leq 1; \langle u, \Delta u \rangle \leq 1\}$$

is compact (see e.g. [RS80, Thm. XIII.73]). This implies that Y is compact if we are able to prove that there are strictly positive constants ε , c_1 and c_2 such that $u \in Y$ implies

$$\|u\| \leq c_1 \quad \text{and} \quad \langle u, \Delta u \rangle \leq c_2.$$

We take any element u in Y and write it as $u = \chi T f$. We have

$$\|u\| \leq \|\chi\|_\infty \|T\| \|f\| \leq c_1.$$

Recall that we assumed the vector fields L_i appearing in the construction of $\tilde{\Delta}$ span \mathbf{R}^n at any point and that a_0 is a strictly positive function. Together with the compactness of the support of u , this implies that there are constants C and k_1 such that

$$\begin{aligned} |\langle u, \Delta u \rangle| &\leq C |\langle u, \tilde{\Delta} u \rangle| = C |\langle u, \tilde{\Delta} \chi T f \rangle| \leq C \|u\| \|\tilde{\Delta} \chi T f\| \\ &\leq C \|\chi \tilde{\Delta} T f\| + C \|[\tilde{\Delta}, \chi] T f\| \leq k_1 + C \|[\tilde{\Delta}, \chi] T f\|, \end{aligned} \quad (7.6)$$

where the last inequality is a consequence of $T = (1 + \tilde{\Delta})^{-1}$. We therefore only need to bound the term containing the commutator of $\tilde{\Delta}$ and χ . Explicit calculation yields

$$[\tilde{\Delta}, \chi] = \sum_{i=1}^{\tilde{N}} \left(-2[L_i, \chi]L_i + [L_i, [L_i, \chi]] + (L_i^* + L_i)[L_i, \chi] \right) \equiv \sum_{i=1}^{\tilde{N}} \eta_i L_i + \eta_0,$$

where the η_i are bounded functions with $\text{supp } \eta_i \subset \mathcal{K}$. So the only terms that remain to be bounded are of the form $\|\eta_i L_i T f\|$. As η_i is bounded, it is enough to bound $\|L_i T f\|$. We have

$$\|L_i T f\|^2 = \langle T f, L_i^* L_i T f \rangle \leq \langle T f, \tilde{\Delta} T f \rangle \leq \|f\|^2. \quad (7.7)$$

This completes the proof of the statement about the relative compactness of χ .

This implies that we can add to H any function with compact support without changing its essential spectrum (see [RS80, Thm. XIII.14]). But the assumption we made concerning V and the positivity of $\tilde{\Delta}$ imply that for any constant C , we can raise the spectrum of $H + \chi$ above C by taking for χ a smooth function satisfying

$$\chi(x) = \begin{cases} C & x \in \mathcal{K}_C, \\ 0 & d(x, \mathcal{K}_C) > 1. \end{cases}$$

Therefore, the essential spectrum of H is empty and thus H has compact resolvent. \square

It is now easy to give the

Proof of Theorem 3.6. By Proposition 3.7 and 3.8, we can choose a constant ε small enough to have, for every $f \in \mathcal{C}_0^\infty(\mathcal{X})$, the estimate

$$\|\tilde{\Delta}^\varepsilon f\| \leq C(\|Kf\| + \|f\|) \quad \text{and} \quad \|G^\varepsilon f\| \leq C(\|Kf\| + \|f\|).$$

We moreover define

$$H \equiv \tilde{\Delta} + G. \quad (7.8)$$

By the proof of Proposition 3.8, we see that the assumptions of Lemma 7.1 are satisfied. We can thus write

$$\begin{aligned} \|H^\varepsilon f\|^2 &= \langle f, H^{2\varepsilon} f \rangle = \langle f, (\tilde{\Delta} + G)^{2\varepsilon} f \rangle \leq \langle f, \tilde{\Delta}^{2\varepsilon} f \rangle + \langle f, G^{2\varepsilon} f \rangle + C\|f\|^2 \\ &\leq \|\tilde{\Delta}^\varepsilon f\|^2 + \|G^\varepsilon f\|^2 + C\|f\|^2 \leq C(\|Kf\| + \|f\|)^2. \end{aligned}$$

Because G is confining, we can apply Proposition 7.2 to see that H , and therefore also H^ε , have compact resolvent. Therefore Corollary 5.9 applies, showing that K has compact resolvent. \square

Remark 7.3. The proof still works under slightly weaker assumptions. The coupling between the ends of the chain and the heat baths does not have to be of the dipolar type. It is enough for example that F_L and F_R belong to some \mathcal{F}_β with $\beta < n$. Moreover, the potentials V_1 and V_2 can be different for each particle. We only have to impose that Assumptions 1–3 can be satisfied for every particle with the same constants ℓ , m and n .

Remark 7.4. Throughout this paper, we restricted ourselves to the one-dimensional case, i.e. each particle had only one degree of freedom. It is not very hard to generalize the results of this paper to the d -dimensional case. It is straightforward to generalize Assumptions 1 and 2, where V' is now a vector. In Assumption 3, the inverse of V_2'' has to be read as the inverse matrix. A matrix or vector-valued function is said to belong to \mathcal{F}_β if each of its components belong to \mathcal{F}_β .

The only point that could cause some trouble is the expression (6.1), because the $V_2''(\tilde{q}_j)$ are now matrices which do not commute, so the expression for $\partial_{q_j} V_L^{(i)}(\tilde{q})$ will show terms of the form

$$V_2''(\tilde{q}_i) V_2''(\tilde{q}_{i-1}) \cdots \cdots V_2'''(\tilde{q}_j) \cdots \cdots V_2''(\tilde{q}_1),$$

where V_2''' is a trilinear form. Such a term can be written as

$$V_2''(\tilde{q}_i) V_2''(\tilde{q}_{i-1}) \cdots \cdots V_2'''(\tilde{q}_j) V_2''(\tilde{q}_{j+1})^{-1} \cdots \cdots V_2''(\tilde{q}_i)^{-1} V_L^{(i)}.$$

If we want to get expressions similar to (6.2) and (6.3), we have to make $|\alpha|$ very big (of the order of N), but this is not a problem.

Remark 7.5. One important assumption was that $m > n$, in other words, the interparticle coupling is stronger at infinity than the single particle potential. If this is not satisfied, our proof does not work. There could be some physical reason behind this. If a stationary state exists, this means that even if the chain is in a state of very high energy, the mean time to reach a region with low energy is finite (see e.g. [Has80]). But if $m < n$, the relative strength of the coupling versus the one-body potential goes to zero at high energy. The consequence is that there is almost no energy transmitted between particles. Since the only points where dissipation occurs are the ends of the chain, we see that the higher the energy of the chain is, the slower this energy will be dissipated. Probably this is not sufficient to destroy the existence of a stationary state, but it could explain why the proof does not work in this situation. It is even possible that this phenomenon destroys the compactness of the resolvent of K .

8. The Invariant Measure

This section is devoted to the proof of Theorem 3.9. Throughout this section, we denote by \mathcal{T}^t the semigroup generated by the system of stochastic differential equations (3.7). We also assume that Assumptions 1–3 are satisfied, so Propositions 3.7 and 3.8 hold, as well as Theorem 3.6. The proof of Theorem 3.9 is divided into three separate propositions, showing respectively the following properties of the invariant measure μ :

- (i) Existence and smoothness.
- (ii) Decay properties.
- (iii) Uniqueness and strict positivity.

Proposition 8.1. *If Assumptions 1–3 are satisfied, the Markov process given by (3.7) possesses an invariant measure μ . It has a density h , which is a C^∞ function on \mathbf{R}^{2N+4} .*

Proof. By Theorem 3.6, we know that K has compact resolvent. This implies also the compactness of the resolvent of $L_{\mathcal{H}}$ and thus of L_0 . Since G grows algebraically at infinity, we see that the constant function 1 belongs to \mathcal{H}_0 . Moreover, we notice that $L_0 1 = 0$, thus the operator L_0 has an eigenvalue 0, which is isolated because of the compactness of its resolvent. This in turn implies that L_0^* also has an isolated eigenvalue 0. We denote the corresponding eigenvector by g and normalize it so that $\langle g, 1 \rangle_{\mathcal{H}_0} = 1$. Since L_0^* is hypoelliptic, g must be C^∞ .

Assume first that $g \geq 0$. We then define

$$h(p, q, r) = Z_0^{-1} g(p, q, r) e^{-2\beta_0 G}, \tag{8.1}$$

where Z_0 is the normalization constant appearing in the definition of \mathcal{H}_0 . Set $\mu(dx) = h(x) dx$; we want to check that μ is the invariant measure we are looking for. Notice that $\mu(dx)$ is a probability measure because

$$\int \mu(dx) = Z_0^{-1} \int e^{-2\beta_0 G(x)} g(x) dx = \langle g, 1 \rangle_{\mathcal{H}_0} = 1.$$

Let A be a Borel set of \mathbf{R}^{2N+4} . Then the characteristic function χ_A of A belongs to \mathcal{H}_0 . We have

$$\begin{aligned} ((\mathcal{T}^t)^* \mu)(A) &= \int (\mathcal{T}^t \chi_A)(x) \mu(dx) = Z_0^{-1} \int e^{-2\beta_0 G(x)} g(x) (\mathcal{T}^t \chi_A)(x) dx \\ &= Z_0^{-1} \int e^{-2\beta_0 G(x)} ((\mathcal{T}_0^t)^* g)(x) \chi_A(x) dx = \mu(A), \end{aligned}$$

thus μ is an invariant measure for the Markov process defined by (3.7).

The argument showing that it was indeed justified to assume g positive can be taken over from [EPR99a, Prop. 3.6]. \square

We next turn to the decay properties of the invariant measure h . We first introduce a convenient family of Hilbert spaces.

Definition 8.2. *Choose $\gamma \in \mathbf{R}$. We define the Hilbert space $W^{(\gamma)}$ as*

$$W^{(\gamma)} \equiv L^2(\mathcal{X}, G^{2\gamma}(x) dx) = \mathcal{D}(G^\gamma).$$

We will denote by $\langle \cdot, \cdot \rangle_{(\gamma)}$ and $\| \cdot \|_{(\gamma)}$ the corresponding scalar product and norm. We also define

$$W^{(\infty)} \equiv \bigcap_{\gamma > 0} W^{(\gamma)},$$

which is the set of all functions that decay at infinity faster than any polynomial.

We already know that h is a C^∞ function, so we want to show that it is possible to write

$$h(p, q, r) = \tilde{h}(p, q, r)e^{-\beta_0 G(p, q, r)}, \quad \tilde{h} \in W^{(\infty)}.$$

The function \tilde{h} being an eigenfunction of the operator K , the decay properties of the invariant measure are a consequence of the following result.

Proposition 8.3. *The eigenfunctions of K and K^* belong to $C^\infty(\mathcal{X}) \cap W^{(\infty)}$.*

We will show Proposition 8.3 only for the eigenfunctions of K . It is a simple exercise left to the reader to retrace the proof for the eigenfunctions of K^* . We already know that K and K^* are hypoelliptic, so their eigenfunctions belong to $C^\infty(\mathcal{X})$. It remains to be proven that they also belong to $W^{(\infty)}$.

To prove the proposition, we will show the implication

$$f \in W^{(\gamma)} \quad \text{and} \quad Kf \in W^{(\gamma)} \quad \Rightarrow \quad f \in W^{(\gamma+\varepsilon)}, \quad (8.2)$$

which immediately implies that the eigenvectors of K belong to $W^{(\infty)}$. For this purpose, we introduce the family of operators K_γ defined by

$$\begin{aligned} K_\gamma : \mathcal{D}(K_\gamma) &\rightarrow W^{(\gamma)} \\ f &\mapsto Kf, \end{aligned}$$

where $\mathcal{D}(K_\gamma)$ is given by

$$\mathcal{D}(K_\gamma) = \{f \in W^{(\gamma)} \mid Kf \in W^{(\gamma)}\}.$$

The expression Kf has to be understood in the sense of distributions.

We have the following preliminary result.

Lemma 8.4. $C_0^\infty(\mathcal{X})$ is a core for K_γ .

Proof. The proof uses the tools developed in Appendix B and is postponed to Appendix C. \square

The key lemma for the proof of Proposition 8.3 is the following.

Lemma 8.5. *There are an $\varepsilon > 0$ and constants $C_\gamma > 0$ such that for every $\gamma > 0$ and every $u \in \mathcal{D}(K_\gamma)$, the relation*

$$\|G^\varepsilon u\|_{(\gamma)}^2 \leq C_\gamma (\|K_\gamma u\|_{(\gamma)}^2 + \|u\|_{(\gamma)}^2) \quad (8.3)$$

holds.

Proof. Since we know that $\mathcal{C}_0^\infty(\mathcal{X})$ is a core for K_γ , it suffices to show (8.3) for $u \in \mathcal{C}_0^\infty(\mathcal{X})$. Let L be the first-order differential operator associated to a divergence-free vector field. Then we have for $f, g \in \mathcal{C}_0^\infty$,

$$\begin{aligned} \langle Lf, g \rangle_{(\gamma)} &= -\langle f, LG^{2\gamma}g \rangle = -\langle f, G^{2\gamma}Lg \rangle - 2\gamma\langle f, G^{2\gamma}G^{-1}(LG)g \rangle \\ &= -\langle f, Lg \rangle_{(\gamma)} - 2\gamma\langle f, G^{-1}(LG)g \rangle_{(\gamma)}. \end{aligned}$$

We write this symbolically as

$$L_\gamma^* = -L_\gamma - 2\gamma G^{-1}(LG).$$

We can use the latter equality to show that there are constants $c_\gamma^{(1)}$ and $c_\gamma^{(2)}$ such that

$$-\frac{(L_\gamma)^2 + (L_\gamma^*)^2}{2} = L_\gamma^*L_\gamma + c_\gamma^{(1)}G^{-2}(LG)^2 + c_\gamma^{(2)}G^{-1}(L^2G).$$

Using the explicit form of K , this in turn yields the useful relation

$$\begin{aligned} \operatorname{Re}\langle u, Ku \rangle_{(\gamma)} &= c_L^2 \|\partial_{r_L} u\|_{(\gamma)}^2 + c_R^2 \|\partial_{r_R} u\|_{(\gamma)}^2 \\ &\quad + a_L^2 \|(r_L - \lambda_L q_0)u\|_{(\gamma)}^2 + a_R^2 \|(r_R - \lambda_R q_N)u\|_{(\gamma)}^2 + \langle u, f_K u \rangle_{(\gamma)}, \end{aligned} \tag{8.4}$$

where f_K is some bounded function.

We now have the tools to prove the validity of (8.3). We use Proposition 3.7 to write

$$\begin{aligned} \|u\|_{(\gamma+\varepsilon)}^2 &= \|G^\varepsilon G^\gamma u\|^2 \leq C(\|K G^\gamma u\|^2 + \|G^\gamma u\|^2) \\ &\leq C(\|G^\gamma Ku\|^2 + \|[K, G^\gamma]u\|^2 + \|G^\gamma u\|^2). \end{aligned}$$

An explicit computation yields

$$[K, G^\gamma]u = G^\gamma(f_L \partial_{r_L} + f_R \partial_{r_R} + f_0)u,$$

for some smooth bounded functions f_L, f_R and f_0 . We are thus able to write

$$\|u\|_{(\gamma+\varepsilon)}^2 \leq C(\|Ku\|_{(\gamma)}^2 + \|u\|_{(\gamma)}^2 + c_L^2 \|\partial_{r_L} u\|_{(\gamma)}^2 + c_R^2 \|\partial_{r_R} u\|_{(\gamma)}^2). \tag{8.5}$$

Using (8.4), we can write

$$\begin{aligned} c_L^2 \|\partial_{r_L} u\|_{(\gamma)}^2 + c_R^2 \|\partial_{r_R} u\|_{(\gamma)}^2 &\leq |\operatorname{Re}\langle u, Ku \rangle_{(\gamma)}| + C\|u\|_{(\gamma)}^2 \\ &\leq C(\|Ku\|_{(\gamma)}^2 + \|u\|_{(\gamma)}^2). \end{aligned}$$

This, together with (8.5), completes the proof of the assertion. \square

Proof of Proposition 8.3. Lemma 8.5 immediately shows that $\mathcal{D}(K_\gamma) \subset W^{(\gamma+\varepsilon)}$ for every $\gamma > 0$. This proves the assertion (8.2).

Let f be an eigenfunction of K . We know that $f \in L^2(\mathcal{X})$ and, because it is an eigenvector of K , we have $Kf \in L^2$. Thus, by (8.2), $f \in W^{(\varepsilon)}$. Of course $Kf \in W^{(\varepsilon)}$ as well, so $f \in W^{(2\varepsilon)}$. This can be continued *ad infinitum*, and so we have $f \in W^{(\infty)}$, which is the desired result. \square

Finally, we want to show the strict positivity and the uniqueness of the invariant measure. The proof of this result will only be sketched, as it simply retraces the proof of Theorem 3.6 in [EPR99b].

Proposition 8.6. *The density h of the invariant measure μ is a strictly positive function. Moreover, the invariant measure is unique.*

Sketch of proof. The idea is to show that the control system associated with the stochastic differential equation (3.7) is strongly completely controllable. This means that, given an initial condition x_0 , a time τ and an endpoint x_τ , it is possible to find a realization of the Wiener process w such that $\xi(\tau; x_0, w) = x_\tau$. The main assumption needed to show that is that the gradient of the two-body potential is a diffeomorphism. This is ensured by Assumption 3.

The consequence is that, for every time τ , every initial condition x_0 and every open set U , the transition probability $P(\tau, x_0, U)$ is strictly positive. Because μ is invariant, we have

$$\mu(U) = \int P(t, x, U) \mu(dx) > 0.$$

This implies the strict positivity of h . Uniqueness follows from an elementary ergodicity argument. \square

A. Proof of Lemma 5.6

Throughout this appendix, we will make use of the same notations as in Sect. 5, i.e. $\mathcal{H} = \mathbb{L}^2(\mathbf{R}^n)$, $\mathcal{D} = \mathcal{C}_0^\infty(\mathbf{R}^n)$ and \mathfrak{D} is the set of differential operators with smooth coefficients.

Moreover, \mathcal{A} denotes some finite subset of \mathfrak{D} and is identified with closed operators on \mathcal{H} . The operator Λ^2 is defined as

$$\Lambda^2 \equiv 1 + \sum_{A \in \mathcal{A}} A^* A. \tag{A.1}$$

We will moreover assume that Hypotheses 1 and 3 concerning \mathcal{A} and \mathfrak{F} hold, i.e. $A, B \in \mathcal{A}$ and $f \in \mathfrak{F}$ imply

$$[A, B] \in \mathcal{Y}^1(\mathcal{A}), \quad A^* \in \mathcal{Y}^1(\mathcal{A}), \quad [A, f] \in \mathfrak{F}. \tag{A.2}$$

In order to prepare the proof of Lemma 5.6, we need a few auxiliary results.

Lemma A.1. *Let \mathcal{A} , \mathfrak{F} , \mathcal{D} and Λ be as above and assume Hypotheses 1 and 3 hold. Then, if $A \in \mathcal{Y}_{\mathfrak{F}}^j(\mathcal{A})$, the operator $A\Lambda^{-j}$ is bounded.*

The proof of this lemma will be a consequence of

Lemma A.2. *Let \mathcal{A} , \mathfrak{F} , \mathcal{D} and Λ be as above and assume Hypotheses 1 and 3 hold. Then, if $A_1, A_2 \in \mathcal{A}$, the operators $A_1\Lambda^{-1}$ and $A_1A_2\Lambda^{-2}$ are bounded.*

Proof. Let us show first that $A_1\Lambda^{-1}$ is bounded. Since \mathcal{D} is a core for Λ , it suffices to show that there is a constant C such that

$$\|A_1 f\|^2 \leq C \|\Lambda f\|^2 \quad \forall f \in \mathcal{D}.$$

This is an immediate consequence of

$$\|\Lambda f\|^2 = \|f\|^2 + \sum_{A \in \mathcal{A}} \|A f\|^2.$$

In order to show that $A_1 A_2 \Lambda^{-2}$ is bounded, we will show that there are constants τ and C such that

$$\|A_1 A_2 f\|^2 \leq C \|\Lambda^2 f + (\tau - 1)f\|^2. \quad (\text{A.3})$$

We can write the following equality:

$$\begin{aligned} \|(\Lambda^2 - 1)f + \tau f\|^2 &= \tau^2 \|f\|^2 + 2\tau \sum_{A \in \mathcal{A}} \|Af\|^2 + \sum_{A, B \in \mathcal{A}} \langle f, A^* A B^* B f \rangle \\ &= \tau^2 \|f\|^2 + 2\tau \sum_{A \in \mathcal{A}} \|Af\|^2 \\ &\quad + \sum_{A, B \in \mathcal{A}} (\|ABf\|^2 + \langle f, [A^* A, B^*] B f \rangle). \end{aligned}$$

We can write the operator intervening in the last term as

$$[A^* A, B^*] B = A^* [A, B^*] B + [B, A]^* A B.$$

Because of Hypothesis 1, this implies that there are positive constants C_{ABC} such that

$$\begin{aligned} \|(\Lambda^2 - 1)f + \tau f\|^2 &\geq \tau^2 \|f\|^2 + 2\tau \sum_{A \in \mathcal{A}} \|Af\|^2 + \sum_{B, C \in \mathcal{A}} \|BCf\|^2 \\ &\quad - \sum_{A, B, C \in \mathcal{A}} C_{ABC} \|Af\| \|BCf\|. \end{aligned}$$

If we use now

$$2xy \leq x^2 s^2 + \frac{y^2}{s^2} \quad x, y \geq 0, \quad s > 0,$$

we see that we can choose τ big enough to have

$$\|(\Lambda^2 - 1)f + \tau f\|^2 \geq \tau^2 \|f\|^2 + \frac{1}{2} \sum_{A \in \mathcal{A}} \|Af\|^2 + \frac{1}{2} \sum_{B, C \in \mathcal{A}} \|BCf\|^2.$$

This immediately implies (A.3). \square

This lemma can now be used to prove Lemma A.1.

Proof of Lemma A.1. We want to show that $A \in \mathcal{Y}_{\frac{1}{8}}^i(\mathcal{A})$ implies $A \Lambda^{-i}$ bounded. We already treated the cases $i = 1$ and $i = 2$. For the other cases, we proceed by induction. Let us fix $j > 2$ and assume the assertion has been proved for $i < j$. Then the operators of the form

$$A_1 A_2 \dots A_j \Lambda^{-j} \quad A_i \in \mathcal{A}, \quad (\text{A.4})$$

are bounded. We distinguish two cases.

$j = 2n$. We write the operator of (A.4) as

$$A_1 A_2 \Lambda^{-2} \cdot \Lambda^2 A_3 A_4 \Lambda^{-4} \cdot \dots \cdot \Lambda^{2n-2} A_{2n-2} A_{2n} \Lambda^{-2n}.$$

We show that operators of the form

$$\Lambda^{2m-2} A B \Lambda^{-2m} \quad A, B \in \mathcal{A}, \quad m \leq n,$$

are bounded. We write

$$\Lambda^{2m-2} A B \Lambda^{-2m} = A B \Lambda^{-2} + [\Lambda^{2m-2}, A B] \Lambda^{-(2m-1)} \Lambda^{-1}.$$

The first term is bounded by Lemma A.2. The second term is bounded by noticing that $[\Lambda^{2m-2}, A B] \in \mathcal{Y}_{\mathfrak{F}}^{2m-1}(\mathcal{A})$ and using the induction hypothesis. $j = 2n + 1$. We write the operator of (A.4) as

$$A_1 A_2 \Lambda^{-2} \cdot \Lambda^2 A_3 A_4 \Lambda^{-4} \cdot \dots \cdot \Lambda^{2n} A_{2n+1} \Lambda^{-2n-1}.$$

The first terms are bounded exactly the same way as before. Concerning the last term, we have

$$\Lambda^{2n} A_{2n+1} \Lambda^{-2n-1} = A_{2n+1} \Lambda^{-1} + [\Lambda^{2n}, A_{2n+1}] \Lambda^{-2n} \Lambda^{-1},$$

which is bounded by Lemma A.2 and the induction hypothesis, noticing that the commutator belongs to $\mathcal{Y}^{2n}(\mathcal{A})$.

This completes the proof of the lemma. \square

We need another result from [EPR99a].

Lemma A.3. *Let $\{A(z)\} \subset \mathfrak{B}(\mathcal{H})$ be a family of uniformly bounded operators, $\Lambda \geq 1$ a self-adjoint operator and let $F(\lambda, z)$ be a real, positive bounded function. Then*

$$\left\| \int_0^\infty A(z) F(\Lambda, z) f dz \right\| \leq \sup_{y \geq 0} \|A(y)\| \|f\| \int_0^\infty \sup_{\lambda \geq 1} F(\lambda, z) dz, \quad \forall f \in \mathcal{H}. \tag{A.5}$$

If furthermore $A = A(z)$ is independent of z , one has the bound

$$\left\| \int_0^\infty A F(\Lambda, z) f dz \right\| \leq \|A\| \|f\| \sup_{\lambda \geq 1} \int_0^\infty F(\lambda, z) dz, \quad \forall f \in \mathcal{H}. \tag{A.6}$$

Lemma A.4. *Let Λ, \mathfrak{F} and \mathcal{A} be as above and assume Hypotheses 1 and 3 hold. If $X \in \mathcal{Y}_{\mathfrak{F}}^j(\mathcal{A})$, then the operators*

$$\Lambda^\beta X \Lambda^\gamma \quad \text{with} \quad \beta + \gamma \leq -j$$

are bounded.

If $Y \in \mathcal{L}$ is such that $[Y, \Lambda^2] \in \mathcal{Y}_{\mathfrak{F}}^j(\mathcal{A})$, then the operators

$$\Lambda^\beta [\Lambda^\alpha, Y] \Lambda^\gamma \quad \text{with} \quad \alpha + \beta + \gamma \leq 2 - j$$

are bounded.

If $X, Y \in \mathcal{L}$ are such that

$$[X, \Lambda^2] \in \mathcal{Y}_{\mathfrak{F}}^j(\mathcal{A}), \quad [Y, \Lambda^2] \in \mathcal{Y}_{\mathfrak{F}}^k(\mathcal{A}) \quad \text{and} \quad [[\Lambda^2, X], Y] \in \mathcal{Y}_{\mathfrak{F}}^{j+k-2}(\mathcal{A}),$$

then the operators

$$\Lambda^\beta [[\Lambda^\alpha, X], Y] \Lambda^\gamma \quad \text{with} \quad \alpha + \beta + \gamma \leq 4 - j - k$$

are bounded.

Proof. Let us prove the first assertion. The case $\gamma = 0$ is handled by noticing that

$$\Lambda^\beta X = \Lambda^{\beta+j} (X^* \Lambda^{-j})^*,$$

and that both operators of the latter product are bounded by Lemma A.1. The case $\beta = 0$ is handled in the same way by considering the adjoint.

The proof for the other cases follows exactly [EPR99a]. We will demonstrate the techniques involved by proving the third assertion, assuming the first two assertions hold. The second assertion can be proved in a similar way without using the third one.

We will first assume that $\alpha \in (-2, 0)$. In this case, we can write (see e.g. [Kat80, § V.3.11])

$$\Lambda^\alpha = C_\alpha \int_0^\infty z^{\alpha/2} (z + \Lambda^2)^{-1} dz, \quad C_\alpha = -\frac{\sin(\pi\alpha/2)}{\pi}. \tag{A.7}$$

We notice moreover that it is possible to write

$$\begin{aligned} [(z + \Lambda^2)^{-1}, X], Y &= (z + \Lambda^2)^{-1} [[\Lambda^2, X], Y] (z + \Lambda^2)^{-1} \\ &\quad + (z + \Lambda^2)^{-1} [\Lambda^2, X] (z + \Lambda^2)^{-1} [\Lambda^2, Y] (z + \Lambda^2)^{-1} \tag{A.8} \\ &\quad + (z + \Lambda^2)^{-1} [\Lambda^2, Y] (z + \Lambda^2)^{-1} [\Lambda^2, X] (z + \Lambda^2)^{-1}. \end{aligned}$$

If we substitute the expression (A.7) in $\Lambda^\beta [[\Lambda^\alpha, X], Y] \Lambda^\gamma$ and use (A.8), we get three terms, which we call T_1, T_2 and T_3 , and which will be estimated separately.

Term T_1 . This term is given by

$$T_1 = C_\alpha \int_0^\infty z^{\alpha/2} \frac{\Lambda^\beta}{z + \Lambda^2} [[\Lambda^2, X], Y] \frac{\Lambda^\gamma}{z + \Lambda^2} dz.$$

We define $B = [[\Lambda^2, X], Y] \in \mathcal{Y}^{j+k-2}(\mathcal{A})$ and write

$$\begin{aligned} T_1 &= C_\alpha \int_0^\infty z^{\alpha/2} \Lambda^\beta B \frac{\Lambda^\gamma}{(z + \Lambda^2)^2} dz + C_\alpha \int_0^\infty z^{\alpha/2} \frac{\Lambda^\beta}{z + \Lambda^2} [\Lambda^2, B] \frac{\Lambda^\gamma}{(z + \Lambda^2)^2} dz \\ &\equiv C_\alpha (T_{11} + T_{12}). \end{aligned}$$

The term T_{11} is estimated by writing, for any $f \in \mathcal{H}$,

$$\begin{aligned} \|T_{11} f\| &= \left\| \Lambda^\beta B \Lambda^{2-j-k-\beta} \int_0^\infty z^{\alpha/2} \frac{\Lambda^{\gamma+\beta+j+k-2}}{(z + \Lambda^2)^2} f dz \right\| \\ &\leq \|f\| \left\| \Lambda^\beta B \Lambda^{2-j-k-\beta} \right\| \sup_{\lambda \geq 1} \int_0^\infty z^{\alpha/2} \frac{\lambda^{\gamma+\beta+j+k-2}}{(z + \lambda^2)^2} dz \\ &= \|f\| \left\| \Lambda^\beta B \Lambda^{2-j-k-\beta} \right\| \sup_{\lambda \geq 1} \int_0^\infty s^{\alpha/2} \frac{\lambda^{\alpha+\gamma+\beta+j+k-4}}{(s + 1)^2} ds. \end{aligned}$$

Since the assumption yields $B \in \mathcal{Y}^{j+k-2}(\mathcal{A})$, the norm is bounded. The integral is also bounded because, by assumption, we have $\alpha + \gamma + \beta \leq 4 - j - k$.

To bound T_{12} , we observe that $[\Lambda^2, B] \in \mathcal{Y}^{j+k-1}(\mathcal{A})$. Using (A.6), we find the bound

$$\begin{aligned} \|T_{12} f\| &= \left\| \int_0^\infty z^{\alpha/2} \frac{\Lambda^\beta}{z + \Lambda^2} [\Lambda^2, B] \Lambda^{3-j-k-\beta} \frac{\Lambda^{\gamma+\beta+j+k-3}}{(z + \Lambda^2)^2} f dz \right\| \\ &\leq \|f\| \sup_{y>0} \left\| \frac{\Lambda^\beta}{y + \Lambda^2} [\Lambda^2, B] \Lambda^{3-j-k-\beta} \right\| \int_0^\infty z^{\alpha/2} \sup_{\lambda \geq 1} \frac{\lambda^{\gamma+\beta+j+k-3}}{(z + \lambda^2)^2} dz. \end{aligned}$$

This expression is bounded when $\alpha + \beta + \gamma \leq 4 - j - k$ and $\alpha \in (-2, 0)$. This can be seen by making as before the substitution $z \mapsto \lambda^2 s$.

Before we go on, we introduce the notation $\Lambda_z \equiv (z + \Lambda^2)^{-1}$.

Term T_2 . This term is given by

$$T_2 = C_\alpha \int_0^\infty z^{\alpha/2} \frac{\Lambda^\beta}{z + \Lambda^2} A \frac{1}{z + \Lambda^2} B \frac{\Lambda^\gamma}{z + \Lambda^2} dz,$$

where we defined

$$A = [\Lambda^2, X] \quad \text{and} \quad B = [\Lambda^2, Y].$$

Since $[\Lambda_z, B] = \Lambda_z[B, \Lambda^2]\Lambda_z$, the term appearing under the integral can be written as

$$\Lambda^\beta \Lambda_z A \Lambda_z B \Lambda_z \Lambda^\gamma = \Lambda^\beta \Lambda_z A B \Lambda_z^2 \Lambda^\gamma + \Lambda^\beta \Lambda_z A \Lambda_z [B, \Lambda^2] \Lambda_z^2 \Lambda^\gamma.$$

According to this, the term T_2 is split into two terms T_{21} and T_{22} . We have

$$\|T_{21}f\| \leq \|f\| \sup_{y>0} \|\Lambda^\beta \Lambda_y A B \Lambda^{-\beta-j-k}\| \int_0^\infty s^{\alpha/2} \sup_{\lambda \geq 1} \frac{\lambda^{\alpha+\beta+\gamma+j+k-4}}{(s+1)^2} ds.$$

The integral is bounded by hypothesis. The norm is also bounded, because $AB \in \mathcal{Y}^{j+k}(\mathcal{A})$. For the second term, we have

$$\|T_{22}f\| \leq \|f\| \sup_{y>0} \|\Lambda^\beta \Lambda_y A \Lambda_y [\Lambda^2, B] \Lambda^{-\beta-j-k}\| \int_0^\infty s^{\alpha/2} \sup_{\lambda \geq 1} \frac{\lambda^{\alpha+\beta+\gamma+j+k-4}}{(s+1)^2} ds.$$

This is bounded in the same fashion, noticing that

$$\begin{aligned} \sup_{y>0} \|\Lambda^\beta \Lambda_y A \Lambda_y [\Lambda^2, B] \Lambda^{-\beta-j-k}\| &\leq \sup_{x>0} \|\Lambda^\beta \Lambda_x A \Lambda^{-\beta-j}\| \\ &\quad \times \sup_{y>0} \left\| \frac{\Lambda^2}{y + \Lambda^2} \Lambda^{j+\beta-2} [\Lambda^2, B] \Lambda^{-\beta-j-k} \right\|. \end{aligned}$$

Term T_3 . It can be bounded in the same way as T_2 by symmetry.

We now have to check the assertion for the other values of α . If $\alpha = 0$ or $\alpha = 2$, it holds trivially. For $\alpha > 0$, we proceed by induction, using the equality

$$\begin{aligned} [[\Lambda^{\alpha+2}, X], Y] &= \Lambda^2 [[\Lambda^\alpha, X], Y] + \Lambda^\alpha [[\Lambda^2, X], Y] \\ &\quad + [\Lambda^\alpha, Y][\Lambda^2, X] + [\Lambda^2, Y][\Lambda^\alpha, X]. \end{aligned} \tag{A.9}$$

For $\alpha = -2$, the assertion is proved using equality (A.8) with $z = 0$. For $\alpha < -2$, we also proceed by induction, using (A.9) with 2 replaced by -2 . This completes the proof of Lemma 5.6. \square

B. Proof of Proposition 2.4

Proposition B.1. \mathcal{T}^l , as defined in (2.10), extends uniquely to a quasi-bounded strongly continuous semigroup on $L^2(\mathcal{X}, dx)$. Its generator L acts like \mathcal{L} on functions in $C_0^\infty(\mathcal{X})$.

Proof. See the proof of Lemma A.1 in [EPR99a]. \square

We now turn to the question of the domain of the generator L . Recall that \mathfrak{L} is the set of all formal expressions of the form

$$\sum_{|l| \leq k} a_l(x) D^l, \quad k \geq 0, \quad a \in C^\infty(\mathbf{R}^n).$$

To any element $L \in \mathfrak{L}$ having the above form, we associate its formal adjoint $L^* \in \mathfrak{L}$ in an obvious way. In the sequel, the notation $\langle f, g \rangle$ will be used to denote the scalar product in L^2 if $f, g \in L^2$ and the evaluation $f(g)$ if f is a distribution and $g \in C_0^\infty(\mathbf{R}^n)$. We hope this slight ambiguity will not be too misleading.

We associate to every $L \in \mathfrak{L}$ the operator $T_L : \mathcal{D}(T_L) \rightarrow L^2(\mathbf{R}^n)$ by

$$(T_L f)(x) = Lf(x) \quad \text{and} \quad \mathcal{D}(T_L) = \{f \in L^2 \mid Lf \in L^2\},$$

where Lf has to be understood in the sense of distributions, i.e.

$$(Lf)(g) \equiv f(L^*g) \quad \text{for all} \quad g \in C_0^\infty(\mathbf{R}^n).$$

We also define the operator $S_L : \mathcal{D}(S_L) \rightarrow L^2(\mathbf{R}^n)$ by

$$S_L = \overline{T_L \upharpoonright C_0^\infty}.$$

The operators T_L and S_L are usually called the *minimal operator* and the *maximal operator* constructed from the *formal operator* L . The following result is classical, so we do not give its proof here:

Proposition B.2. For every $L \in \mathfrak{L}$, we have $T_L^* = S_L^*$ and $S_L^* = T_L^*$. In particular, this shows that T_L is closed. \square

We prove now the quasi m -dissipativity of $S_{\mathcal{L}}$. We define

$$\tilde{\mathcal{L}} \equiv \mathcal{L} - \sum_{i=1}^M \gamma_i - 1.$$

By definition, if $S_{\tilde{\mathcal{L}}}$ is strictly m -dissipative, $S_{\mathcal{L}}$ is quasi m -dissipative. It is well-known that an equivalent characterization of strict m -dissipativity is that

- (a) $S_{\tilde{\mathcal{L}}}$ is strictly dissipative and
- (b) $\text{Range}(S_{\tilde{\mathcal{L}}}) = \mathcal{H}$.

Proposition B.3. Assume Assumption 0 holds. Then $S_{\tilde{\mathcal{L}}}$ is strictly m -dissipative.

Remark B.4. It is clear that the statement holds if we consider the minimal operator in $L^2(\mathcal{K}, dx)$, where \mathcal{K} is some compact domain of \mathcal{X} . The idea is to approximate \mathcal{X} by a sequence of increasing compact domains and to control the rest terms.

This proposition fills a gap in [EPR99a], since the statement “ $\text{Re}(f, L^*f) = -\frac{1}{2} \|\sigma^T \nabla f\|^2 + (f, \text{div } b f) \leq B \|f\|^2$ ” in the proof of Lemma A.1 is not justified for every $f \in \mathcal{D}(L^*)$.

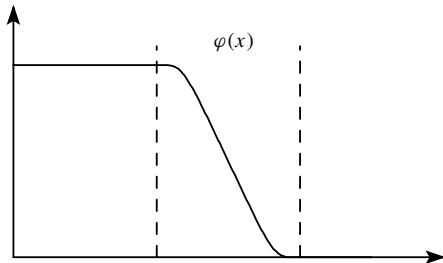


Fig. B.1.

Proof. Property (a) is immediate. By the closed-range theorem, property (b) is equivalent to the statement

$$(b') \quad f \in L^2 \text{ and } \tilde{\mathcal{L}}^* f = 0 \text{ imply } f = 0.$$

Assume on the contrary that there exists a non-vanishing function $f \in L^2$ for which $\tilde{\mathcal{L}}^* f = 0$ holds in the sense of distributions. Since $\tilde{\mathcal{L}}^*$ is hypoelliptic, f must be a C^∞ function. Let us choose some function $\varphi \in C_0^\infty(\mathbf{R}_+)$ such that $\varphi(x) = 1$ if $x \in [0, 1]$. We also define

$$\begin{aligned} \varphi_n : \mathcal{X} &\rightarrow \mathbf{R} \\ x &\mapsto \varphi(G(x)/n). \end{aligned}$$

By assumption, $\tilde{\mathcal{L}}^* f = 0$, so we have

$$0 = 2\text{Re}\langle \varphi_n f, \tilde{\mathcal{L}}^* f \rangle = \langle \varphi_n f, \tilde{\mathcal{L}}^* f \rangle + \langle \tilde{\mathcal{L}}^* f, \varphi_n f \rangle.$$

Since $\varphi_n \in C_0^\infty$ and all the other functions are C^∞ , we can make all the formal manipulations we want. In particular, we have

$$\langle \tilde{\mathcal{L}}^* f, \varphi_n f \rangle = \langle f, \tilde{\mathcal{L}} \varphi_n f \rangle \quad \Rightarrow \quad \langle f, (\varphi_n \tilde{\mathcal{L}}^* + \tilde{\mathcal{L}} \varphi_n) f \rangle = 0. \tag{B.1}$$

Recall that $\tilde{\mathcal{L}}$ is given by

$$\begin{aligned} \tilde{\mathcal{L}} &= \sum_{i=1}^M \lambda_i^2 \gamma_i T_i \partial_{r_i}^2 - \sum_{i=1}^M \gamma_i (r_i - \lambda_i^2 F_i(p, q)) \partial_{r_i} \\ &\quad + X^{H_S} - \sum_{i=1}^M r_i X^{F_i} - \sum_{i=1}^M \gamma_i - 1 \\ &\equiv \sum_{i=1}^M \zeta_i \partial_{r_i}^2 + Y_0 - 1. \end{aligned} \tag{B.2}$$

Straightforward computation yields

$$\begin{aligned}
\varphi_n \tilde{\mathcal{L}}^* + \tilde{\mathcal{L}}\varphi_n &= 2 \sum_{i=1}^M \zeta_i \partial_{r_i} \varphi_n \partial_{r_i} + \sum_{i=1}^M \zeta_i (\partial_{r_i}^2 \varphi_n) + [Y_0, \varphi_n] - \varphi_n \\
&= 2 \sum_{i=1}^M \zeta_i \partial_{r_i} \varphi_n \partial_{r_i} \\
&\quad + \sum_{i=1}^M \zeta_i \left(\frac{1}{n} (\partial_{r_i}^2 G) \varphi''(G/n) + \frac{1}{n^2} (\partial_{r_i} G)^2 \varphi'(G/n) \right) \\
&\quad + \frac{1}{n} \sum_{i=1}^M \frac{\gamma_i}{\lambda_i^2} (r_i - \lambda_i^2 F_i(p, q))^2 \varphi'(G/n) - \varphi_n \\
&\equiv 2 \sum_{i=1}^M \zeta_i \partial_{r_i} \varphi_n \partial_{r_i} + \Phi_n - \varphi_n.
\end{aligned} \tag{B.3}$$

Using Assumption 0, we next verify that $|\Phi_n(x)| \leq \tilde{C}$ for all $x \in \mathcal{X}$ and for all $n \geq 1$. We define

$$c_1 \equiv \sup_{x \geq 0} \varphi''(x) \quad \text{and} \quad c_2 \equiv \sup_{x \geq 0} x \varphi'(x).$$

An elementary computation shows that Assumption 0 implies that there are constants $c_3, \dots, c_5 > 0$ for which

$$|\partial_{r_i}^2 G(x)| \leq c_3, \quad |\partial_{r_i} G(x)|^2 \leq c_4 G(x), \quad \text{and} \quad (r_i - \lambda_i^2 F_i(p, q))^2 \leq c_5 G(p, q, r).$$

We thus have

$$\begin{aligned}
|\Phi_n(x)| &\leq \sum_{i=1}^M \left(\frac{\zeta_i c_3}{n} |\varphi''(G/n)| + \frac{\zeta_i c_4}{n} |(G/n) \varphi'(G/n)| + \frac{\gamma_i c_5}{\lambda_i^2} |(G/n) \varphi'(G/n)| \right) \\
&\leq \sum_{i=1}^M \left(\zeta_i \frac{c_1 c_3 + c_2 c_4}{n} + \frac{\gamma_i c_2 c_5}{\lambda_i^2} \right) \leq \tilde{C},
\end{aligned}$$

as asserted. Moreover, the first part of Assumption 0 implies that there exist constants $C, \alpha > 0$ such that

$$\text{supp } \Phi_n \subset \{x \in \mathcal{X} \mid \|x\|^\alpha \geq n/C\}. \tag{B.4}$$

Substituting (B.3) back into (B.1), we get

$$0 = -2 \sum_{i=1}^M \zeta_i \|\sqrt{\varphi_n} \partial_{r_i} f\|^2 - \|\sqrt{\varphi_n} f\|^2 + \int_{\mathcal{X}} \Phi_n(x) |f(x)|^2 dx. \tag{B.5}$$

Since $f \in L^2(\mathcal{X})$, one has

$$\lim_{n \rightarrow \infty} \|\sqrt{\varphi_n} f\|^2 = \|f\|^2.$$

Moreover, the uniform boundedness of Φ_n together with property (B.4) imply that

$$\lim_{n \rightarrow \infty} \int_{\mathcal{X}} \Phi_n(x) |f(x)|^2 dx = 0.$$

This supplies the required contradiction to (B.5), thus establishing the strict m -dissipativity of $S_{\tilde{\mathcal{L}}}$. \square

We complete now the

Proof of Proposition 2.4. It only remains to be proved that $L = S_{\mathcal{L}}$ and that $L^* = S_{\mathcal{L}^*}$.

It is clear that the generator L of \mathcal{T}^t satisfies $S_{\mathcal{L}} \subset L$. Since $S_{\mathcal{L}}$ is quasi m -dissipative, i.e. has no proper quasi dissipative extension, and since the generator of a quasi-bounded semigroup is always quasi m -dissipative, we must have $L = S_{\mathcal{L}}$.

Concerning the adjoint, we have by Proposition B.2, $L^* = T_{\mathcal{L}^*}$. It is possible to retrace the above argument for \mathcal{L}^* to show that $S_{\mathcal{L}^*}$ is quasi m -dissipative. Since L^* is also quasi m -dissipative and $S_{\mathcal{L}^*} \subset L^*$, we must have $L^* = S_{\mathcal{L}^*}$. \square

C. Proof of Lemma 8.4

Using the technique developed in Appendix B, we can now turn to the proof of Lemma 8.4. Recall that K is given by (3.10) and that

$$W^{(\gamma)} = L^2(\mathcal{X}, G^{2\gamma} dx).$$

Moreover, K_{γ} is the maximal operator constructed from K when considering it as a differential operator in $W^{(\gamma)}$. We have

Proposition C.1. $C_0^{\infty}(\mathcal{X})$ is a core for K_{γ} .

Proof. We introduce the unitary operator $U : W^{(\gamma)} \rightarrow L^2(\mathcal{X})$ defined by

$$(Uf)(x) = G^{\gamma}(x)f(x).$$

We also define $K_{\gamma}^0 \equiv \overline{K_{\gamma} \upharpoonright C_0^{\infty}(\mathcal{X})}$. The operators K_{γ} and K_{γ}^0 are unitarily equivalent to the operators \tilde{K}_{γ} and \tilde{K}_{γ}^0 respectively by the following relations:

$$\begin{array}{ccc} \mathcal{D}(K_{\gamma}) \xrightarrow{K_{\gamma}} W^{(\gamma)} & & \mathcal{D}(K_{\gamma}^0) \xrightarrow{K_{\gamma}^0} W^{(\gamma)} \\ U \downarrow \uparrow U^{-1} & U \downarrow \uparrow U^{-1} & U \downarrow \uparrow U^{-1} \\ \mathcal{D}(\tilde{K}_{\gamma}) \xrightarrow{\tilde{K}_{\gamma}} L^2(\mathcal{X}) & & \mathcal{D}(\tilde{K}_{\gamma}^0) \xrightarrow{\tilde{K}_{\gamma}^0} L^2(\mathcal{X}) \end{array}$$

By construction, \tilde{K}_{γ} is maximal. Thus, by Proposition B.2, its adjoint \tilde{K}_{γ}^* is minimal. It is immediate that the formal expressions for \tilde{K}_{γ}^* and K_{γ}^0 are given by

$$\tilde{K}_{\gamma}^* = G^{-\gamma} K^* G^{\gamma} \quad \text{and} \quad \tilde{K}_{\gamma}^0 = G^{\gamma} K G^{-\gamma}.$$

It is now a simple exercise to retrace the proof of Proposition B.3 to see that \tilde{K}_{γ}^* and \tilde{K}_{γ}^0 are both m -accretive. The remark of Sect. 2.2 concerning the adjoints of m -accretive operators implies that \tilde{K}_{γ} is also m -accretive. Since $\tilde{K}_{\gamma}^0 \subset \tilde{K}_{\gamma}$, we must have $\tilde{K}_{\gamma}^0 = \tilde{K}_{\gamma}$ and thus $K_{\gamma}^0 = K_{\gamma}$. This proves the assertion. \square

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