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Self-adjointness of some infinite-dimensional elliptic operators and application to stochastic quantization

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Abstract. We consider an operator $\mathring{K}\varphi = L\varphi - \langle CDU(x), D\varphi \rangle$ in a Hilbert space H, where L is an Ornstein–Uhlenbeck operator, $U \in W^{1,4}(H, \mu)$ and μ is the invariant measure associated with L. We show that \mathring{K} is essentially self-adjoint in the space $L^2(H, \nu)$ where ν is the "Gibbs" measure $\nu(dx) = Z^{-1}e^{-2U(x)}dx$. An application to Stochastic quantization is given.

1. Introduction

We are concerned with the following operator in a separable Hilbert space (norm $|\cdot|$, inner product $\langle\cdot,\cdot\rangle$):

$$\mathring{K}\varphi = L\varphi - \langle C^{1/2}DU(x), C^{1/2}D\varphi \rangle, \qquad (1.1)$$

defined on $D(\mathring{K}) = D(L) \cap C_h^2(H)$ and where the operator *L*, defined by

$$L\varphi = \frac{1}{2} \operatorname{Tr} \left[CD^2 \varphi \right] + \langle Ax, D\varphi \rangle, \quad \varphi \in D(L), \tag{1.2}$$

is the Ornstein–Uhlenbeck operator, see Section 2 for precise definitions.

Here A and C are linear operators in H, and U a is mapping from H into \mathbb{R} . Moreover $D\varphi$ represents the Fréchet derivative of φ .

Such operators arise in several applications as the Landau–Ginzburg equations [3], and in stochastic quantization. An application will be given at the end of this paper.

Let us formulate our assumptions.

Hypothesis 1.1.

(i) A : D(A) ⊂ H → H is self-adjoint strictly negative.
(ii) C = (-A)^{-ε} for some ε ∈ (0, 1).

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(iii) The operator $Q = \frac{1}{2}(-A)^{-(1+\varepsilon)}$ is of trace class.

Hence we shall denote by $\{-\alpha_k\}$ the sequence of (negative) eigenvalues of *A* and by $\{e_k\}$ the corresponding complete orthonormal system. We have also, by denoting $\lambda_k = \frac{1}{2}\alpha_k^{-(1+\varepsilon)}$

$$Qe_k = \lambda e_k, \qquad k \in \mathbb{N}.$$

Moreover we shall denote by μ the gaussian measure $\mathcal{N}(0, Q)$ of mean 0 and covariance operator Q.

Now we consider the probability measure

$$v(dx) = Z^{-1} e^{-2U(x)} \mu(dx),$$

where

$$Z = \int_H e^{-2U(x)} \mu(dx),$$

under the following assumptions

Hypothesis 1.2.

(i) e^{-U} belongs to $L^p(H, \mu)$ for any $p \ge 1$. (ii) $U: H \to \mathbb{R}$ belongs to $W^{1,4}_C(H, \nu)$. (¹)

Our goal is to show that, under Hypotheses 1.1 and 1.2, \mathring{K} is essentially self-adjoint in $L^2(H, \nu)$. In this way we are able to construct a semigroup e^{tK} , $t \ge 0$ were Kis the closure of \mathring{K} . We notice that formally e^{tK} , $t \ge 0$ is the transition semigroup corresponding to the differential stochastic equation

$$dX = (AX - CDU)dt + C^{1/2}dW(t), \quad X(0) = x,$$
(1.3)

where W is a cylindrical Wiener process taking values in H. We notice that under assumptions above we are not able to find even a weak solution of equation (1.3).

The problem of self-adjointness of the operator \mathring{K} has been studied by several authors under different assumptions using the Dirichlet forms theory, see e.g [17], [1], [16], and references therein. The existence of a self-adjoint extension K of \mathring{K} follows from the closability of the Dirichlet form naturally associated with \mathring{K} . Then the main problem consists in proving that K is the closure of \mathring{K} (Uniqueness problem).

Our approach is different. We show that the operator \hat{K} is symmetric, and that the image of $\lambda - \hat{K}$ is dense on $L^2(H, \nu)$ for $\lambda > 0$. This will imply, by the Lumer-Phillips theorem, see [18], that \hat{K} is closable and its closure K is self-adjoint (We will denote by D(K) its domain).

¹ Spaces $W_C^{1,4}(H, \nu)$ are introduced in the next section. Moreover if H and K are Hilbert spaces we denote by $C_b(H; K)$ the Banach space of all uniformly continuous and bounded mappings from H into K, endowed with the sup norm $\|\cdot\|_0$. Moreover, for any $k \in \mathbb{N}$, $C_b^k(H; K)$ will represent the Banach space of all mappings from H into K, that are uniformly continuous and bounded together with their Fréchet derivatives of order less or equal to k endowed with their natural norm $\|\cdot\|$. If $K = \mathbb{R}$ we set $C_b(H; K) = C_b(H)$ and $C_b^k(H; K) = C_b^k(H)$

The main tools in order to prove density of $(\lambda - \mathring{K})(L^2(H, \nu))$ are an approximation of *U* by smooth functions and an *a priori* estimate on $W_C^{1,4}(H, \nu)$. This estimate is similar to one proved in the papers [17] and [16].

We can also show, see Remark 3.4 below, that the set $\mathscr{E}_A(H)$ of all functions φ of the form

$$\varphi(x) = \operatorname{Re} \sum_{k=1}^{n} a_k e^{i \langle x, h_k \rangle},$$

where $n \in \mathbb{N}$, $h_1, ..., h_n \in D(A)$ and $a_1, a_2, ..., a_n \in \mathbb{C}$, is a core for K.

We notice that our assumptions are close to that of [16], but our method seems to be simpler and can be applied to non gradient Dirichlet operator, by replacing symmetry with dissipativity, see Remark 3.5.

As in [16] we give finally an application to Stochastic Quantization in dimensions 2.

Stochastic Quantization has been studied by several authors see [15], [5], [2], [14], [20]. In particular in [15], a transition semigroup for equation (1.3) was built by giving a meaning to Girsanov formula for $\varepsilon < 1/10$.

2. Preliminary results

We first introduce the Ornstein–Uhlenbeck semigroup R_t , $t \ge 0$, as a family of bounded operators in $C_b(H)$:

$$R_t\varphi(x) = \int_H \varphi(y)\mathcal{N}(e^{tA}x, Q_t)(dy), \ \varphi \in C_b(H),$$
(2.1)

where

$$Q_t = Q\left(1 - e^{2tA}\right), \ t \ge 0, \tag{2.2}$$

and $\mathcal{N}(e^{tA}x, Q_t)$ is the gaussian measure on *H* of mean $e^{tA}x$ and covariance operator Q_t .

Proposition 2.1. For any $\varphi \in C_b(H)$ and any t > 0, $h \in H$, $R_t \varphi$ is differentiable in the direction $C^{1/2}h$ and we have

$$\langle DR_t\varphi, C^{1/2}h\rangle = \int_H \langle \Lambda(t)h, Q_t^{-1/2}y\rangle\varphi(e^{tA}x+y)\mathcal{N}(0, Q_t)(dy), \qquad (2.3)$$

where

$$\Lambda(t) = \sqrt{2} \ (-A)^{1/2} (1 - e^{2tA})^{-1/2} e^{tA}, \ t \ge 0.$$
(2.4)

Moreover

$$|C^{1/2}DR_t\varphi(x)| \le t^{-\frac{1}{2}} \|\varphi\|_0, \ t \ge 0.$$
(2.5)

Proof. (2.3) follows easily from the Cameron-Martin formula, see e.g. [11]. Let us check (2.5). Since

$$\Lambda(t)e_{k} = \sqrt{2} \; \frac{\alpha_{k}^{1/2}e^{-\alpha_{k}t}}{(1 - e^{-2\alpha_{k}t})^{1/2}} \; e_{k},$$

we have

$$\|\Lambda(t)\| = \sup_{k \in \mathbb{N}} \sqrt{2} \ \frac{\alpha_k^{1/2} e^{-\alpha_k t}}{(1 - e^{-2\alpha_k t})^{1/2}} \le t^{-1/2}.$$

Now the conclusion follows integrating with respect to μ .

The semigroup R_t is not strongly continuous in $C_b(H)$. However one can define its infinitesimal generator L, see [6], as the unique linear operator $L: D(L) \subset$ $C_b(H) \rightarrow C_b(H)$ whose resolvent is given by

$$R(\lambda, L)\varphi(x) = \int_0^{+\infty} e^{-\lambda t} R_t \varphi(x) dt, \quad x \in H, \ \varphi \in C_b(H).$$

The following result is a consequence of (2.5).

Proposition 2.2. For any $\varphi \in D(L)$ we have that $C^{1/2}D\varphi \in C_b(H; H)$ and

$$|C^{1/2}D\varphi(x)| \le \sqrt{\pi} \ (\|\varphi\|_0 + \|L\varphi\|_0), \ x \in H.$$
(2.6)

Proof. Let $\varphi \in D(L)$ and set $f = \varphi - L\varphi$. Then from (2.5) it follows that

$$|C^{1/2}DR(\lambda, L)f(x)| \le \sqrt{\pi} ||f||_0,$$

that yields the conclusion.

Let us also recall the following identity, see e.g. [13]

$$\int_{H} L\varphi \,\varphi \,d\nu = -\frac{1}{2} \,\int_{H} |C^{1/2} D\varphi|^2 d\nu, \quad \varphi \in D(L).$$
(2.7)

We now define the Sobolev space $W_C^{1,2}(H, \nu)$. For any $k \in \mathbb{N}$, and for any $\varphi \in C_b^1(H)$, we denote by $D_k \varphi$ the derivative of φ on the direction e_k , and we set $x_k = \langle x, e_k \rangle$. The following identity is well known, see e.g. [13],

$$\int_{H} D_{k} \varphi \psi d\mu = -\int_{H} \varphi D_{k} \psi d\mu + \frac{1}{\lambda_{k}} \int_{H} x_{k} \varphi \psi d\mu, \ \varphi, \psi \in C_{b}^{1}(H), \quad (2.8)$$

and so we obtain the result,

Lemma 2.3. Let $\varphi, \psi \in C_h^1(H)$. Then for any $k \in \mathbb{N}$ the following identity holds

$$\int_{H} D_{k}\varphi\psi d\nu = -\int_{H} \varphi D_{k}\psi d\nu + 2\int_{H} D_{k}U\varphi\psi d\nu + \frac{1}{\lambda_{k}}\int_{H} x_{k}\varphi\psi d\nu.$$
(2.9)

Moreover D_k is closable in $L^2(H, v)$; we shall still denote by D_k its closure.

We can finally define the Sobolev space $W_C^{1,2}(H, \nu)$ by setting

$$\begin{split} W^{1,2}_C(H,\nu) &= \Big\{ \varphi \in L^2(H,\nu) : \ D_i \varphi \in L^2(H,\nu) \quad \forall \, i \in \mathbb{N} \\ &\sum_{i=1}^{\infty} \alpha_i^{-\varepsilon} \int_H |D_i \varphi(x)|^2 \nu(dx) < +\infty \Big\}. \end{split}$$

We conclude this section with a result needed later.

Lemma 2.4. Let $\beta : \mathbb{R} \to \mathbb{R}$ be Lipschitz continuous $(^2)$, $\varphi \in C_b^1(H)$ and p > 1. Then $\beta \circ \varphi \in W_C^{1,p}(H, \nu)$.

Proof. Let $\{\beta_n\} \subset C^1(\mathbb{R})$ be uniformly convergent to β and such that $\sup_{n \in \mathbb{N}} \|\beta_n\|_{\text{Lip}} \leq \|\beta\|_{\text{Lip}}$. Then $\beta_n \circ \varphi \to \beta \circ \varphi$ uniformly and

$$\int_{H} |C^{1/2} D\beta_n \circ \varphi|^2 d\nu = \int_{H} |D\beta_n(\varphi(x))|^2 |C^{1/2} D\varphi(x)|^2 d\nu \le \|\beta\|_{\operatorname{Lip}}^2 \|\varphi\|_1^2.$$

Thus, by a standard argument, we have $\beta \circ \varphi \in W^{1,p}_C(H,\nu)$.

Proposition 2.5. Let $p \ge 1$, and let $\varphi \in W_C^{1,p}(H, \nu)$. Then for any constant $\kappa > 0$ we have $\min\{\varphi, \kappa\} \in W_C^{1,p}(H, \nu)$.

Proof. Let $\{\varphi_n\}$ be a sequence in $C_b^1(H)$ convergent to φ in $W_C^{1,p}(H, \nu)$. Then $\psi_n = \min\{\varphi_n, \kappa\}, n \in \mathbb{N}$ belongs to $W_C^{1,p}(H, \nu)$ by Lemma 2.4. Therefore we have $\psi_n \to \min\{\varphi, \kappa\}$ in $L^p(H, \nu)$, and moreover

$$\|C^{1/2}D\psi_n\|_{L^p(H,\nu)} \le \|C^{1/2}D\varphi_n\|_{L^p(H,\nu)} \le \sup_{n\in\mathbb{N}} \|C^{1/2}D\varphi_n\|_{L^p(H,\nu)} < +\infty.$$

The same argument as before implies the conclusion.

3. The main result

Here we assume that Hypotheses 1.1 and 1.2 hold. We are concerned with the operator \mathring{K} defined by (1.1). By (2.7) it follows immediately that for any $\varphi, \psi \in D(\mathring{K})$ we have

$$\int_{H} \mathring{K}\varphi\psi d\nu = -\frac{1}{2} \int_{H} \langle CD\varphi, D\psi \rangle d\nu.$$
(3.1)

Therefore \mathring{K} is symmetric and consequently closable in $L^2(H, \nu)$.

We now introduce an approximating problem. Let $\{U_n\}$ be a sequence in $C_b^{\infty}(H)$ convergent to U in $W_C^{1,4}(H, \nu)$. Such a sequence can be easily constructed by setting

$$U_n(x) = S_{1/n}[V_n(x)],$$

² $\|\beta\|_{\text{Lip}} = \sup_{x,y\in\mathbb{R}} \frac{|\beta(x)-\beta(y)|}{|x-y|}$

where S_t is an auxiliary strong Feller Ornstein–Uhlenbeck semigroup, see [11, (9.50)], (³) and

$$V_n(x) = \begin{cases} U(x) & \text{if } |U(x)| \le n, \\ \\ \frac{U(x)}{|U(x)|}n & \text{if } |U(x)| > n, \end{cases}$$

and recalling Proposition 2.5.

We define

$$K_n \varphi = L \varphi - \langle C D U_n, D \varphi \rangle, \ \varphi \in D(L).$$
 (3.2)

Then it is easy to check, by a simple fixed point argument taking into account Proposition 2.2, that the resolvent set of K_n contains the half line $(0, +\infty)$, and (since DU_n is regular) that K_n is the infinitesimal generator of the transition semigroup corresponding to the differential stochastic equation

$$\begin{cases} dX_n(t,x) = (AX_n(t,x) - DU_n(X_n(t,x)))dt + \sqrt{C} \ dW(t), \\ X_n(0,x) = x. \end{cases}$$
(3.3)

Consequently

$$R(\lambda, K_n)f(x) = \int_0^{+\infty} e^{-\lambda t} \mathbb{E}(f(X_n(t, x))dt, f \in C_b(H),$$
(3.4)

Let now $f \in C_b^2(H)$ and let φ_n be the solution to

$$\lambda \varphi_n - K_n \varphi_n = f. \tag{3.5}$$

It is easy to see that $\varphi_n \in D(L) \cap C_b^2(H)$, and

$$\lambda \varphi_n - \mathring{K} \varphi_n + \langle C^{1/2} (DU - DU_n), C^{1/2} D \varphi_n \rangle = f.$$
(3.6)

In order to prove that \mathring{K} is essentially self-adjoint we will show that the image of $\lambda - \mathring{K}$ is dense in $L^2(H, \nu)$. For this we need an estimate of $|C^{1/2}D\varphi_n|$ in $L^4(H, \nu)$. This is provided by the following result.

Proposition 3.1. Let $f \in C_b^2(H)$. Then the solution φ_n to (3.5) belongs to $D(\mathring{K})$ and there is a positive constant κ depending only on λ , $||f||_{\infty}$, and $||C^{1/2}DU||_{L^4(H,\nu)}$ such that

$$\int_{H} \langle C D \varphi_n, D \varphi_n \rangle^2 d\nu \le \kappa.$$
(3.7)

Proof. We proceed in several steps.

Step 1. For any $\varphi, \psi \in D(\mathring{K})$ we have.

$$\int_{H} K_{n}\varphi \ \psi d\nu = -\frac{1}{2} \ \int_{H} \langle CD\varphi, D\psi \rangle d\nu + \int_{H} \langle CD\varphi, D(U-U_{n}) \rangle \psi d\nu.$$
(3.8)

³ We can choose $S_t \varphi(x) = \int_H \varphi(e^{\frac{t}{2}Q^{-1}}x + y) \mathcal{N}(0, Q(1 - e^{tQ^{-1}}))(dy)$

The simple verification is left to the reader.

Step 2. For any $\varphi \in D(\mathring{K})$ we have

$$\int_{H} \langle K_n D\varphi, CD\varphi \rangle d\nu = -\frac{1}{2} \int_{H} \operatorname{Tr} \left[(CD^2 \varphi)^2 \right] d\nu + \int_{H} \langle D^2 \varphi \ CD\varphi, C(DU - DU_n) \rangle d\nu.$$
(3.9)

We have in fact, taking into account (3.8),

$$\begin{split} \int_{H} \langle K_n D\varphi, CD\varphi \rangle d\nu &= \sum_{i=1}^{\infty} \alpha_i^{-\varepsilon} \int_{H} K_n D_i \varphi \ D_i \varphi d\nu \\ &= -\frac{1}{2} \sum_{i=1}^{\infty} \alpha_i^{-\varepsilon} \int_{H} \langle CDD_i \varphi, DD_i \varphi \rangle d\nu \\ &+ \sum_{i=1}^{\infty} \alpha_i^{-\varepsilon} \int_{H} \langle CDD_i \varphi, D(U-U_n) \rangle D_i \varphi d\nu, \end{split}$$

and the conclusion follows.

Step 3. For any $\varphi \in D(\mathring{K})$ we have

$$2\int_{H} (K_{n}\varphi)^{2} dv + \int_{H} \langle D^{2}\varphi \cdot CD\varphi, CDU_{n} \rangle dv + 2\int_{H} K_{n}\varphi \langle D\varphi, CDU_{n} \rangle dv + \int_{H} \langle D^{2}\varphi \cdot CD\varphi, C(DU - DU_{n}) \rangle dv - 2\int_{H} K_{n}\varphi \langle D\varphi, C(DU - DU_{n}) \rangle dv - 2\int_{H} \langle CD\varphi, D(U - U_{n}) \rangle \langle CD\varphi, DU_{n} \rangle dv = \frac{1}{2}\int_{H} \operatorname{Tr} \left[(CD^{2}\varphi)^{2} \right] dv - \int_{H} \langle ACD\varphi, D\varphi \rangle dv.$$
(3.10)

In fact, setting in (3.8) $\psi = K_n \varphi$ we find

$$\begin{split} \int_{H} (K_{n}\varphi)^{2} d\nu &= -\frac{1}{2} \int_{H} \langle CD\varphi, DK_{n}\varphi \rangle d\nu + \int_{H} \langle CD\varphi, D(U-U_{n}) \rangle K_{n}\varphi d\nu \\ &= -\frac{1}{2} \int_{H} \langle CD\varphi, [D, K_{n}]\varphi \rangle d\nu - \frac{1}{2} \int_{H} \langle CD\varphi, CK_{n}\varphi \rangle d\nu \\ &+ \int_{H} \langle CD\varphi, D(U-U_{n}) \rangle K_{n}\varphi d\nu \end{split}$$

Now, taking into account (3.9) and (3.10), we find

$$\int_{H} (K_{n}\varphi)^{2} d\nu = -\frac{1}{2} \int_{H} \langle CD\varphi, AD\varphi \rangle d\nu + \frac{1}{2} \int_{H} \langle CD\varphi, D\{\langle CDU, D\varphi \rangle\} \rangle d\nu$$
$$= -\frac{1}{2} \int_{H} \langle D^{2}\varphi \ CDU, D\varphi \rangle d\nu + \frac{1}{4} \int_{H} \operatorname{Tr} \left[(CD^{2}\varphi)^{2} \right] d\nu$$

$$= -\frac{1}{2} \int_{H} \langle D^{2}\varphi \ CD\varphi, C(DU - DU_{n}) \rangle d\nu + \int_{H} \langle CD\varphi, DU - DU_{n} \rangle K_{n}\varphi d\nu$$
(3.11)

Setting in (3.8) $\psi = \langle CDU, D\varphi \rangle$ we find

$$\begin{split} &\int_{H} \langle CD\varphi, D\left\{ \langle CDU, D\varphi \rangle \right\} \rangle d\nu \\ &-2 \int_{H} K_{n} \varphi \langle CDU, D\varphi \rangle d\nu + 2 \int_{H} \langle CD\varphi, D(U-U_{n}) \rangle \langle CDU_{n}, D\varphi \rangle d\nu. \end{split}$$

By substituting this in (3.11) the conclusion follows.

Step 4. For any $\varphi \in D(\mathring{K})$ we have

$$\int_{H} \langle CD\varphi, D\varphi \rangle^{2} d\nu = -2 \int_{H} K_{n} \varphi \ \varphi \ \langle CD\varphi, D\varphi \rangle d\nu$$
$$-2 \int_{H} \varphi \ \langle D^{2}\varphi \cdot CD\varphi, CD\varphi \rangle d\nu$$
$$+2 \int_{H} \langle CD\varphi, DU - DU_{n} \rangle \langle CD\varphi, D\varphi \rangle \varphi d\nu. \quad (3.12)$$

Setting in (3.8) $\psi = \varphi \langle C D \varphi, D \varphi \rangle$ we find

$$\begin{split} \int_{H} K_{n}\varphi \; \varphi \; \langle CD\varphi, D\varphi \rangle d\nu &= -\frac{1}{2} \int_{H} \langle CD\varphi, D(\varphi \langle CD\varphi, D\varphi \rangle) \rangle d\nu \\ &+ \langle CD\varphi, DU - DU_{n} \rangle \varphi \langle CD\varphi, D\varphi \rangle, \end{split}$$

that yields (3.12).

Now it is easy to see that

$$\|\varphi_n\|_{\infty} \le \frac{1}{\lambda} \|f\|_{\infty}, \tag{3.13}$$

and

$$\|K_n\varphi_n\|_{\infty} \le 2\|f\|_{\infty}.$$
(3.14)

Step 5. We have

$$\|C^{1/2}D\varphi_n\|_{L^2(H,\nu)} \le \frac{2}{\lambda} \|f\|_{\infty} \left[\sqrt{2\lambda} + \|C^{1/2}D(U-U_n)\|_{L^2(H,\nu)}\right].$$
 (3.15)

We first note that, multiplying (3.5) by φ_n and integrating, we have

$$\lambda \int_{H} \varphi_n^2 dv - \int_{H} K_n \varphi_n \, \varphi_n dv = \int_{H} f \varphi_n dv.$$

Taking into account (3.8), (3.13), and (3.14), it follows

$$\begin{split} \frac{1}{2} \int_{H} \langle CD\varphi_n, CD\varphi_n \rangle d\nu &\leq \int_{H} f\varphi_n d\nu + \int_{H} \langle CD\varphi_n, D(U-U_n) \rangle \varphi_n d\nu \\ &\leq \frac{1}{\lambda} \|f\|_{\infty}^2 + \frac{1}{\lambda} \|f\|_{\infty} \int_{H} |C^{1/2}D\varphi_n| |C^{1/2}D(U-U_n)| d\nu \\ &\leq \frac{1}{\lambda} \|f\|_{\infty}^2 + \frac{1}{\lambda} \|f\|_{\infty} \left(\int_{H} |C^{1/2}D\varphi_n|^2 d\nu\right)^{1/2} \\ &\times \left(\int_{H} |C^{1/2}D(U-U_n)|^2 d\nu\right)^{1/2}, \end{split}$$

that yields the conclusion.

Step 6. We have

$$\left| \int_{H} \langle D^{2} \varphi \cdot Cu, Cv \rangle dv \right| \leq \left(\int_{H} \operatorname{Tr} \left[(CD^{2} \varphi)^{2} \right] dv \right)^{1/2} \left(\int_{H} \langle Cu, u \rangle \langle Cv, v \rangle dv \right)^{1/2} dv \right)^{1/2} dv$$

Step 7. Conclusion.

We now estimate $\int_H \langle C D \varphi_n, D \varphi_n \rangle^2 d\nu$ starting from (3.12), that we write in the form

$$\int_{H} \langle C D \varphi_n, D \varphi_n \rangle^2 \, d\nu = \Lambda_1 + \Lambda_2 + \Lambda_3. \tag{3.16}$$

By (3.15) it follows

$$\Lambda_{1} \leq \frac{2}{\lambda} \|f\|_{\infty}^{2} \int_{H} |C^{1/2} D\varphi_{n}|^{2} d\nu$$

$$\leq \frac{4}{\lambda^{3}} \|f\|_{\infty}^{4} \left[\sqrt{2\lambda} + \|C^{1/2} D(U - U_{n})\|_{L^{2}(H,\nu)}\right]$$
(3.17)

Moreover, taking into account (3.15) we find

$$\Lambda_2 \le \frac{1}{\lambda} \|f\|_{\infty}^2 \left(\int_H \text{Tr} \left[(CD^2\varphi_n)^2 \right] d\nu \right)^{1/2} \left(\int_H |C^{1/2}D\varphi_n|^4 d\nu \right)^{1/2}, \quad (3.18)$$

and

$$\Lambda_{3} \leq \frac{2}{\lambda} \| f \|_{\infty} \int_{H} |C^{1/2} D\varphi_{n}|^{3} |DU - DU_{n}| d\nu$$

$$\leq \frac{2}{\lambda} \| f \|_{\infty} \left(\int_{H} |D(U - U_{n})|^{4} d\nu \right)^{1/4} \left(\int_{H} |C^{1/2} D\varphi_{n}|^{4} d\nu \right)^{3/4}.$$
(3.19)

Substituting (3.17), (3.18) and (3.19) in (3.16), we find

$$\begin{split} &\int_{H} |C^{1/2} D\varphi_{n}|^{4} d\nu \leq \frac{4}{\lambda^{3}} \|f\|_{\infty}^{4} \left[\sqrt{2\lambda} + \|C^{1/2} D(U-U_{n}))\|_{L^{2}(H,\nu)}\right] \\ &+ \frac{1}{\lambda} \|f\|_{\infty}^{2} \left(\int_{H} \operatorname{Tr}\left[(C D^{2} \varphi_{n})^{2}\right] d\nu\right)^{1/2} \left(\int_{H} |C^{1/2} D\varphi_{n}|^{4} d\nu\right)^{1/2} \qquad (3.20) \\ &+ \frac{2}{\lambda} \|f\|_{\infty} \left(\int_{H} |D(U-U_{n})|^{4} d\nu\right)^{1/4} \left(\int_{H} |C^{1/2} D\varphi_{n}|^{4} d\nu\right)^{3/4}. \end{split}$$

Consequently, there exists a constant $\kappa_1 > 0$ depending on λ , $||f||_{\infty}$, and $||C^{1/2}DU||_{L^4(H,\nu)}$, but not on *n*, such that

$$\int_{H} |C^{1/2} D\varphi_n|^4 d\nu \le \kappa_1 \left(1 + \int_{H} \operatorname{Tr} \left[(C D^2 \varphi_n)^2 \right] d\nu \right).$$
(3.21)

Now we estimate the second hand side.

By (3.10), using the Hölder inequality, we find

$$\begin{split} &\frac{1}{2} \int_{H} \operatorname{Tr} \left[(CD^{2}\varphi_{n})^{2} \right] d\nu \leq 4 \|f\|_{\infty}^{2} \\ &+ \left(\int_{H} \operatorname{Tr} \left[(CD^{2}\varphi_{n})^{2} \right] d\nu \right)^{1/2} \left(\int_{H} |C^{1/2}D\varphi_{n}|^{2} |C^{1/2}DU_{n}|^{2} d\nu \right)^{1/2} \\ &+ 2 \|f\|_{\infty} \int_{H} |C^{1/2}D\varphi_{n}| |C^{1/2}DU_{n}| d\nu \\ &+ \left(\int_{H} \operatorname{Tr} \left[(CD^{2}\varphi_{n})^{2} \right] d\nu \right)^{1/2} \left(\int_{H} |C^{1/2}D\varphi_{n}|^{2} |C^{1/2}D(U-U_{n})|^{2} d\nu \right)^{1/2} \\ &+ 4 \|f\|_{\infty} \int_{H} |C^{1/2}D\varphi_{n}| |C^{1/2}D(U-U_{n})| d\nu \\ &+ 2 \int_{H} |C^{1/2}D\varphi_{n}|^{2} |C^{1/2}DU_{n}| |C^{1/2}D(U-U_{n})| d\nu. \end{split}$$

Using the Hölder inequality we find

$$\begin{split} &\frac{1}{2} \int_{H} \operatorname{Tr} \left[(CD^{2}\varphi_{n})^{2} \right] d\nu \leq 4 \|f\|_{\infty}^{2} \\ &+ \left(\int_{H} \operatorname{Tr} \left[(CD^{2}\varphi_{n})^{2} \right] d\nu \right)^{\frac{1}{2}} \left(\int_{H} |C^{1/2}D\varphi_{n}|^{4}d\nu \right)^{\frac{1}{4}} \left(\int_{H} |C^{1/2}DU_{n}|^{4}d\nu \right)^{\frac{1}{4}} \\ &+ 2 \|f\|_{\infty} \left(\int_{H} |C^{1/2}D\varphi_{n}|^{2}d\nu \right)^{\frac{1}{2}} \left(\int_{H} |C^{1/2}DU_{n}|^{2}d\nu \right)^{\frac{1}{2}} \\ &+ \left(\int_{H} \operatorname{Tr} \left[(CD^{2}\varphi_{n})^{2} \right] d\nu \right)^{\frac{1}{2}} \left(\int_{H} |C^{1/2}D\varphi_{n}|^{4}d\nu \int_{H} |C^{1/2}D(U-U_{n})|^{4}d\nu \right)^{\frac{1}{4}} \end{split}$$

$$+4\|f\|_{\infty} \left(\int_{H} |C^{1/2} D\varphi_{n}|^{2} d\nu\right)^{\frac{1}{2}} \left(\int_{H} |C^{1/2} D(U-U_{n})|^{2} d\nu\right)^{\frac{1}{2}} \\ +2\left(\int_{H} |C^{1/2} D\varphi_{n}|^{4} d\nu\right)^{\frac{1}{2}} \left(\int_{H} |C^{1/2} DU_{n}|^{4} d\nu \cdot \int_{H} |C^{1/2} D(U-U_{n})|^{4} d\nu\right)^{\frac{1}{4}}.$$

Therefore there exists $\kappa_2 > 0$ depending only on λ and $||f||_{\infty}$ such that

$$\int_{H} \operatorname{Tr} \left[(CD^{2}\varphi_{n})^{2} \right] d\nu
\leq \kappa_{2} \left[1 + \left(\int_{H} \operatorname{Tr} \left[(CD^{2}\varphi_{n})^{2} \right] d\nu \right)^{1/2} \left(\int_{H} |C^{1/2}D\varphi_{n}|^{4} d\nu \right)^{1/4} + \left(\int_{H} |C^{1/2}D\varphi_{n}|^{4} d\nu \right)^{1/2} \right].$$
(3.22)

Setting

$$x = \int_{H} |C^{1/2} D\varphi_n|^4 \, d\nu, \quad y = \int_{H} |C^{1/2} D\varphi_n|^4 \, d\nu$$

we find by (3.21) and (3.22)

$$y \le \kappa_1(1+x), \ x \le \kappa_2(1+y^{1/2}+x^{1/2}y^{1/4}).$$

This gives the conclusion.

Now we are ready to prove

Theorem 3.2. Assume that Hypotheses 1.1 and 1.2 holds. Then \mathring{K} is essentially self-adjoint in $L^2(H, v)$. Moreover, denoting by K its closure, we have

$$\int_{H} K\varphi \psi d\nu = -\frac{1}{2} \int_{H} \langle CD\varphi, D\psi \rangle d\nu. \qquad \varphi, \psi \in D(K).$$
(3.23)

Proof. Since \mathring{K} is symmetric, it is sufficient to show that the image of $\lambda - \mathring{K}$ is dense in $L^2(H, \nu)$, see [18].

Let $f \in C_b^2(H)$ and let φ_n be the solution of (3.5). By the Hölder inequality, we have

$$\left[\int_{H} |\langle C^{1/2}(DU - DU_n), C^{1/2}D\varphi_n\rangle|^2 d\nu\right]^2$$

$$\leq \int_{H} |C^{1/2}(DU - DU_n)|^4 d\nu \int_{H} |C^{1/2}D\varphi_n|^4 d\nu$$

Taking into account Proposition 3.1 we obtain

$$\lim_{n\to\infty}\int_{H}|\langle C^{1/2}(DU-DU_n),C^{1/2}D\varphi_n\rangle|^2d\nu=0,$$

so that

$$\lambda \varphi_n - \mathring{K} \varphi_n \to f$$

in $L^2(H, \nu)$ and, since $C_b^2(H)$ is dense in $L^2(H, \nu)$, it follows that \mathring{K} is essentially self-adjoint. The last statement follows from (3.1).

Remark 3.3. By (3.23) for any $\varphi \in L^2(H, \nu)$ we have

$$\int_{H} |e^{Kt}\varphi|^2 \, d\nu + \int_0^t ds \int_{H} |C^{1/2} D e^{Kt}\varphi|^2 \, d\nu = \int_{H} |\varphi|^2 \, d\nu.$$
(3.24)

Using (3.24) it is not difficult to show

$$\lim_{t \to \infty} e^{Kt} \varphi = \int_H \varphi \, d\nu \qquad \text{in } L^2(H, \nu).$$

Thus the measure ν is strongly mixing.

Remark 3.4. Let *K* be the closure of \mathring{K} . Then obviously $D(L) \cap C_b^2(H)$ is a core for *K*. Since $\mathscr{E}_A(H) \subset D(L_2) \cap C_b^2(H)(^4)$ and is dense in $D(L_2)$, endowed with the graph norm, then it is not difficult to show that $\mathscr{E}_A(H)$ is a core for *K*.

Remark 3.5. Let $F \in L^4(H, \nu; H)$ and consider the operator

$$\mathring{N}\varphi = L\varphi - \langle F(x), D\varphi \rangle, \ \varphi \in C_b^1(H) \cap D(L).$$
(3.25)

Assume that there exists a Borel measure v on H such that

$$\int_{H} \mathring{N}\varphi\varphi d\nu = -\frac{1}{2} \int_{H} |C^{1/2}D\varphi|^2 d\nu.$$
(3.26)

Then all previous considerations can be repeated and we can conclude that \mathring{N} is closable on $L^2(H, \nu)$ and its closure is *m*-dissipative.

4. Application to stochastic quantization

Let S be the square $[0, 2\pi]^2$. Let us consider the differential stochastic equation

$$\begin{cases} dX = \frac{1}{2} \ ((\Delta - 1)X + \sigma X^n) \, dt + dW(t), \\ X(0) = x, \end{cases}$$
(4.1)

where *n* is odd, $\sigma > 0$, and *W* is the cylindrical Wiener process on $L^2(\mathbb{S})$, W(t) is defined as

$$W(t)(\xi) = \sum_{h \in \mathbb{Z}^2} e_h(\xi) \eta_h,$$

where $\{\eta_h\}$ is a sequence of gaussian random variables $\mathcal{N}(0, 1)$ mutually independent and $\{e_h\}$ is the complete orthonormal system in $L^2(\mathbb{S})$ defined by

$$e_h(\xi) = \frac{1}{2\pi} e^{i\langle h, \xi \rangle}, \qquad h = (h_1, h_2) \in \mathbb{Z}^2, \ \xi \in \mathbb{S}.$$

Moreover Δ is the realization of the Laplace operator with 2π -periodic boundary conditions.

⁴ Here we denote by $D(L_2)$ the domain of Ornstein-Uhlenbeck operator in $L^2(H, \mu)$

Equation (4.1) describes a gradient system having formally the Gibbs invariant probability measure:

$$\nu(dx) = Z^{-1} e^{\frac{1}{n+1} \int_{\mathbb{S}} x^{n+1}(\xi) d\xi} \mathcal{N}(0, C)(dx)$$

where $C = (1 - \Delta)^{-1}$. Notice that this definition is not meaningful since

Tr
$$C = \sum_{h \in \mathbb{Z}^2} \frac{1}{1 + |h|^2} = +\infty.$$

Equation (4.1) takes in account an interacting field. The corresponding free field is described by

$$\begin{cases} dZ = \frac{1}{2} (\Delta - 1) Z dt + dW(t), \\ Z(0, x) = x, \end{cases}$$
(4.2)

so that

$$Z(t,x) = e^{\frac{1}{2}(\Delta-1)t}x + \int_0^t e^{\frac{1}{2}(\Delta-1)(t-s)} dW_s$$

Since

$$\mathbb{E}|Z(t,0)|^2 = \operatorname{Tr} [C(1-e^{(\Delta-1)t})] = +\infty.$$

then Z(t) does not live in $L^2(\mathbb{S})$. However it lives in $H^{-1}(\mathbb{S})$ defined as the completion of $L^2(\mathbb{S})$ with respect to the inner product

$$\langle x, y \rangle_{-1} = \sum_{h \in \mathbb{Z}^2} (1 + |h|^2)^{-1} \langle x, e_h \rangle \overline{\langle y, e_h \rangle}.$$

We have in fact

$$\mathbb{E}|Z(t,0)|_{-1}^2 = \operatorname{Tr} \left[C^2(1-e^{(\Delta-1)t})\right] < +\infty,$$

as easily checked.

Now we interpret equation (4.1) as an equation in $H^{-1}(\mathbb{S})$. For this we set $W_1(t) = C^{-1/2}W(t)$ so that W_1 is a cylindrical Wiener process in $H^{-1}(\mathbb{S})$ and replace the Gibbs measure by

$$\nu(dx) = Z^{-1} \exp\left(\frac{1}{n+1} \int_{\mathbb{S}} : x^{n+1}(\xi) : d\xi\right) \ \mathcal{N}(0, C)(dx), \tag{4.3}$$

where : $x^{n+1}(\xi)$: is the Wick product that we recall below and Z is the normalization constant. In this way (4.1) is replaced by

$$\begin{cases} dX = \frac{1}{2} \left((\Delta - 1)X + \sigma C : X^n : \right) dt + C^{1/2} dW_1(t), \\ X(0) = x. \end{cases}$$
(4.4)

We notice that the factor C in front to : X^n : is necessary to make (4.4) a gradient system (having (4.3) as invariant measure).

We fix now $n \in \mathbb{N}$ and recall the definition of the Wick monomial of order *n*:

$$: x^n := G_n(x) : H^{-1}(\mathbb{S}) \to H^{-1}(\mathbb{S}).$$

where $G_n(x)$ is the limit in $L^2(H^{-1}(\mathbb{S}), \mu, H^{-1}(\mathbb{S}))$ of the functions $G_{n,N}(x)$ defined as

$$G_{n,N}(x)(\xi) = \sqrt{n!}\rho_N^n \, \mathbb{H}_n\left(\frac{1}{\rho_N}\sum_{|h|\leq N} \langle x, e_h \rangle e_h(\xi)\right),$$

where \mathbb{H}_n denotes the Hermite polynomial of order *n* (see [10]) and

$$\rho_N^2 = \frac{1}{(2\pi)^2} \sum_{|h| \le N} \frac{1}{1 + |h|^2}.$$

Moreover, setting

$$U_n(x) = \frac{1}{n+1} \int_{\mathbb{S}} G_{n+1}(x)(\xi) d\xi,$$

we have that $U_n \in W_C^{1,p}(H^{-1}(\mathbb{S}), \mu)$ for all $p \ge 1$, and $DU_n(x) = G_n(x)$. Also if *n* is a positive odd integer

$$e^{-pU_n(x)}$$
 is in $L^1(H^{-1}(\mathbb{S}), \mu)$

for all $p \ge 1$. For a proof of these results we refer to Simon [22]; see also [10].

Finally fix $\sigma > 0$ and a positive odd integer *n*. Then consider the following linear operator in $H = H^{-1}(\mathbb{S})$:

$$\mathring{K}\varphi = \frac{1}{2} \operatorname{Tr} \left[CD^2\varphi \right] - \frac{1}{2} \langle C^{-1}x, D\varphi \rangle_{-1} - \frac{\sigma}{2} \langle CG_n(x), D\varphi \rangle_{-1}.$$

Setting $A = -\frac{1}{2} C^{-1}$, $U = -\frac{\sigma}{2} U_n$, we can apply Theorem 3.2 and conclude that \mathring{K} is essentially selfadjoint on $L^2(H, \nu)$.

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