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On the support of a hypoelliptic diffusion on the Heisenberg group

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Abstract. We provide an elementary proof of the support of the law of a hypoelliptic Brownian motion on the Heisenberg group \mathbb{H} . We consider a control norm associated to left-invariant vector fields on \mathbb{H} , and describe the support in terms of the space of finite energy horizontal curves.

1. Introduction

The purpose of this paper is to describe the support of the law of a hypoellptic diffusion g_t on the Heisenberg group \mathbb{H} . The novelty of this paper is the norm on the path space $W_0(\mathbb{H})$ of \mathbb{H} -valued continuous curves starting at the identity that is used for the support. The group \mathbb{H} is the simplest example of a Carnot group and it comes with a natural left-invariant distance, the Carnot-Carathéodory distance d_{cc} . This is the control distance associated to the left-invariant vector fields on \mathbb{H} , see Definition 2.4. The corresponding control norm d_c is defined as $d_c(x) := d_{cc}(x, e)$, where $e \in \mathbb{H}$ is the identity. Our main result is to prove a support theorem for the hypoelliptic Brownian motion g_t with respect to the norm $\max_{0 \leq t \leq 1} d_c(\gamma(t))$ for $\gamma \in W_0(\mathbb{H})$. As pointed out in Remark 2.15 it is sufficient to describe the support with respect to some equivalent norms.

The support of a diffusion was first studied by Stroock and Varadhan (1972), which we now describe briefly. Let X_t be an \mathbb{R}^d -valued diffusion which is solution to the stochastic differential equation

$$dX_t = \sigma\left(t, X_t\right) \circ dW_t + b\left(t, X_t\right) dt, \quad X_0 = 0, \tag{1.1}$$

where $\sigma(t, x)$ is a $d \times \ell$ matrix whose entries are functions of $(t, x) \in [0, 1] \times \mathbb{R}^d$, and b(t, x) is a vector in \mathbb{R}^d , and W_t is an ℓ -dimensional Brownian motion, and $\circ d$ denotes the stochastic differential in Stratonovich's sense. We can view the process $\{X_t\}_{0 \leq t \leq 1}$ as a $W_0(\mathbb{R}^d)$ -valued random variable, where $W_0(\mathbb{R}^d)$ is the space of \mathbb{R}^d -valued continuous paths starting at zero. Let μ be the law of $\{X_t\}_{0 \leq t \leq 1}$ and S_{μ} its support. If H denotes the subset of $W_0(\mathbb{R}^d)$ consisting of absolutely continuous paths, then to any $\phi \in H$ one can associate a deterministic path x_{ϕ} as being the solution to the ordinary

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differential equation

$$x'_{\phi}(t) = \sigma(t, x_{\phi}(t)) \phi'(t) dt + b(t, x_{\phi}(t)) dt,$$
(1.2)
$$x_{\phi}(0) = 0.$$

We follow Kunita (1978) and refer to solutions to (1.2) as controlled systems. Then

$$S_{\mu} = \overline{\{x_{\phi}, \ \phi \in H\}}^{\infty}, \tag{1.3}$$

where the closure is taken in the uniform topology in $W_0(\mathbb{R}^d)$.

Note that the hypoelliptic Brownian motion g_t can be viewed as an \mathbb{R}^3 -valued stochastic process. This is not a Gaussian process and it satisfies the stochastic differential equation

$$\begin{pmatrix} dg_1(t) \\ dg_2(t) \\ dg_3(t) \end{pmatrix} = \begin{pmatrix} 1 & 0 \\ 0 & 1 \\ -\frac{1}{2}B_2(t) & \frac{1}{2}B_1(t) \end{pmatrix} \begin{pmatrix} dB_1(t) \\ dB_2(t) \end{pmatrix}, \quad g_0 = (0, 0, 0).$$

Stroock-Varadhan original support theorem (1.3) was proven under the assumption that σ is C^2 in space and C^1 in time, bounded together with its partial derivatives of order one and two, and b is globally Lipschitz and bounded. In a series of papers by Gyöngy (1988a,b, 1994), and by Gyöngy and Pröhle (1990) (1.3) is proven for processes driven by continuous semi-martingales under milder assumptions on σ and b. In particular, (1.3) for the law of $\{g_t\}_{0 \leq t \leq 1}$ with respect to the uniform norm $\max_{0 \leq t \leq 1} |\gamma(t)|_{\mathbb{R}^3}$ on $W_0(\mathbb{R}^3)$ follows from Gyöngy and Pröhle (1990, Theorem 3.1). Moreover, (1.3) for the hypoelliptic Brownian motion can also be proven by rough path theory and continuity of the Lyons-Itô map, that is, the solution map of a rough differential equation Friz and Victoir (2010, Section 10.3, Section 13.7). In the current paper we prove (1.3) for the hypoelliptic Brownian motion on the Heisenberg group. Differently from Gyöngy and Pröhle (1990), we replace the Euclidean norm in \mathbb{R}^3 by the control norm d_c , which is a more natural norm and it is consistent with the geometry of the Heisenberg group. Our proof does not rely on rough path theory and it is based on a time change argument.

We mention that (1.3) for diffusion processes on Hilbert spaces is proven in Aida (1990); Gyöngy (1989). We also mention that in Ledoux et al. (2002) a rough paths approach is used, and a support theorem in the *p*-variational topology is proven.

One can ask under what condition the closure in (1.3) coincides with the whole path space $W_0(\mathbb{R}^d)$. This question has been addressed in Kunita (1978), where the author gives nearly necessary and sufficient conditions for

$$W_0(\mathbb{R}^d) = \overline{\{x_\phi, \ \phi \in H\}}^{\infty}$$
(1.4)

to hold. We prove (1.4) for the hypoelliptic Brownian motion. Our proof is explicit and it relies on the group structure of \mathbb{H} .

The main results of this paper are contained in Theorem 2.14, where we prove (1.3) and (1.4) for the hypoelliptic Brownian motion on the Heisenberg group. More precisely, if $H(\mathbb{H})$ denotes the set of finite energy horizontal curves, then we prove that

$$W_0(\mathbb{H}) = \overline{H(\mathbb{H})}^{d_c} = \mathcal{S}_{\mu}, \qquad (1.5)$$

where μ is the law of $\{g_t\}_{0 \leq t \leq 1}$ and the closure is taken with respect to the norm $\max_{0 \leq t \leq 1} d_c(\gamma(t))$ for $\gamma \in W_0(\mathbb{H})$.

First, we show that $S_{\mu} \subset \overline{H(\mathbb{H})}^{d_c}$ by constructing a family of stochastic processes $\{g_{\delta}\}_{\delta>0}$ that approximates g in the sense that the law μ_{δ} of g_{δ} converges weakly to the law μ of g. We further study relations between the measures μ_{δ} and μ , and prove that they are singular. In particular we show that the space $H(\mathbb{H})$ of finite energy horizontal curves has μ -measure zero. The space $H(\mathbb{H})$ can then be viewed as a Cameron-Martin space in a non-Gaussian setting by (1.5) and since $\mu(H(\mathbb{H})) = 0$. We then prove that $\overline{H(\mathbb{H})}^{d_c} \subset S_{\mu}$ using Theorem 3.6 and the explicit form of the process g_t . Namely, $g_t = (B_t, A_t)$, where B_t is a two-dimensional standard Brownian motion and A_t is the corresponding Lévy's stochastic area. Our proof relies on the classical identity $A_t = b_{\tau(t)}$, where b_t is a one-dimensional standard Brownian motion independent of B_t , and $\tau(t)$ is a stopping time.

The paper is organized as follows. In Section 2 we describe the Heisenberg group \mathbb{H} , and the corresponding control norm and hypoelliptic Brownian motion, and state the main result of the paper Theorem 2.14. Section 3 then contains the proof of Theorem 2.14.

2. The setting and the main result

2.1. Heisenberg group as Lie group. The Heisenberg group \mathbb{H} as a set is $\mathbb{R}^3 \cong \mathbb{R}^2 \times \mathbb{R}$ with the group multiplication given by

$$(\mathbf{v}_1, z_1) \cdot (\mathbf{v}_2, z_2) := \left(x_1 + x_2, y_1 + y_2, z_1 + z_2 + \frac{1}{2}\omega \left(\mathbf{v}_1, \mathbf{v}_2 \right) \right),$$
where $\mathbf{v}_1 = (x_1, y_1), \mathbf{v}_2 = (x_2, y_2) \in \mathbb{R}^2,$ and
 $\omega : \mathbb{R}^2 \times \mathbb{R}^2 \longrightarrow \mathbb{R}, \quad \omega \left(\mathbf{v}_1, \mathbf{v}_2 \right) := x_1 y_2 - x_2 y_1.$

The identity in \mathbb{H} is e = (0, 0, 0) and the inverse is given by $(\mathbf{v}, z)^{-1} = (-\mathbf{v}, -z)$. The Lie algebra of \mathbb{H} can be identified with the space $\mathbb{R}^3 \cong \mathbb{R}^2 \times \mathbb{R}$ with the Lie bracket defined by

$$\left[\left(\mathbf{a}_{1}, c_{1}
ight), \left(\mathbf{a}_{2}, c_{2}
ight)
ight] = \left(0, \omega \left(\mathbf{a}_{1}, \mathbf{a}_{2}
ight)
ight)$$

The set $\mathbb{R}^3 \cong \mathbb{R}^2 \times \mathbb{R}$ with this Lie algebra structure will be denoted by \mathfrak{h} .

Let us now recall some basic notation for Lie groups. Suppose G is a Lie group, then the left and right multiplication by an element $k \in G$ are denoted by

$$L_k: G \longrightarrow G, \qquad \qquad g \longmapsto k^{-1}g$$
$$R_k: G \longrightarrow G, \qquad \qquad g \longmapsto gk$$

Recall that the tangent space T_eG can be identified with the Lie algebra \mathfrak{g} of left-invariant vector fields on G, that is, vector fields X on G such that $dL_k \circ X = X \circ L_k$, where dL_k is the differential of L_k . More precisely, if A is a vector in T_eG , then we denote by $\tilde{A} \in \mathfrak{g}$ the (unique) left-invariant vector field such that $\tilde{A}(e) = A$. A left-invariant vector field is determined by its value at the identity, namely, $\tilde{A}(k) = dL_k \circ \tilde{A}(e)$.

For the Heisenberg group the differential of left and right multiplication can be described explicitly as follows.

Proposition 2.1. Let $k = (k_1, k_2, k_3) = (\mathbf{k}, k_3)$ and $g = (g_1, g_2, g_3) = (\mathbf{g}, g_3)$ be two elements in \mathbb{H} . Then, for every $v = (v_1, v_2, v_3) = (\mathbf{v}, v_3)$ in $T_g \mathbb{H}$, the differentials of the left and right multiplication are given by

$$dL_{k}: T_{g}\mathbb{H} \longrightarrow T_{k^{-1}g}\mathbb{H},$$

$$dR_{k}: T_{g}\mathbb{H} \longrightarrow T_{gk}\mathbb{H},$$

$$dL_{k}(v) = \left(v_{1}, v_{2}, v_{3} + \frac{1}{2}\omega(\mathbf{v}, \mathbf{k})\right),$$

$$dR_{k}(v) = \left(v_{1}, v_{2}, v_{3} + \frac{1}{2}\omega(\mathbf{v}, \mathbf{k})\right).$$
(2.1)

2.2. Heisenberg group as a sub-Riemannian manifold. The Heisenberg group \mathbb{H} is the simplest non-trivial example of a sub-Riemannian manifold. We define X, Y and Z as the unique left-invariant

vector fields satisfying $X_e = \partial_x$, $Y_e = \partial_y$ and $Z_e = \partial_z$, that is,

$$X = \partial_x - \frac{1}{2}y\partial_z,$$

$$Y = \partial_y + \frac{1}{2}x\partial_z,$$

$$Z = \partial_z.$$

Note that the only non-zero Lie bracket for these left-invariant vector fields is [X, Y] = Z, so the vector fields $\{X, Y\}$ satisfy Hörmander's condition. We define the *horizontal distribution* as $\mathcal{H} := \operatorname{span} \{X, Y\}$ fiberwise, thus making \mathcal{H} a sub-bundle in the tangent bundle $T\mathbb{H}$. To finish the description of the Heisenberg group as a sub-Riemannian manifold we need to equip the horizontal distribution \mathcal{H} with an inner product. For any $p \in \mathbb{H}$ we define the inner product $\langle \cdot, \cdot \rangle_{\mathcal{H}_p}$ on \mathcal{H}_p so that $\{X(p), Y(p)\}$ is an orthonormal (horizontal) frame at any $p \in \mathbb{H}$. Vectors in \mathcal{H}_p will be called *horizontal*, and the corresponding norm will be denoted by $\|\cdot\|_{\mathcal{H}_p}$.

In addition, Hörmander's condition ensures that a natural sub-Laplacian on the Heisenberg group

$$\Delta_{\mathcal{H}} = X^2 + Y^2 \tag{2.2}$$

is a hypoelliptic operator by Hörmander (1967). We recall now another important object in sub-Riemannian geometry, namely, horizontal curves.

Definition 2.2. An absolutely continuous path $t \mapsto \gamma(t) \in \mathbb{H}$, for a.e. $t \in [0, 1]$ is said to be horizontal if $\gamma'(t) \in \mathcal{H}_{\gamma(t)}$ for a.e. t, that is, the tangent vector to $\gamma(t)$ at every point $\gamma(t)$ is horizontal. Equivalently we can say that γ is horizontal if $c_{\gamma}(t) := dL_{\gamma(t)}(\gamma'(t)) \in \mathcal{H}_e$ for a.e. t.

Note that for $\gamma(t) = (\mathbf{x}(t), z(t))$ we have

$$c_{\gamma}(t) = dL_{\gamma(t)}\left(\gamma'(t)\right) = \left(\mathbf{x}'(t), z'(t) - \frac{1}{2}\omega(\mathbf{x}(t), \mathbf{x}'(t))\right), \qquad (2.3)$$

where we used Proposition 2.1. Equation (2.3) can be used to characterize horizontal curves in terms of the components as follows. The curve γ is horizontal if and only if, for a.e. $0 \leq t \leq 1$

$$z'(t) - \frac{1}{2}\omega(\mathbf{x}(t), \mathbf{x}'(t))) = 0.$$
(2.4)

Definition 2.3. We say that a horizontal curve $t \mapsto \gamma(t) \in \mathbb{H}$, $0 \leq t \leq 1$ has finite energy if

$$\|\gamma\|_{H(\mathbb{H})}^{2} := \int_{0}^{1} |c_{\gamma}(s)|_{\mathcal{H}_{e}}^{2} ds = \int_{0}^{1} |dL_{\gamma(s)}(\gamma'(s))|_{\mathcal{H}_{e}}^{2} ds < \infty.$$
(2.5)

We denote by $H(\mathbb{H})$ the space of finite energy horizontal curves starting at the identity. The inner product corresponding to the norm $\|\cdot\|_{H(\mathbb{H})}$ is denoted by $\langle \cdot, \cdot \rangle_{H(\mathbb{H})}$. Note that the Heisenberg group as a sub-Riemannian manifold comes with a natural left-invariant distance.

Definition 2.4. For any $x, y \in \mathbb{H}$ the Carnot-Carathéodory distance is defined as

$$d_{cc}(x,y) := \inf \left\{ \int_0^1 |c_\gamma(s)|_{\mathcal{H}_e} ds, \\ \gamma : [0,1] \longrightarrow \mathbb{H}, \gamma(0) = x, \gamma(1) = y, \gamma \text{ is horizontal} \right\}.$$

Another consequence of Hörmander's condition for left-invariant vector fields X, Y and Z is that we can apply the Chow–Rashevskii theorem. As a result, given two points in \mathbb{H} there exists a horizontal curve connecting them, and therefore the Carnot-Carathéodory distance is finite on \mathbb{H} . The Carnot-Carathéodory distance in Definition 2.4 is the control distance associated to the vector fields X, Y, and Z on \mathbb{H} (Bonfiglioli et al., 2007, Definition 5.2.2). The control norm $d_c : \mathbb{H} \to \mathbb{R}$ is then defined as $d_c(x) = d_{cc}(x, e)$. Note that

$$d_c(y^{-1}x) = d_{cc}(x,y), (2.6)$$

by left-invariance of X, Y, and Z, and the definition of d_{cc} . The control norm is an example of a homogeneous norm.

Definition 2.5. Let $\rho : \mathbb{H} \to [0, \infty)$ be a continuous function with respect to the Euclidean topology. Then ρ is a homogeneous norm if it satisfies the following properties

$$\rho(\delta_{\lambda}(x)) = \lambda \rho(x), \text{ for every } \lambda > 0, \text{ and } x \in \mathbb{H},$$

 $\rho(x) = 0 \text{ if and only if } x = e,$

where $\delta_{\lambda}(x) := (\lambda x_1, \lambda x_2, \lambda^2 x_3).$

If ρ_1 and ρ_2 are two homogeneous norms, then there exists a constant c > 0 such that

$$c^{-1}\rho_1(x) \leqslant \rho_2(x) \leqslant c\rho_1(x), \tag{2.7}$$

for every $x \in \mathbb{H}$, Bonfiglioli et al. (2007, Proposition 5.1.4). We consider the following homogeneous norm

$$|x| := \left((x_1^2 + x_2^2)^2 + x_3^2 \right)^{\frac{1}{4}}, \qquad (2.8)$$

for every $x = (x_1, x_2, x_3) \in \mathbb{H}$. By (2.6) and (2.7) it follows that

$$c^{-1}d_{cc}(x,y) \leqslant |y^{-1}x| \leqslant cd_{cc}(x,y),$$

for any $x, y \in \mathbb{H}$.

Finally, we need to describe a hypoelliptic Brownian motion with values in \mathbb{H} .

Definition 2.6. An \mathbb{H} -valued Markov process g_t is called a hypoelliptic Brownian motion if its generator is the sub-Laplacian $\frac{1}{2}\Delta_{\mathcal{H}}$ defined by Equation (2.2).

One can write a stochastic differential equation for g_t . This form is the standard stochastic differential equation for a Lie group-valued Brownian motion, namely,

$$dL_{g_t}(dg_t) = (dB_1(t), dB_2(t), 0),$$

$$g_0 = e,$$

where $B_t = (B_1(t), B_2(t))$ is a standard two-dimensional Brownian motion. An explicit solution is given by

$$g_t := (B_t, A_t), \tag{2.9}$$

where $A_t := \frac{1}{2} \int_0^t \omega (B_s, dB_s)$ is the Levy's stochastic area. Note that we used the Itô integral rather than the Stratonovich integral. However, these two integrals are equal since the symplectic form ω is skew-symmetric, and therefore Lévy's stochastic area functional is the same for both integrals.

Notation 2.7. Throughout the paper we fix a filtered probability space $(\Omega, \mathcal{F}, \mathcal{F}_t, \mathbb{P})$. We denote the expectation under \mathbb{P} by \mathbb{E} .

2.3. *The Wiener meaure.* We recall here the definition of Wiener measure, and collect some notations that will be used throughout the paper.

Notation 2.8 (Topology on \mathbb{H}). Let $|\cdot|$ be the homogeneous norm in (2.8). We consider the topology on \mathbb{H} whose open balls centered at the identity are $\{x \in \mathbb{H}, |x| < r\}$.

Note that by (2.7) all homogeneous norms induce the same topology.

Notation 2.9 (Standard Wiener space). We denote by $W_0(\mathbb{R}^n)$ the space of \mathbb{R}^n -valued continuous functions on [0,1] starting at 0. This space comes with the norm

$$\|h\|_{W_0(\mathbb{R}^n)} := \max_{0 \le t \le 1} |h(t)|_{\mathbb{R}^n}, \quad h \in W_0(\mathbb{R}^n),$$

and the associated distance $d_{W_0(\mathbb{R}^n)}(h,k) = \max_{0 \leq t \leq 1} |h(t) - k(t)|_{\mathbb{R}^n}$, where $|\cdot|_{\mathbb{R}^n}$ is the Euclidean norm.

Definition 2.10 (Wiener space over \mathbb{H}). The Wiener space over \mathbb{H} , denoted by $W_0(\mathbb{H})$, is the space of \mathbb{H} -valued continuous functions on [0, 1] starting at identity in \mathbb{H} .

Once a homogeneous norm ρ on \mathbb{H} is fixed, one can introduce a topology on $W_0(\mathbb{H})$ in the following way. We endow $W_0(\mathbb{H})$ with the following norm

$$\|\eta\|_{\rho} := \max_{0 \le t \le 1} \rho(\eta(t)), \quad \eta \in W_0\left(\mathbb{H}\right),$$

and the associated distance is $\max_{0 \leq t \leq 1} \rho(\eta(t)^{-1}\gamma(t))$ for any $\eta, \gamma \in W_0(\mathbb{H})$.

Definition 2.11. Let $W_0(\mathbb{H})$ be the Wiener space over \mathbb{H} , and $\{g_t\}_{0 \leq t \leq 1}$ be the hypoelliptic Brownian motion defined by (2.9). We call its law the horizontal Wiener measure and we denote it by μ .

The process g_t can be viewed as a $W_0(\mathbb{H})$ -valued random variable, that is,

 $g: \Omega \longrightarrow W_0(\mathbb{H}), \qquad \omega \mapsto \{t \mapsto g_t(\omega)\}.$

The measure μ is then given by $\mu(E) = \mathbb{P}(g^{-1}(E)) = \mathbb{P}(g \in E)$ for any Borel set E in $W_0(\mathbb{H})$. We denote the support of μ by S_{μ} , that is, S_{μ} is the smallest closed subset of $W_0(\mathbb{H})$ having μ -measure one.

Remark 2.12. Note that even though the hypoelliptic Brownian motion g_t is an \mathbb{R}^3 -valued stochastic process, it is not a Gaussian process, and its law μ is not a Gaussian measure on $W_0(\mathbb{H})$. Moreover, contrary to the Euclidean case, the space $W_0(\mathbb{H})$ is not a Banach space. It is easy to see that the space $W_0(\mathbb{H})$ is closed under the norm $\max_{0 \leq t \leq 1} \rho(\gamma(t))$ for $\gamma \in W_0(\mathbb{H})$, where ρ is a homogeneous norm on \mathbb{H} , but $W_0(\mathbb{H})$ is not a linear space.

Let us denote by π the projection

$$\pi: W_0(\mathbb{H}) \longrightarrow W_0(\mathbb{R}^2), \ \pi(\gamma) = (\gamma_1, \gamma_2), \tag{2.10}$$

for any $\gamma = (\gamma_1, \gamma_2, \gamma_3) \in W_0(\mathbb{H}).$

Remark 2.13. Let $\phi = (\phi_1, \phi_2, \phi_3) \in H(\mathbb{H})$ be a finite energy horizontal curve as in Definition 2.3. Then $\pi(\phi)$ is in the Cameron-Martin space on \mathbb{R}^2 , that is, $\pi(\phi)$ is an absolutely continuous \mathbb{R}^2 -valued curve starting at zero such that

$$\int_0^1 |\pi(\phi)'(s)|_{\mathbb{R}^2}^2 ds < \infty.$$

2.4. Main result. Now we have all the ingredients needed to state the main result of this paper, that is, we describe the support of the Wiener measure for the hypoelliptic Brownian motion g_t in terms of horizontal paths.

Theorem 2.14. Let $W_0(\mathbb{H})$ be the Wiener space over \mathbb{H} , and μ be the horizontal Wiener measure on $W_0(\mathbb{H})$, and $H(\mathbb{H})$ be the space of horizontal curves with finite energy. Then

$$\mathcal{S}_{\mu} = \overline{H\left(\mathbb{H}\right)}^{d_{c}} = W_{0}\left(\mathbb{H}\right),$$

where the closure is taken with respect to the norm $\max_{0 \le t \le 1} d_c(\gamma(t))$, for $\gamma \in W_0(\mathbb{H})$, and d_c is the control norm induced by the Carnot-Carathéodory distance.

Remark 2.15. It is enough to prove Theorem 2.14 for the homogeneous norm $|\cdot|$ given by (2.8). Indeed, if ρ_1 and ρ_2 are two homogeneous norms on \mathbb{H} and $\|\cdot\|_{\rho_1}$, $\|\cdot\|_{\rho_2}$ denote the corresponding norms on $W_0(\mathbb{H})$, that is,

$$\|\gamma\|_{\rho_i} := \max_{0 \le t \le 1} \rho_i(\gamma(t)), \quad \gamma \in W_0(\mathbb{H}).$$

Then

$$\overline{H\left(\mathbb{H}\right)}^{\rho_{1}}=\overline{H\left(\mathbb{H}\right)}^{\rho_{2}},$$

since $\|\cdot\|_{\rho_1}$, $\|\cdot\|_{\rho_2}$ are equivalent by (2.7), and hence Theorem 2.14 holds for any homogeneous norm as soon as it holds for one norm.

Let us denote by $H(\mathbb{H})$ the closure of $H(\mathbb{H})$ with respect to the norm $\max_{0 \leq t \leq 1} \rho(\gamma(t)), \gamma \in W_0(\mathbb{H})$, and for ρ a homogeneous norm on \mathbb{H} . By Remark 2.15, $\overline{H(\mathbb{H})}$ is independent of ρ .

3. Proof of Theorem 2.14

We will divide the proof of Theorem 2.14 in two steps. First, we introduce a family of processes that approximates $\{g_t\}_{0 \le t \le 1}$. This is used in Corollary 3.4 to show that the support S_{μ} is contained in $\overline{H(\mathbb{H})}$. The reverse inclusion is proven in Corollary 3.7 which follows from Theorem 3.6. In Proposition 3.10 we prove that $\overline{H(\mathbb{H})} = W_0(\mathbb{H})$, which concludes the proof of Theorem 2.14.

3.1. Approximation of the hypoelliptic Brownian motion. The aim of this step is to show that the support S_{μ} of the law of $\{g_t\}_{0 \leq t \leq 1}$ is contained in $\overline{H(\mathbb{H})}$. This will be accomplished by constructing a horizontal piecewise approximation $g_{\delta}(t)$ of g_t such that $\mu_{\delta} \to \mu$ weakly, where μ_{δ} is the law of $g_{\delta}(t)$. Different approximations of a Brownian motion have been extensively studied over the decades, see for example Wong and Zakai (1965), Kunita (1974), Nakao and Yamato (1978), Ikeda et al. (1977/78), and Ikeda and Watanabe (1989, Chapter 6, Section 7) for more details. We are not able to refer to all the vast literature on the subject, but we mentioned some results which are closer and more relevant to the techniques we use in this paper.

Let $\{B_{\delta}\}_{\delta>0}$ be an approximation of a two-dimensional Brownian motion, that is,

$$\mathbb{E}\left[\max_{0\leqslant t\leqslant 1}|B_{\delta}(t)-B_{t}|_{\mathbb{R}^{2}}^{2}\right]\longrightarrow 0, \quad \text{as } \delta\to 0,$$
(3.1)

such that

$$\mathbb{E}\left[\max_{0\leqslant t\leqslant 1}A_{\delta}(t)-A_{t}|_{\mathbb{R}}^{2}\right]\longrightarrow 0 \quad \text{as } \delta\to 0,$$
(3.2)

where

$$A_{\delta}(t) := \frac{1}{2} \int_0^t \left(B_{1,\delta}(s) B'_{2,\delta}(s) - B_{2,\delta}(s) B'_{1,\delta}(s) \right) ds.$$
(3.3)

Let f_1 and f_2 be differentiable functions on [0,1] such that $f_i(0) = 0$ and $f_i(1) = 1$ for i = 1, 2. Set

$$B_{i,\delta}(t) := B_i(k\delta) + f_i\left(\frac{t-k\delta}{\delta}\right) \left(B_i(k\delta+\delta) - B_i(k\delta)\right), \ k\delta \le t < (k+1)\delta, \tag{3.4}$$

then by Ikeda and Watanabe (1989, Theorem 7.1) the family $\{B_{\delta}\}_{\delta>0}$ satisfies (3.1) and (3.2). Let us define now a sequence of processes $g_{\delta}(t)$ on \mathbb{H} .

Definition 3.1. Let $B_{\delta}(t)$ be an approximation of a two-dimensional Brownian motion satisfying (3.2). For each δ , t, and ω we set

$$g_{\delta}(t) = \left(g_{1,\delta}(t), g_{2,\delta}(t), g_{3,\delta}(t)\right),$$

where

$$g_{1,\delta}(t) = B_{1,\delta}(t)$$

$$g_{2,\delta}(t) = B_{2,\delta}(t)$$

$$g_{3,\delta}(t) = A_{\delta}(t).$$
(3.5)

Let $C_p^2(\mathbb{R}^2)$ be the space of piecewise continuously twice differentiable curves in \mathbb{R}^2 starting at zero, and set

$$H_p\left(\mathbb{H}\right) := \left\{\gamma : [0,1] \longrightarrow \mathbb{H}, \, \pi(\gamma) \in C_p^2(\mathbb{R}^2), \, \gamma_3(t) = \frac{1}{2} \int_0^t \omega\left(\pi(\gamma)(s), \pi(\gamma)'(s)\right) ds\right\},$$

where $\gamma = (\pi(\gamma), \gamma_3)$, that is, $H_p(\mathbb{H})$ is the set of piecewise continuously twice differentiable horizontal curves. Clearly we have that

$$\overline{H\left(\mathbb{H}\right)} = \overline{H_p\left(\mathbb{H}\right)}$$

We can view g_{δ} as a $H_p(\mathbb{H})$ -valued random variable, that is,

$$g_{\delta}: \Omega \longrightarrow H_p(\mathbb{H}), \quad \omega \mapsto \{t \mapsto g_{\delta}(t,\omega)\},$$

$$(3.6)$$

and hence we can induce a probability measure μ_{δ} on $W_0(\mathbb{H})$ by

$$\mu_{\delta}(E) := \mathbb{P}\left(g_{\delta}^{-1}\left(E \cap H_p\left(\mathbb{H}\right)\right)\right)$$

for any Borel set E in $W_0(\mathbb{H})$.

Proposition 3.2. Let $S_{\mu_{\delta}}$ be the support of the measure μ_{δ} . Then

$$\mathcal{S}_{\mu_{\delta}} \subset \overline{H_p\left(\mathbb{H}\right)} = \overline{H\left(\mathbb{H}\right)}.$$

Proof: By 3.6 we have that $g_{\delta}(\Omega) \subset H_p(\mathbb{H})$ and hence

$$\Omega \subset g_{\delta}^{-1}g_{\delta}\left(\Omega\right) \subset g_{\delta}^{-1}\left(H_{p}\left(\mathbb{H}\right)\right) \subset \Omega.$$

Therefore by the definition of μ_{δ} it follows that

$$1 = \mathbb{P}\left(g_{\delta}^{-1}\left(H_{p}\left(\mathbb{H}\right)\right)\right) = \mu_{\delta}\left(H_{p}\left(\mathbb{H}\right)\right) \leqslant \mu_{\delta}\left(\overline{H_{p}\left(\mathbb{H}\right)}\right) \leqslant 1,$$

and the proof is complete since $S_{\mu_{\delta}}$ is the smallest closed subset of $W_0(\mathbb{H})$ having μ_{δ} -measure one.

We can now state and prove the main result of this section, that is, that the family $\{g_{\delta}\}_{\delta>0}$ is an approximation of the hypoelliptic Brownian motion g in the sense that $\mathbb{E}\left[\max_{0 \leq t \leq 1} \rho(g_{\delta}(t)^{-1}g_t)^2\right] \longrightarrow 0$ as $\delta \to 0$ for any homogeneous norm ρ . As a consequence, the support of the measure μ is contained in $\overline{H(\mathbb{H})}$.

Theorem 3.3. Let $\{g_{\delta}\}_{\delta>0}$ be the sequence defined by (3.5), and ρ be a homogeneous norm on \mathbb{H} . Then

$$\lim_{\delta \to 0} \mathbb{E} \left[\max_{0 \le t \le 1} \rho(g_{\delta}(t)^{-1}g_t)^2 \right] = 0.$$
(3.7)

Proof: By (2.7) and (2.8) we have that

$$\begin{aligned} \max_{0 \leqslant t \leqslant 1} \rho(g_{\delta}(t)^{-1}g_{t})^{4} &\leqslant C \max_{0 \leqslant t \leqslant 1} |g_{\delta}^{-1}(t)g_{t}|^{4} \\ &\leqslant C \max_{0 \leqslant t \leqslant 1} |B_{t} - B_{\delta}(t)|^{4} + C \max_{0 \leqslant t \leqslant 1} \left| A_{t} - A_{\delta}(t) - \frac{1}{2}\omega \left(B_{\delta}(t), B_{t} \right) \right|^{2} \\ &\leqslant C \left(\max_{0 \leqslant t \leqslant 1} |B_{t} - B_{\delta}(t)|^{2} + \max_{0 \leqslant t \leqslant 1} \left| A_{t} - A_{\delta}(t) - \frac{1}{2}\omega \left(B_{\delta}(t), B_{t} \right) \right| \right)^{2}, \end{aligned}$$

and hence

$$\mathbb{E}\left[\max_{0\leqslant t\leqslant 1}\rho(g_{\delta}(t)^{-1}g_{t})^{2}\right]$$

$$\leqslant C\mathbb{E}\left[\max_{0\leqslant t\leqslant 1}|B_{t}-B_{\delta}(t)|^{2}\right]+C\mathbb{E}\left[\max_{0\leqslant t\leqslant 1}|A_{t}-A_{\delta}(t)|\right]+C\mathbb{E}\left[\max_{0\leqslant t\leqslant 1}\left|\frac{1}{2}\omega\left(B_{\delta}(t),B_{t}\right)\right|\right]$$

$$\leqslant C\mathbb{E}\left[\max_{0\leqslant t\leqslant 1}|B_{t}-B_{\delta}(t)|^{2}\right]+C\mathbb{E}\left[\max_{0\leqslant t\leqslant 1}|A_{t}-A_{\delta}(t)|^{2}\right]^{\frac{1}{2}}+C\mathbb{E}\left[\max_{0\leqslant t\leqslant 1}\left|\frac{1}{2}\omega\left(B_{\delta}(t),B_{t}\right)\right|\right],$$

for some constant C (which varies from line to line). By (3.1) and (3.2), we only need to show that

$$\mathbb{E}\left[\max_{0\leqslant t\leqslant 1} \left| \frac{1}{2}\omega\left(B_{\delta}(t), B_{t}\right) \right| \right] \longrightarrow 0, \text{ as } \delta \mapsto 0.$$

Since B_i is independent of $B_{j,\delta} - B_j$ when $i \neq j$, and

$$\begin{split} \max_{0 \le t \le 1} \left| \frac{1}{2} \omega \left(B_{\delta}(t), B_{t} \right) \right| &\leq \frac{1}{2} \max_{0 \le t \le 1} |B_{1}(t)| \max_{0 \le t \le 1} |B_{2,\delta}(t) - B_{2}(t)| \\ &+ \frac{1}{2} \max_{0 \le t \le 1} |B_{2}(t)| \max_{0 \le t \le 1} |B_{1,\delta}(t) - B_{1}(t)|, \end{split}$$

the proof is complete.

Corollary 3.4. We have that $\mu_{\delta} \rightarrow \mu$ weakly. In particular

$$\mathcal{S}_{\mu} \subset \overline{H\left(\mathbb{H}\right)}.\tag{3.8}$$

Proof: Let us first show that $\{g_{\delta}\}_{\delta>0}$ converges to g in probability in $W_0(\mathbb{H})$. For any fixed $\varepsilon > 0$ we have that

$$\mathbb{P}\left(\max_{0\leqslant t\leqslant 1}\rho(g_{\delta}(t)^{-1}g_{t})>\varepsilon\right)\leqslant \frac{1}{\varepsilon^{2}}\mathbb{E}\left[\max_{0\leqslant t\leqslant 1}\rho(g_{\delta}(t)^{-1}g_{t})^{2}\right]$$

which goes to zero by Theorem 3.3. Therefore $\{g_{\delta}\}_{\delta>0}$ converges to g in distribution, and hence μ_{δ} converges weakly to μ in $W_0(\mathbb{H})$. Thus, for any closed set F in $W_0(\mathbb{H})$ we have that

$$\mu(F) \geqslant \limsup_{\delta \to 0} \mu_{\delta}(F).$$

In particular, for $F = \overline{H_p(\mathbb{H})}$ and by Proposition 3.2 it follows that

$$\mu\left(\overline{H_{p}\left(\mathbb{H}\right)}\right) \geqslant \limsup_{\delta \to 0} \mu_{\delta}(\overline{H_{p}\left(\mathbb{H}\right)}) = 1$$

Since S_{μ} is the smallest closed subset having μ -measure one, we have that $S_{\mu} \subset \overline{H_p(\mathbb{H})} = \overline{H(\mathbb{H})}$.

We conclude this section showing that for each fixed δ , the measures μ and μ_{δ} are singular.

Proposition 3.5. For each δ the measures μ and μ_{δ} are singular.

Proof: From the proof of Proposition 3.2 we know that $\mu_{\delta}(H(\mathbb{H})) = 1$. It is then enough to show that $\mu(H(\mathbb{H})) = 0$. Let us denote by ν the law of a two-dimensional standard Brownian motion. By definition of g_t and π (2.10), the following diagram commutes

$$\Omega \xrightarrow{g} W_0(\mathbb{H})$$

$$\downarrow^{\pi}$$

$$W_0(\mathbb{R}^2),$$

and for any Borel set E in $W_0(\mathbb{R}^2)$ we have that

$$\nu(E) := \mathbb{P}(B^{-1}(E)) = \mathbb{P}(g^{-1} \circ \pi^{-1}(E)) = \mu(\pi^{-1}(E))$$

Moreover, from Remark 2.13 we know that $\pi(H(\mathbb{H}))$ is the Cameron-Martin space over \mathbb{R}^2 , which is known to have ν -measure zero, see Gross (1967) for more details. Therefore we can conclude that

$$\mu\left(H\left(\mathbb{H}\right)\right) \leqslant \mu\left(\pi^{-1}\pi\left(H\left(\mathbb{H}\right)\right)\right) = \nu\left(\pi\left(H\left(\mathbb{H}\right)\right)\right) = 0.$$

3.2. Support of the Wiener measure. The goal of this section is to prove that $\overline{H(\mathbb{H})} \subset S_{\mu}$ which will follow from Theorem 3.6. Moreover, in Proposition 3.10 we show that $\overline{H(\mathbb{H})} = W_0(\mathbb{H})$.

Theorem 3.6. Let $\phi = (\pi(\phi), \phi_3) = \in H_p(\mathbb{H})$. For $\delta > 0$ let us denote by $E_{\delta,\phi}$ the event

$$E_{\delta,\phi} := \left\{ \max_{0 \leqslant t \leqslant 1} |B_t - \pi(\phi)(t)|_{\mathbb{R}^2} < \delta \right\}.$$

Then for any $\varepsilon > 0$

$$\lim_{\delta \to 0} \mathbb{P}\left(\max_{\substack{0 \le t \le 1\\ \forall \pi}} \rho(\phi(t)^{-1}g_t) > \varepsilon \,|\, E_{\delta,\phi}\right) = 0,$$

where ρ is a homogeneous norm on \mathbb{H} .

Proof: By (2.7) it is enough to prove it for the homogeneous norm given by (2.8). For $\phi \in H_p(\mathbb{H})$ we have that

$$\begin{aligned} \max_{0 \leqslant t \leqslant 1} |\phi(t)^{-1}g_t|^4 &\leqslant \max_{0 \leqslant t \leqslant 1} |B_t - \pi(\phi)(t)|_{\mathbb{R}^2}^4 \\ &+ \max_{0 \leqslant t \leqslant 1} \left| \frac{1}{2} \int_0^t \omega \left(B_s - \pi(\phi)(s), dB_s - \pi(\phi)'(s) ds \right) + \int_0^t \omega \left(B_s - \pi(\phi)(s), \pi(\phi)'(s) \right) ds \right|^2 \\ &\leqslant \left(\max_{0 \leqslant t \leqslant 1} |B_t - \pi(\phi)(t)|_{\mathbb{R}^2}^2 \right) \\ &+ \max_{0 \leqslant t \leqslant 1} \left| \frac{1}{2} \int_0^t \omega \left(B_s - \pi(\phi)(s), dB_s - \pi(\phi)'(s) ds \right) + \int_0^t \omega \left(B_s - \pi(\phi)(s), \pi(\phi)'(s) \right) ds \right| \right)^2. \end{aligned}$$

Therefore on the event $E_{\delta,\phi}$ we have that

$$\begin{split} &\max_{0\leqslant t\leqslant 1} |\phi(t)^{-1}g_t|^2 \leqslant \max_{0\leqslant t\leqslant 1} |B_t - \pi(\phi)(t)|_{\mathbb{R}^2}^2 \\ &+ \max_{0\leqslant t\leqslant 1} \left| \frac{1}{2} \int_0^t \omega \left(B_s - \pi(\phi)(s), dB_s - \pi(\phi)'(s) ds \right) + \int_0^t \omega \left(B_s - \pi(\phi)(s), \pi(\phi)'(s) \right) ds \right| \\ &\leqslant \delta^2 + \max_{0\leqslant t\leqslant 1} \left| \int_0^t \omega \left(B_s - \pi(\phi)(s), \pi(\phi)'(s) \right) ds \right| + \max_{0\leqslant t\leqslant 1} \left| \frac{1}{2} \int_0^t \omega \left(B_s - \pi(\phi)(s), dB_s - \pi(\phi)'(s) ds \right) \right| \\ &\leqslant \delta^2 + \delta C_\phi + \max_{0\leqslant t\leqslant 1} \left| \frac{1}{2} \int_0^t \omega \left(B_s - \pi(\phi)(s), dB_s - \pi(\phi)'(s) ds \right) \right|, \end{split}$$

where $C_{\phi} := \int_0^1 |\phi_1'(s)| + |\phi_2'(s)| ds$. It then follows that

$$\mathbb{P}\left(\max_{0\leqslant t\leqslant 1} |\phi(t)^{-1}g_t| > \varepsilon |E_{\delta,\phi}\right)$$

$$\leqslant \mathbb{P}\left(\max_{0\leqslant t\leqslant 1} \left|\frac{1}{2} \int_0^t \omega \left(B_s - \pi(\phi)(s), dB_s - \pi(\phi)'(s)ds\right)\right| > \varepsilon^2 - \delta C_\phi - \delta^2 |E_{\delta,\phi}\right).$$

Note that this last expression only depends on the process $B_t - \pi(\phi)(t)$. Since $\phi = (\pi(\phi), \phi_3) \in H_p(\mathbb{H})$, by Remark 2.13 we know that $\pi(\phi)$ belongs to the Cameron-Martin space over \mathbb{R}^2 . Therefore from the Cameron-Martin-Girsanov Theorem there exists a probability measure \mathbb{Q}^{ϕ} such that the process $B_t^{\phi} := B_t + \pi(\phi)(t)$ is a Brownian motion under \mathbb{Q}^{ϕ} . More precisely there exists an exponential martingale \mathcal{E}^{ϕ} such that

$$\mathbb{Q}^{\phi}(A) = \mathbb{E}\left[\mathcal{E}^{\phi}\mathbf{1}_{A}\right] \quad \forall A \in \mathcal{F},$$

where $\mathcal{E}^{\phi} = \exp\left(-\int_{0}^{1} \langle \pi(\phi)'(s), dB_{s} \rangle_{\mathbb{R}^{2}} ds - \frac{1}{2}\int_{0}^{1} |\pi(\phi)'(s)|_{\mathbb{R}^{2}}^{2} ds\right)$. Note that
 $d\left(B_{t} - \pi(\phi)(t)\right) = dB_{t} - \pi(\phi)'(t)dt$, and
 $dB_{t} = dB_{t}^{\phi} - \pi(\phi)'(t)dt$,

that is, the law of $B_t - \pi(\phi)(t)$ under \mathbb{P} is the same as the law of B_t under \mathbb{Q}^{ϕ} . Therefore we can write

$$\mathbb{P}(E_{\delta,\phi}) = \mathbb{P}\left(\max_{0 \le t \le 1} |B_t - \pi(\phi)(t)|_{\mathbb{R}^2} < \delta\right)$$
$$= \mathbb{Q}^{\phi}\left(\max_{0 \le t \le 1} |B_t|_{\mathbb{R}^2} < \delta\right) = \mathbb{E}\left[\mathcal{E}^{\phi} \mathbf{1}_{E_{\delta}}\right] = \mathbb{E}\left[\mathcal{E}^{\phi} |E_{\delta}\right] \mathbb{P}(E_{\delta})$$

where we set $E_{\delta} := \{ \max_{0 \leq t \leq 1} |B_t|_{\mathbb{R}^2} < \delta \}$. Similarly we have that

$$\mathbb{P}\left(\max_{0\leqslant t\leqslant 1} \left| \frac{1}{2} \int_0^t \omega \left(B_s - \pi(\phi)(s), dB_s - \pi(\phi)'(s) ds \right) \right| > \varepsilon^2 - \delta C_\phi - \delta^2, \ E_{\delta,\phi} \right)$$
$$= \mathbb{E}\left[\mathcal{E}^{\phi} | F_{\delta,\phi}^{\varepsilon} \cap E_{\delta} \right] \mathbb{P}\left(F_{\delta,\phi}^{\varepsilon} \cap E_{\delta} \right),$$

where $F_{\delta,\phi}^{\varepsilon} := \left\{ \max_{0 \le t \le 1} \left| \frac{1}{2} \int_0^t \omega \left(B_s, dB_s \right) \right| > \varepsilon^2 - \delta C_{\phi} - \delta^2 \right\}$. Therefore it follows that

$$\mathbb{P}\left(\max_{0\leqslant t\leqslant 1} |\phi(t)^{-1}g_{t}| > \varepsilon |E_{\delta,\phi}\right) \\
\leqslant \mathbb{P}\left(\max_{0\leqslant t\leqslant 1} \left|\frac{1}{2}\int_{0}^{t} \omega \left(B_{s} - \pi(\phi)(s), dB_{s} - \pi(\phi)'(s)ds\right)\right| > \varepsilon^{2} - \delta C_{\phi} - \delta^{2} |E_{\delta,\phi}\right) \\
= \frac{\mathbb{P}\left(F_{\delta,\phi}^{\varepsilon} \cap E_{\delta}\right) E\left[\mathcal{E}^{\phi}|F_{\delta,\phi}^{\varepsilon} \cap E_{\delta}\right]}{\mathbb{P}\left(E_{\delta}\right) \mathbb{E}\left[\mathcal{E}^{\phi}|E_{\delta}\right]} = \mathbb{P}\left(F_{\delta,\phi}^{\varepsilon}|E_{\delta}\right) \times \frac{E\left[\mathcal{E}^{\phi}|F_{\delta,\phi}^{\varepsilon} \cap E_{\delta}\right]}{\mathbb{E}\left[\mathcal{E}^{\phi}|E_{\delta}\right]} \tag{3.9}$$

We will show later in the paper, see Lemma 3.9, that for any $\varepsilon > 0$ and any $\phi \in H_p(\mathbb{H})$ we have that

$$\lim_{\delta \to 0} \frac{\mathbb{E}\left[\mathcal{E}^{\phi} \mid F_{\delta,\phi}^{\varepsilon} \cap E_{\delta}\right]}{\mathbb{E}\left[\mathcal{E}^{\phi} \mid E_{\delta}\right]} = 1.$$
(3.10)

In light of 3.9 and 3.10, the proof will be completed once we show that

$$\begin{split} &\lim_{\delta \to 0} \mathbb{P}\left(F_{\delta,\phi}^{\varepsilon} \mid E_{\delta}\right) := \\ &\lim_{\delta \to 0} \mathbb{P}\left(\max_{0 \leqslant t \leqslant 1} \left|\frac{1}{2} \int_{0}^{t} \omega\left(B_{s}, dB_{s}\right)\right| > \varepsilon^{2} - \delta C_{\phi} - \delta^{2} \left|\max_{0 \leqslant t \leqslant 1} |B_{t}|_{\mathbb{R}^{2}} < \delta\right) = 0. \end{split}$$

The process $A_t := \frac{1}{2} \int_0^t \omega(B_s, dB_s)$ is a square integrable martingale with zero mean, and therefore there exists a one dimensional Brownian motion b_t such that

$$b_{\tau(t)} = \frac{1}{2} \int_0^t \omega\left(B_s, dB_s\right) \,,$$

where $\tau(t) = \frac{1}{4} \int_0^t B_1(s)^2 + B_2(s)^2 ds$. Moreover it is known that b_t is independent of B_t Ikeda and Watanabe (1989, Chapter 6 p. 470). Hence we have that

$$\begin{split} & \mathbb{P}\left(\max_{0\leqslant t\leqslant 1}|b_{\tau(t)}|>\varepsilon^{2}-\delta C_{\phi}-\delta^{2}|,\ \max_{0\leqslant t\leqslant 1}|B_{t}|_{\mathbb{R}^{2}}<\delta\right)\\ & \leqslant \mathbb{P}\left(\max_{0\leqslant t\leqslant \frac{1}{4}\delta^{2}}|b_{t}|>\varepsilon^{2}-\delta C_{\phi}-\delta^{2}|,\ \max_{0\leqslant t\leqslant 1}|B_{t}|_{\mathbb{R}^{2}}<\delta\right)\\ & = \mathbb{P}\left(\max_{0\leqslant t\leqslant \frac{1}{4}\delta^{2}}|b_{t}|>\varepsilon^{2}-\delta C_{\phi}-\delta^{2}\right)=\mathbb{P}\left(\max_{0\leqslant t\leqslant 1}|b_{t}|>2(\frac{\varepsilon^{2}}{\delta}-C_{\phi}-\delta)\right), \end{split}$$

which goes to zero as δ goes to zero.

Corollary 3.7. $\overline{H(\mathbb{H})} \subset \mathcal{S}_{\mu}$.

Proof: Let us first prove that for any $\phi \in H_p(\mathbb{H})$ and $\varepsilon > 0$ we have that $\mu(B_{\varepsilon}(\phi)) > 0$, where

$$B_{\varepsilon}(\phi) := \left\{ \gamma \in W_0\left(\mathbb{H}\right), \ \max_{0 \leqslant t \leqslant 1} |\phi^{-1}(t)\gamma(t)| < \varepsilon \right\}$$

Indeed, for any $\phi \in H_p(\mathbb{H})$ and $\varepsilon > 0$ we have that

$$\mu(B_{\varepsilon}(\phi)) := \mathbb{P}\left(g \in B_{\varepsilon}(\phi)\right) = \mathbb{P}\left(\max_{0 \le t \le 1} |\phi(t)^{-1}g_t| < \varepsilon\right)$$

$$\geqslant \mathbb{P}\left(\max_{0 \le t \le 1} |\phi(t)^{-1}g_t| < \varepsilon \,|\, E_{\delta,\phi}\right) \mathbb{P}\left(E_{\delta,\phi}\right),$$

where $E_{\delta,\phi} := \{\max_{0 \leq t \leq 1} |B_t - \pi(\phi)(t)|_{\mathbb{R}^2} < \delta\}$. From Theorem (3.6) there exists a δ_0 such that for every $\delta \in (0, \delta_0)$

$$\mathbb{P}\left(\max_{0 \le t \le 1} |\phi(t)^{-1}g_t| < \varepsilon \,|\, E_{\delta,\phi}\right) \ge \frac{1}{2},$$

for any $\varepsilon > 0$. Combining everything together we have that

$$\mu\left(B_{\varepsilon}(\phi)\right) \geqslant \frac{1}{2} \mathbb{P}\left(\sup_{0 \leqslant t \leqslant 1} |B_t - \pi(\phi)(t)|_{\mathbb{R}^2} < \frac{\delta_0}{2}\right),$$

and the latter is positive since $\pi(\phi)$ is in the Cameron-Martin space over \mathbb{R}^2 . Therefore, if O is any open set in $W_0(\mathbb{H})$ with $\mu(O) = 0$ then $O \subset H_p(\mathbb{H})^c$, and hence

$$\bigcup_{\substack{O \text{ open}\\ \mu(O)=0}} O \subset H_p\left(\mathbb{H}\right)^c, \text{ that is, } \mathcal{S}_{\mu} := \bigcap_{\substack{F \text{ closed}\\ \mu(F)=1}} F \supset H_p\left(\mathbb{H}\right),$$

and since S_{μ} is closed, we have that $\mathcal{S}_{\mu} \supset \overline{H_p(\mathbb{H})} = \overline{H(\mathbb{H})}$.

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The proof of Theorem 3.6 will be completed once we show (3.10). Before proceeding to the proof of (3.10), we need the following lemma whose proof can be found in Ikeda and Watanabe (1989, pp. 536-537).

Lemma 3.8 (pp. 536-537 in Ikeda and Watanabe (1989)). Let I_1, \ldots, I_n be n random variables on a probability space $(\Omega, \mathcal{F}, \mathbb{P})$. Let $\{A_{\delta}\}_{0 < \delta < 1}$ be a family of events in \mathcal{F} and a_1, \ldots, a_n be n numbers. If for every real number c and every $1 \leq i \leq n$

$$\limsup_{\delta \to 0} \mathbb{E} \left[\exp(c I_i) | A_{\delta} \right] \leqslant \exp(c a_i),$$

then

$$\lim_{\delta \to 0} \mathbb{E}\left[\exp\left(\sum_{i=1}^{n} I_i\right) | A_{\delta}\right] = \exp\left(\sum_{i=1}^{n} a_i\right).$$

Lemma 3.9. Let E_{δ} and $F_{\delta,\phi}^{\varepsilon}$ be given as in the proof of Theorem 3.6. Then

$$\lim_{\delta \to 0} \frac{\mathbb{E}\left[\mathcal{E}^{\phi} \mid F_{\delta,\phi}^{\varepsilon} \cap E_{\delta}\right]}{\mathbb{E}\left[\mathcal{E}^{\phi} \mid E_{\delta}\right]} = 1.$$

Proof: Let us first prove that

$$\lim_{\delta \to 0} \mathbb{E}\left[\mathcal{E}^{\phi} \mid E_{\delta}\right] = \exp\left(-\frac{1}{2}\int_{0}^{1} |\pi(\phi)'(s)|_{\mathbb{R}^{2}}^{2} ds\right).$$
(3.11)

Since $\mathcal{E}^{\phi} = \exp\left(-\int_{0}^{1} \langle \pi(\phi)'(s), dB_s \rangle_{\mathbb{R}^2} ds - \frac{1}{2} \int_{0}^{1} |\pi(\phi)'(s)|_{\mathbb{R}^2}^2 ds\right)$, by Lemma 3.8 and the definition of E_{δ} , it is enough to show that for any real number c and i = 1, 2

$$\limsup_{\delta \to 0} \mathbb{E}\left[\exp\left(-c \int_0^1 \phi_i'(s) dB_i(s)\right) \, \Big| \, \max_{0 \leqslant t \leqslant 1} |B_t|_{\mathbb{R}^2} < \delta \right] \leqslant 1.$$

For $\phi \in H_p(\mathbb{H})$ we can write $\int_0^1 \phi'_i(s) dB_i(s) = \phi'_i(1)B_i(1) - \int_0^1 \phi''_i(s)B_i(s) ds$, and hence on the event E_{δ} we have that

$$\exp\left(-c\int_0^1\phi_i'(s)dB_i(s)\right)\leqslant\exp\left(-ck_{\phi}\delta\right),$$

for some finite constant k_{ϕ} only depending on ϕ . Therefore we have that

$$\limsup_{\delta \to 0} \mathbb{E}\left[\exp\left(-c\int_0^1 \phi_i'(s)dB_i(s)\right) \ \bigg| \max_{0 \leqslant t \leqslant 1} |B_t|_{\mathbb{R}^2} < \delta\right] \leqslant \mathbb{E}\left[\limsup_{\delta \to 0} e^{-ck_{\phi}\delta}|E_{\delta}\right] \leqslant 1.$$

In a similar way it can be shown that

$$\lim_{\delta \to 0} \mathbb{E}\left[\mathcal{E}^{\phi} \,|\, F^{\varepsilon}_{\delta,\phi} \cap E_{\delta}\right] = \exp\left(-\frac{1}{2} \int_{0}^{1} |\pi(\phi)'(s)|^{2}_{\mathbb{R}^{2}} ds\right),$$

and the proof is completed.

The following Proposition completes the proof of Theorem 2.14.

Proposition 3.10. We have that $\overline{H(\mathbb{H})} = W_0(\mathbb{H})$.

Proof: Any element in $W_0(\mathbb{H})$ can be approximated with piecewise linear curves in the uniform topology. It is then enough to prove that for any piecewise linear curve ξ there exists a sequence of horizontal finite energy curves $\{\phi_n^{\xi}\}_{n\in\mathbb{N}}$ such that $d_{W_0(\mathbb{H})}(\phi_n^{\xi},\xi) \to 0$. Let us first explain the geometric construction through the following example. Consider the curve $t \to \xi(t) = (0,0,t) \in \mathbb{H}$

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for $t \in [0, 1]$, which is the prototype of a non-horizontal curve. Let us define a family of finite energy horizontal curves ϕ_n by

$$\phi_n(s) := \left(\frac{2}{n}\cos\left(n^2 s\right), \frac{1}{n}\sin\left(n^2 s\right), s\right).$$

Geometrically, the curves ϕ_n are helics that shrink around the ξ as n goes to infinity. Indeed,

$$d_{W_0(\mathbb{H})} (\phi_n, \xi)^4 = \max_{0 \le s \le 1} \left[\left(\left(\frac{2}{n} \cos\left(n^2 s\right) \right)^2 + \left(\frac{1}{n} \sin\left(n^2 s\right) \right)^2 \right) + \left(s - \frac{1}{2} \int_0^t \omega \left(\pi(\phi)_n(u), \pi(\phi)'_n(u) \right) du \right)^2 \right]$$
$$= \max_{0 \le s \le 1} \left(\left(\frac{2}{n} \cos\left(n^2 s\right) \right)^2 + \left(\frac{1}{n} \sin\left(n^2 s\right) \right)^2 \right)^2 \longrightarrow 0,$$

as n goes to infinity. Now, let $\xi(t) = (a_1t, a_2t, a_3t)$ be a linear curve in \mathbb{H} , where $a_1, a_2, a_3 \in \mathbb{R}$. Then set

$$\phi_n(s) := \left(a_1 s + \frac{2}{n}\cos\left(n^2 a_3 s\right), a_2 s + \frac{1}{n}\sin\left(n^2 a_3 s\right), a_3 s - \frac{a_2 s}{n}\cos\left(n^2 a_3 s\right) + \frac{a_1 s}{2n}\sin\left(n^2 a_3 s\right) + \frac{1}{n}\int_0^s 2a_2\cos\left(n^2 a_3 u\right) - a_1\sin\left(n^2 a_3 u\right) du\right).$$

It is easy to check that for any $n \in \mathbb{N}$, ϕ_n is a finite energy horizontal curve such that

$$(\phi_n^{-1}\xi)(s) = \left(-\frac{2}{n}\cos\left(n^2a_3s\right), -\frac{1}{n}\sin\left(n^2a_3s\right), \frac{1}{n}\int_0^s a_1\sin\left(n^2a_3u\right) - 2a_2\cos\left(n^2a_3u\right)du\right),$$

which implies that $d_{W_0(\mathbb{H})}(\phi_n,\xi) \to 0$ as n goes to infinity.

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