ANALYSIS OF WIENER MEASURE ON PATH AND LOOP GROUPS

BRUCE K. DRIVER

Contents

1. Introduction	1
2. Path Group Results	6
3. Pinned Wiener Measure Results	11
4. Comparing Heat Kernel and Pinned Wiener Mea	asures 17
Appendix A. A Wong - Zakai Approximation Theo	orem 19
References	27

1. INTRODUCTION

This article gives a unified explanation of many of the results regarding analysis of Wiener and pinned Wiener measure on path and loop groups which have been developed over the last decade. Most of the results discussed below are outgrowths or have been inspired by Leonard Gross' pioneering papers [36, 37]. Although this article is mostly a survey, the result in Theorem 2.5 below is new.

The main theme of this article is to explain how many properties of Wiener and pinned Wiener measure on path and loop groups can be traced back to the fact that the path group of interest is flat in the sense of Riemannian geometry, see Corollary 1.11 below. This point of view is often best understood using certain heuristic type path integral arguments for Wiener measure. As such, it is useful to give heuristic theorems (labelled as Meta-Theorems below) and heuristic proofs (labelled as Meta-Proofs below) of many of the results. Nevertheless we will always give a corresponding honest theorem capturing the "content" of the intuitive Meta-Theorem. Similarly, along with all Meta-Proofs of honest theorems, we will give references to the literature containing a rigorous proof.

1.1. Notation.

Notation 1.1. Let K be a connected compact Lie group, $\mathfrak{k} := T_e K$ be the Lie algebra of $K, T \langle \cdot, \cdot \rangle_{\mathfrak{k}}$ be an Ad_K -invariant inner product on \mathfrak{k} and let $\langle \cdot, \cdot \rangle_{\mathfrak{k}}$ denote the unique bi-invariant Riemannian metric on K which agrees with $\langle \cdot, \cdot \rangle_{\mathfrak{k}}$ on $\mathfrak{k} := T_e K$. To simplify notation later we will assume that K is a matrix group in which case \mathfrak{k} may also be viewed as a matrix Lie algebra. (Since K is compact, this is no restriction, see for example Theorem 4.1 on p. 136 in [11].) Elements $A \in \mathfrak{k}$ will be

Date: February 18, 2002 File:gross.tex.

This research was partially supported by NSF Grants DMS 96-12651 and DMS 99-71036.

identified with the unique left invariant vector field on K agreeing with A at the identity in K, i.e. if $f \in C^{\infty}(K)$ then

$$Af(x) = \frac{d}{dt}|_0 f(xe^{tA}).$$

Example 1.2. As an example, let K = SO(3) be the group of 3×3 real orthogonal matrices with determinant 1. The Lie algebra of K is $\mathfrak{k} = so(3)$, the set of 3×3 real skew symmetric matrices, and the inner product $\langle A, B \rangle_{\mathfrak{k}} := -\mathrm{tr}(AB)$ is an example of an Ad_K – invariant inner product on \mathfrak{k} .

Our main interest here is the path and loop spaces built on K which we now define.

Notation 1.3. Suppose that $M = \mathfrak{k}$ or K and $o = 0 \in \mathfrak{k}$ or $o = e \in K$ respectively. Let W(M) denote the collection of continuous paths $\sigma : [0,1] \to M$ such that $\sigma(0) = o \in M$. The subset of finite energy paths H(M) consists of those $\sigma \in W(M)$ which are absolutely continuous and satisfy $E_M(\sigma) < \infty$ where

(1.1)
$$E_{\mathfrak{k}}(\sigma) := \int_0^1 |\sigma'(s)|_{\mathfrak{k}}^2 ds < \infty \text{ and } E_K(\sigma) := \int_0^1 \left| [\sigma(s)]^{-1} \sigma'(s) \right|_{\mathfrak{k}}^2 ds.$$

Also let $\mathcal{L}(M)$ and $H_0(M)$ denote the space of loops inside W(M) and H(M) respectively. In particular,

(1.2)
$$W(K) \equiv \{ \sigma \in C \left([0,1] \to K \right) | \sigma \left(0 \right) = e \}$$

and

(1.3)
$$\mathcal{L}(K) \equiv \{ \sigma \in W(K) | \sigma(1) = e \}$$

respectively. Also let $\mathbf{e} \in \mathcal{L}(K)$ denote the constant path at identity $e \in K$, i.e. $\mathbf{e}(s) = e$ for $s \in [0, 1]$.

Because $\langle \cdot, \cdot \rangle_{\mathfrak{k}}$ is Ad_K – invariant, the formula for the **energy**, $E_K(\sigma)$, in Eq. (1.1) may also be expressed as

(1.4)
$$E_K(\sigma) = \int_0^1 \left| \sigma'(s)\sigma(s)^{-1} \right|_{\mathfrak{k}}^2 ds.$$

As usual we will refer to $H(\mathfrak{k})$ equipped with the Hilbertian inner product,

(1.5)
$$(h,k) := \int_0^1 \langle h'(s), k'(s) \rangle \ ds$$

as the Cameron – Martin Hilbert space. Notice that $E_{\mathfrak{k}}(h) = (h, h)$.

Remark 1.4. It is well known that H(K) is a Hilbert Lie group under pointwise multiplication and that the map

$$(x,h) \in H(K) \times H(\mathfrak{k}) \to L_{x*}h \in T(H(K))$$

is a trivialization of the tangent bundle of H(K). (We are using $L_x : H(K) \to H(K)$ to denote left multiplication by x.) This trivialization induces a left-invariant Riemannian metric (\cdot, \cdot) on H(K) given explicitly by

(1.6)
$$(L_{x*}h, L_{x*}h) = \int_0^1 \langle h'(s), h'(s) \rangle \ ds \quad \forall x \in H(K) \text{ and } h \in H(\mathfrak{k}).$$

See Appendix A in [20] and the references therein for more details.

Notation 1.5. In the sequel we will fix a special type (for technical convergence issues later) of orthonormal bases, S and S_0 , for $H(\mathfrak{k})$ and $H_0(\mathfrak{k})$ respectively. Namely we assume,

$$S = \{\ell \xi : \ell \in S_{\mathbb{R}} \text{ and } \xi \in \beta\} \text{ and } S_0 = \{\ell \xi : \ell \in (S_{\mathbb{R}})_0 \text{ and } \xi \in \beta\}$$

where β is an orthonormal basis for \mathfrak{k} and $S_{\mathbb{R}}$ and $(S_{\mathbb{R}})_0$ are orthonormal bases for $H(\mathbb{R})$ and $H_0(\mathbb{R})$ respectively.

Definition 1.6 (Differential Operators). For $h \in H(\mathfrak{k})$ $(h \in H_0(\mathfrak{k}))$, let \tilde{h} denote the left invariant vector field on H(K) $(H_0(K))$ such that $\tilde{h}(e) = h$, i.e. if $f \in C^1(H(K))$ $(f \in C^1(H_0(K)))$ and $x \in H(K)$ $(x \in H_0(K))$ then

$$\tilde{h}f(x) = \left. \frac{d}{dt} \right|_0 f\left(x e^{th} \right)$$

where $(xe^{th})(s) = x(s)e^{th(s)}$ for all $s \in [0,1]$. Also for $f \in C^2(H(K))$, let

$$\|\operatorname{grad} f\|^2 := \sum_{h \in S} \left(\tilde{h}f \right)^2$$
 and $\triangle_{H(K)}f := \sum_{h \in S} \tilde{h}^2 f$

and for $f \in C^2(H_0(K))$ let

$$\|\operatorname{grad}_0 f\|^2 := \sum_{h \in S_0} \left(\tilde{h}f \right)^2$$
 and $\triangle_{H_0(K)} f := \sum_{h \in S_0} \tilde{h}^2 f.$

Similarly for $f \in C^2(H(\mathfrak{k}))$ and $h \in H(\mathfrak{k})$ let

$$\partial_h f(x) := \frac{d}{dt} |_0 f(x+th), \ \|\operatorname{grad} f\|^2 := \sum_{h \in S} (\partial_h f)^2 \ \text{and} \ \triangle_{H(\mathfrak{k})} f := \sum_{h \in S} \partial_h^2 f$$

and for $f \in C^2(H_0(\mathfrak{k}))$ and $h \in H_0(\mathfrak{k})$ let

$$\partial_h f(x) := \frac{d}{dt} |_0 f(x+th), \quad \|\operatorname{grad}_0 f\|^2 := \sum_{h \in S_0} (\partial_h f)^2 \text{ and } \triangle_{H_0(\mathfrak{k})} f := \sum_{h \in S_0} \partial_h^2 f.$$

Notation 1.7. If μ is a probability measure on a measure space (Ω, \mathcal{F}) and $f \in L^1(\mu) = L^1(\Omega, \mathcal{F}, \mu)$, we will often write $\mu(f)$ and sometimes $\mathbb{E}f$ for the integral, $\int_{\Omega} f d\mu$.

1.2. Basis Geometrical Results.

Meta-Theorem 1.8. The Lie groups H(K) and $H_0(K)$ should be thought of as being unimodular, i.e. the "Riemannian volume form" is formally invariant under both left and right translations. As a result, $\Delta_{H(K)}$ and $\Delta_{H_0(K)}$ should be interpreted as the Laplace Beltrami operators on H(K) and $H_0(K)$ respectively.

Proof. (Meta Proof) One way to "verify" this assertion would be to show that $\operatorname{trace}(ad_h) = 0$ for all $h \in H(\mathfrak{k})$ $(h \in H_0(\mathfrak{k}))$. The problem here however is that $ad_h : H(\mathfrak{k}) \to H(\mathfrak{k})$ is not a trace class operator in general. However, if we compute the "trace" using the type of basis described in Notation 1.5 we do get

trace $(ad_h) = 0$:

$$\operatorname{trace}\left(ad_{h}\right) = \sum_{k \in S} \langle [h, k], k \rangle_{H(\mathfrak{k})} = \sum_{k \in S} \int_{0}^{1} \langle [h, k]'(s), k'(s) \rangle_{\mathfrak{k}} ds$$
$$= \sum_{k \in S} \int_{0}^{1} \langle [h', k](s), k'(s) \rangle_{\mathfrak{k}} ds = \sum_{k \in S} \int_{0}^{1} \langle [h'(s), [k(s), k'(s)] \rangle_{\mathfrak{k}} ds = 0$$

where in the fourth equality we have used the Ad_K – invariance of $\langle \cdot, \cdot \rangle$ and in the last equality the fact that [k(s), k'(s)] = 0 for all $s \in [0, 1]$ and $k \in S$. So even though ad_h is not trace class (i.e. the trace is basis dependent) for the type of bases that we consider one should interpret trace (ad_h) to be 0. Similar computations hold in the $H_0(K)$ case as well.

Definition 1.9 (Left and Right Anti-development Motions). The left and right anti-development maps, $b: H(K) \to H(\mathfrak{k})$ and $B: H(K) \to H(\mathfrak{k})$ respectively are defined by

(1.7)
$$b_s(x) := \int_0^s x(s)^{-1} x'(s) ds \text{ and}$$

(1.8)
$$B_s(x) := \int_0^s x'(s)x(s)^{-1}ds.$$

The following theorem appears (in a disguised form) in Theorem 3.14 and Lemma 3.15 of Gross [36], also see Shigekawa [56], Fang and Franchi [30], Driver and Hall [20].

Proposition 1.10. Let $b, B : H(K) \to H(\mathfrak{k})$ be as in Definition 1.9, then

(1.9)
$$\left(B_*\tilde{h}(x)\right)'_s = \tilde{h}B'_s(x) = Ad_{x(s)}h'(s)$$

and

(1.10)
$$\tilde{h}b'_s(x) = h'(s) + [b'_s(x), h(s)].$$

Proof. These are routine computations:

$$\begin{split} \tilde{h}B'_s(x) &= \frac{d}{dt}|_0 \left\{ (xe^{th})'(s)(x(s)e^{th(s)})^{-1} \right\} \\ &= \frac{d}{dt}|_0 \left\{ \left[x'(s)e^{th(s)} + x(s)(e^{th(s)})' \right] e^{-th(s)}x(s)^{-1} \right\} \\ &= \frac{d}{dt}|_0 \left\{ x'(s)x(s)^{-1} + x(s)(e^{th(s)})'e^{-th(s)}x(s)^{-1} \right\} = Ad_{x(s)}h'(s) \end{split}$$

and

$$\begin{split} \tilde{h}b'_s(x) &= \frac{d}{dt}|_0 \left\{ (x(s)e^{th(s)})^{-1}(xe^{th})'(s) \right\} \\ &= \frac{d}{dt}|_0 \left\{ e^{-th(s)}x(s)^{-1} \left[x'(s)e^{th(s)} + x(s)(e^{th(s)})' \right] \right\} \\ &= \frac{d}{dt}|_0 \left\{ e^{-th(s)}b'_s(x)e^{th(s)} + e^{-th(s)}(e^{th(s)})' \right\} \\ &= h'(s) + [b'_s(x), h(s)]. \end{split}$$

The next corollary is the most crucial result for the rest of this paper.

Corollary 1.11. The map $B : H(K) \to H(\mathfrak{k})$ is an isometric isomorphism of infinite-dimensional Riemannian manifolds. In particular, H(K) is flat.

Proof. By Eq. (1.9),

$$\langle B_*\tilde{h}(x), B_*\tilde{h}(x)\rangle_{H(\mathfrak{k})} = \int_0^1 \left| \left(B_*\tilde{h}(x) \right)_s' \right|^2 ds = \int_0^1 \left| Ad_{x(s)}h'(s) \right|^2 ds$$
$$= \int_0^1 \left| h'(s) \right|^2 ds = \langle h, h \rangle_{H(\mathfrak{k})} = \langle \tilde{h}, \tilde{h} \rangle_{T_xH(K)}.$$

Theorem 1.12. The operators, $\Delta_{H(\mathfrak{k})}$ on $H(\mathfrak{k})$ and $\Delta_{H(K)}$ on H(K), are intertwined by B. More precisely if $f : H(K) \to \mathbb{C}$ is a smooth cylinder function, then $\Delta_{H(\mathfrak{k})}(f \circ B^{-1})$ and $\Delta_{H(K)}f$ exist and

(1.11)
$$\Delta_{H(\mathfrak{k})}(f \circ B^{-1}) = (\Delta_{H(K)}f) \circ B^{-1}$$

Proof. (Meta-Proof.) Since *B* is an isometry and, by Meta-Theorem 1.8, $\Delta_{H(\mathfrak{k})}$ and $\Delta_{H(K)}$ should be thought of as the Laplace Beltrami operators on $H(\mathfrak{k})$ and H(K) respectively, we should expect Eq. (1.11) to hold. This argument would be valid in finite dimensions but because of convergence issues it is not a real proof here. We refer the reader to Appendix *A* in Driver and Hall [20] for a detailed account.

The next proposition points out the well known difficulties when dealing with Laplacians in infinite dimensions.

Proposition 1.13. Let μ be standard Wiener measure on $\Omega := C([0,1],\mathbb{R})$. For a cylinder function of the form $f(\omega) = F(\omega|_{\pi})$, where

(1.12)
$$\pi := \{ 0 = s_0 < s_1 < \dots < s_n = 1 \}$$

is a partition of [0,1] and $F : \mathbb{R}^n \to \mathbb{R}$ is a smooth function which along with its derivative has at most polynomial growth, let

$$\Delta f = \triangle_{H(\mathbb{R})} f := \sum_{h \in S(\mathbb{R})} \partial_h^2 f$$

be as defined in Definition 1.6 above with $\mathfrak{k} = \mathbb{R}$. Then Δ is **not** a closable operator on $L^2(\Omega, \mu)$.

I would like to thank Brian Hall for the natural counter example to the closability of Δ used in the following proof.

Proof. Let

$$f(\omega) := \omega_1^2 - \frac{1}{2} = \int_0^1 \omega_s d\omega_s \text{ (Itô - integral)}$$

and for a partition π as in Eq. (1.12) let

$$f_{\pi}(\omega) := \sum_{j=0}^{n-1} \omega(t_j) \left(\omega(t_{j+1}) - \omega(t_j) \right).$$

Then by the basic theory of the Itô integral, $f_{\pi} \to f$ in $L^2(\mu)$ as $|\pi| \to 0$ where $|\pi|$ is the mesh size of π defined by

(1.13)
$$|\pi| := \max\left\{ |s_{i+1} - s_i| : i = 0, 1, 2, \dots, n-1 \right\}.$$

To compute Δf and Δf_{π} we will make use of the identity

$$\sum_{h \in S(\mathbb{R})} h(s)h(t) = \min(s, t)$$

which follows from the reproducing properties of $H(\mathbb{R})$ and Parseval's inequality, see Lemma 3.3 in [23] for example. By this identity and elementary computations,

$$\Delta f_{\pi} = \sum_{h \in S(\mathbb{R})} \partial_{h}^{2} f_{\pi} = 2 \sum_{h \in S(\mathbb{R})} \sum_{j=0}^{n-1} h(t_{j}) \left(h(t_{j+1}) - h(t_{j}) \right)$$
$$= 2 \sum_{j=0}^{n-1} \left[\min(t_{j}, t_{j+1}) - \min(t_{j}, t_{j}) \right] = 0$$

while

$$\Delta f = \sum_{h \in S(\mathbb{R})} \partial_h^2 f = 2 \sum_{h \in S(\mathbb{R})} h^2(1) = 2 \neq 0.$$

If Δ were closable, then $0 = \Delta f_{\pi} \to \Delta f = 2$ as $|\pi| \to 0$ which clearly does not happen.

Acknowledgment. It is a pleasure to thank Brian Hall, Yaozhong Hu, Masha Gordina and Harry Kesten for illuminating discussions relating to this paper. I would especially like to thank my former thesis advisor and more importantly my good friend, Len Gross, for his continued support and his never ending generosity both in mathematics and in life.

2. PATH GROUP RESULTS

Notation 2.1. If (Ω, \mathcal{F}, P) is a probability space and $f \in L^1(\Omega, \mu)$, we will often write $\int_{\Omega} f dP$ as P(f) or sometimes as $\mathbb{E}f$.

Given a Brownian motion $\{W_s : s \in [0, 1]\}$ on \mathfrak{k} with variance determined by $\langle \cdot, \cdot \rangle$, the process k defined by Fisk – Stratonovich stochastic differential equation

(2.1)
$$dk_s = k_s \circ \left(\sqrt{t}dW_s\right) \text{ with } k_0 = e.$$

is a diffusion on K with generator, $t\Delta/2$, where $\Delta = \sum_{A \in \beta} \tilde{A}^2$ where β is an orthonormal basis for \mathfrak{k} . Let dx denote Haar measure on K and

$$p_t(x) = e^{t\Delta/2}\delta_e(x) = \frac{d\text{Law}(k_1)}{dx}$$

be the convolution heat kernel on K.

Definition 2.2 (Wiener measure). The measure, $\mu_t = \text{Law}(k)$, is called **Wiener measure** on W(K) with variance t.

Remark 2.3. Using the Ad_K – invariance of $\langle \cdot, \cdot \rangle_{\mathfrak{k}}$ and Levy's criteria of Brownian motion, it is easily seen that $\mu_t = \operatorname{Law}(\tilde{k})$ as well, where \tilde{k} is the solution to the Fisk – Stratonovich stochastic differential equation

(2.2)
$$d\tilde{k}_s = \left(\sqrt{t}dW_s\right) \circ \tilde{k}_s \text{ with } k_0 = e.$$

Meta-Theorem 2.4 (Informal description of Wiener measure). The measure μ_t may be described *informally* as

(2.3)
$$d\mu_t(x) = \rho_t(x)\mathcal{D}x$$

where

(2.4)
$$\rho_t(x) = \frac{1}{Z_t} e^{-\frac{1}{2t}E_K(x)},$$

 $E_K(x)$ is the energy of x (see Eqs. (1.1) and (1.4)), $\mathcal{D}x$ is the formal Riemann volume measure ("Haar" measure) on H(K) and Z_t is the ill defined "normalization constant,"

(2.5)
$$"Z_t := \int_{W(K)} e^{-\frac{1}{2t}E(x)} \mathcal{D}x."$$

Moreover the "function" $\rho_t(x)$ "solves" the heat equation

(2.6)
$$\partial_t \rho_t = \frac{1}{2} \Delta_{H(K)} \rho_t \text{ with } \rho_0 = \delta_e$$

Proof. (Meta Proof.) We make use of the well known informal description of Wiener measure, P_t , with variance t on $W(\mathfrak{k})$. Namely $dP_t(\omega) = \gamma_t(\omega)\mathcal{D}\omega$ where

$$\gamma_t(\omega)\mathcal{D}\omega = \frac{1}{Z_t}e^{-\frac{1}{2t}E_{\mathfrak{k}}(\omega)}\mathcal{D}\omega,$$

 $\mathcal{D}\omega$ is the "Riemann volume measure" on $H(\mathfrak{k})$ and Z_t is a normalization "constant," see Kuo [47] for example. By analogy with the finite dimensional case we should expect the Gaussian "density," $\gamma_t(\omega)$, to be a solution to the heat equation

$$\partial_t \gamma_t = \frac{1}{2} \Delta_{H(\mathfrak{k})} \gamma_t \text{ with } \gamma_0 = \delta_0.$$

Since $B : H(K) \to H(\mathfrak{k})$ is an isometry of Riemannian manifolds, we conclude that $B_*\mathcal{D}x = \mathcal{D}\omega$, i.e. that *B* takes the Riemann volume element on H(K) to the Riemann volume element on $H(\mathfrak{k})$. Similarly by the definition of *B*, $E_{\mathfrak{k}}$ and E_K , $E_K(x) = E_{\mathfrak{k}}(B(x))$ and therefore, $\gamma_t \circ B = \rho_t$ and ρ_t should satisfy the heat Eq. (2.6) because of Theorem 1.12. Furthermore, by Eq. (2.2), $d\mu_t(x) = d\left(P_t \circ \tilde{B}\right)(x)$ where

$$\tilde{B}(x) := \int_0^{\cdot} dx_s \circ x_s^{-1}$$
 (Fisk – Stratonovich integral),

a "stochastic extension" of B to W(K). Therefore we should have

$$d\mu_t(x) = d\left(P_t \circ B\right)(x) = \gamma_t(B(x)) \cdot B_*^{-1} \mathcal{D}\omega = \rho_t(x) \mathcal{D}x$$

which formally proves Eq. (2.3).

Theorems 2.5 and 2.6 below give rigorous meaning to the results in the Meta-Theorem 2.4. In the first theorem we follow the ideas in Driver and Andersson [9] and replace H(K) by certain finite dimensional approximations.

Theorem 2.5. To each partition π of [0,1] as in Eq. (1.12), let

(2.7)
$$H_{\pi}(K) = \left\{ x \in H(K) : \frac{d}{ds} \left[x'(s)x(s)^{-1} \right] = 0 \text{ if } s \notin \pi \right\}$$
$$= \left\{ x \in H(K) : \frac{d^2}{ds^2} B_s(x) = 0 \text{ if } s \notin \pi \right\}$$

- the space of piecewise exponential maps on K. Then $H_{\pi}(K)$ is a submanifold of H(K) and as such inherits a Riemannian metric from H(K). Let λ_{π}^{K} denote the associated volume form on $H_{\pi}(K)$ and define

$$d\mu_t^{\pi}(x) := \frac{1}{Z_{\pi}} e^{-\frac{1}{2t}E_K(x)} d\lambda_{\pi}^K(x),$$

where Z_{π} is chosen so that μ_t^{π} is a probability measure on $H_{\pi}(K)$. Then $\lim_{|\pi|\to 0} \mu_t^{\pi} = \mu_t$ weakly where $|\pi|$ is defined in Eq. (1.13), i.e.

$$\lim_{|\pi|\to 0} \int_{H_{\pi}(K)} f d\mu_t^{\pi} = \int_{W(K)} f d\mu_t$$

for all bounded continuous functions $f: W(K) \to \mathbb{R}$. (Here we consider W(K) a metric space in the sup – norm topology.)

Proof. This theorem is a fairly direct consequence of Corollary 1.11 and a Wong – Zakai type approximation theorem for solving stochastic differential equations, see Theorem A.4 and Remark A.5 in Appendix A. Here are the details.

Let $H_{\pi}(\mathfrak{k})$ be the subspace of $H(\mathfrak{k})$ given by

$$H_{\pi}(\mathfrak{k}) = \left\{ \omega \in H(\mathfrak{k}) : \frac{d^2}{ds^2} \omega_s = 0 \text{ if } s \notin \pi \right\}.$$

We view $H_{\pi}(\mathfrak{k})$ as a Riemannian manifold with metric determined by the inner product on $H(\mathfrak{k})$, see Eq. (1.5). By the definition of $H_{\pi}(K)$ (Eq. (2.7)), $B(H_{\pi}(K)) = H_{\pi}(\mathfrak{k})$ and by Corollary 1.11, $B: H_{\pi}(K) \to H_{\pi}(\mathfrak{k})$ is an isometry of Riemannian manifolds. Therefore $B_*\lambda_{\pi}^K = \lambda_{\pi}^{\mathfrak{k}}$ where $\lambda_{\pi}^{\mathfrak{k}}$ is the Riemannian volume measure on $H_{\pi}(\mathfrak{k})$. From these remarks and the identity $E_K = E_{\mathfrak{k}} \circ B$, which follows from Eq. (1.4) and the definition of B, we find

(2.8)
$$d(B_*\mu_t^{\pi}) = \frac{1}{Z_{\pi}} e^{-\frac{1}{2t}E_K \circ B} d(B_*\lambda_{\pi}^K) = \frac{1}{Z_{\pi}} e^{-\frac{1}{2t}E_{\mathfrak{k}}} d\lambda_{\pi}^{\mathfrak{k}}$$

where Z_{π} is a normalization constant chosen to make $B_*\mu_t^{\pi}$ a probability measure.

As in Appendix A, let $\{W_s^{\pi} : s \in [0,1]\}$ be the piecewise linear process defined by

$$W_s^{\pi} = W_{s_j} + (s - s_j) \frac{W_{s_{j+1}} - W_{s_j}}{s_{j+1} - s_j} \text{ when } s \in [s_j, s_{j+1}],$$

where W is the \mathfrak{k} – valued Brownian motion used to define k in Eq. (2.1). If $g: H_{\pi}(\mathfrak{k}) \to \mathbb{R}$ is a bounded measurable function, then using the Gaussian finite dimensional distributions of W, we have

(2.9)
$$\mathbb{E}\left[g(\sqrt{t}W^{\pi})\right] = \int_{\mathfrak{g}^n} g(\sqrt{t}\omega^{\pi}_{(\xi_1,\ldots,\xi_n)}) \left[\prod_{j=1}^n p^{\mathfrak{e}}_{s_{j+1}-s_j}(\xi_{j+1}-\xi_j)\right] d\xi_1 \ldots d\xi_n$$

where $p_s^{\mathfrak{k}}(\xi) = (2\pi s)^{-\dim \mathfrak{k}/2} \exp\left(-\frac{1}{2s}|\xi|^2\right)$ and $\omega_{(\xi_1,\ldots,\xi_n)}^{\pi} \in H_{\pi}(\mathfrak{k})$ is determined uniquely by $\omega_{(\xi_1,\ldots,\xi_n)}^{\pi}(s_j) = \xi_j$ for $j = 1, 2, \ldots, n$. Using $\sqrt{t}\omega_{(\xi_1,\ldots,\xi_n)}^{\pi} = \omega_{(\sqrt{t}\xi_1,\ldots,\sqrt{t}\xi_n)}^{\pi}$, we may make a change of variables in Eq. (2.9) to find

(2.10)
$$\mathbb{E}\left[g(\sqrt{t}W^{\pi})\right] = \int_{\mathfrak{k}^n} g(\omega_{(\xi_1,\dots,\xi_n)}^{\pi}) \left[\prod_{j=1}^n p_{t(s_{j+1}-s_j)}^{\mathfrak{k}}(\xi_{j+1}-\xi_j)\right] d\xi_1\dots d\xi_n.$$

Simple algebra shows

$$\begin{split} \prod_{j=1}^{n} p_{t(s_{j+1}-s_{j})}^{\mathfrak{k}}(\xi_{j+1}-\xi_{j}) &= C_{\pi} \exp\left(-\frac{1}{2t} \sum_{j=1}^{n} \frac{|\xi_{j+1}-\xi_{j}|^{2}}{s_{j+1}-s_{j}}\right) \\ &= C_{\pi} \exp\left(-\frac{1}{2t} \sum_{j=1}^{n} \frac{|\xi_{j+1}-\xi_{j}|^{2}}{(s_{j+1}-s_{j})^{2}} \left(s_{j+1}-s_{j}\right)\right) \\ &= C_{\pi} \exp\left(-\frac{1}{2t} E_{\mathfrak{k}}(\omega_{(\xi_{1},\dots,\xi_{n})}^{\pi})\right), \end{split}$$

where

$$C_{\pi} := \prod_{j=1}^{n} \left(\frac{1}{2\pi t \left(s_{j+1} - s_{j}\right)}\right)^{\dim \mathfrak{k}/2}.$$

Since $(\xi_1, \ldots, \xi_n) \in \mathfrak{k}^n \xrightarrow{\phi} \omega_{(\xi_1, \ldots, \xi_n)}^{\pi} \in H_{\pi}(\mathfrak{k})$ is linear, $\phi_*(d\xi_1 \ldots d\xi_n)$ is a translation invariant measure on $H_{\pi}(\mathfrak{k})$ as is $\lambda_{\pi}^{\mathfrak{k}}$, so $\phi_*(d\xi_1 \ldots d\xi_n) = c_{\pi}d\lambda_{\pi}^{\mathfrak{k}}$ for some constant c_{π} . Combining all of these observations shows that Eq. (2.10) may be written as

$$\mathbb{E}\left[g(\sqrt{t}W^{\pi})\right] = \int_{\mathfrak{k}^n} g(\omega_{(\xi_1,\dots,\xi_n)}^{\pi}) C_{\pi} \exp\left(-\frac{1}{2t} E_{\mathfrak{k}}(\omega_{(\xi_1,\dots,\xi_n)}^{\pi})\right) d\xi_1 \dots d\xi_n$$

$$(2.11) \qquad \qquad = \int_{H_{\pi}(\mathfrak{k})} g(\omega) \frac{1}{Z_{\pi}} \exp\left(-\frac{1}{2t} E_{\mathfrak{k}}(\omega)\right) d\lambda_{\pi}^{\mathfrak{k}}(\omega)$$

where $Z_{\pi}^{-1} = C_{\pi}c_{\pi}$ is a normalization constant such that $\frac{1}{Z_{\pi}}\exp\left(-\frac{1}{2t}E_{\mathfrak{k}}(\omega)\right)d\lambda_{\pi}^{\mathfrak{k}}(\omega)$ is a probability measure as can be seen by taking g = 1 in Eq. (2.11). Combining Eq. (2.11) with Eq. (2.8) gives

$$\int_{H_{\pi}(K)} g(B(x)) d\mu_t^{\pi}(x) = \mathbb{E}\left[g(\sqrt{t}W^{\pi})\right]$$

and applying this formula to $g := f \circ B^{-1}$ shows

(2.12)
$$\int_{H_{\pi}(K)} f(x) d\mu_t^{\pi}(x) = \mathbb{E}\left[f \circ B^{-1}(\sqrt{t}W^{\pi})\right].$$

Using the definition of $B, k^{\pi} := B^{-1}(\sqrt{t}W^{\pi})$ satisfies

$$\sqrt{t}W_s^{\pi} = B_s(k^{\pi}) = \int_0^s (k^{\pi}(\sigma))^{-1} \frac{dk^{\pi}(\sigma)}{d\sigma} d\sigma$$

or equivalently

$$\frac{dk^{\pi}(s)}{ds} = k^{\pi}(s)\sqrt{t}\frac{dW_s^{\pi}}{ds} \text{ with } k^{\pi}(0) = e \in K.$$

So by the dominated convergence theorem and a Wong – Zakai type approximation Theorem A.4 and Remark A.5 in Appendix A below,

$$\lim_{|\pi| \to 0} \int_{H_{\pi}(K)} f(x) d\mu_t^{\pi}(x) = \lim_{|\pi| \to 0} \mathbb{E}\left[f(k^{\pi})\right] = \mathbb{E}\left[f(k)\right] = \int_W f(x) d\mu_t(x)$$

where k is the solution to Eq. (2.1). \blacksquare

~

2.1. Consequences of Meta – Theorem 2.4.

Theorem 2.6. If $f: W(K) \to \mathbb{R}$ is a smooth cylinder function then

$$G(t,x) := \int_{W(K)} f(xy^{-1})d\mu_t(y)$$

is a solution to the heat equation,

(2.13)
$$\left(\partial_t - \frac{1}{2}\Delta_{H(K)}\right)G(t,x) = 0 \text{ with } \lim_{t\downarrow 0} G(t,x) = f(x).$$

Proof. (Meta-proof) A rigorous proof of this theorem may be found in Driver and Srimurthy [23]. Here we will give a heuristic proof based on Meta-Theorem 2.4 which implies

$$G(t,x) = \int_{W(K)} f(xy^{-1})\rho_t(y)\mathcal{D}y.$$

So (formally)

$$\begin{split} \frac{1}{2} \Delta_{H(K)} G(t,x) &= \frac{1}{2} \sum_{h \in S} \frac{d^2}{d\tau^2} |_0 G(t,xe^{\tau h}) = \frac{1}{2} \int_{W(K)} \sum_{h \in S} \frac{d^2}{d\tau^2} |_0 f(xe^{\tau h}y^{-1}) \rho_t(y) \mathcal{D}y \\ &= \frac{1}{2} \int_{W(K)} \sum_{h \in S} \frac{d^2}{d\tau^2} |_0 \left[f(xy^{-1}) \rho_t(ye^{\tau h}) \right] \mathcal{D}y \\ &= \int_{W(K)} f(xy^{-1}) \frac{1}{2} \Delta_{H(K)} \rho_t(y) \mathcal{D}y = \int_{W(K)} f(xy^{-1}) \partial_t \rho_t(y) \mathcal{D}y \\ &= \partial_t G(t,x) \end{split}$$

where in the third equality we used the "invariance" of $\mathcal{D}y$ to make the change of variables, $y \to y e^{\tau h}$ and we used Eq. (2.6) in the second to last equality. In analogy with finite dimensional Laplace's method we expect $\lim_{t\downarrow 0} G(t, x) = f(x)$.

The following Logarithmic Sobolev inequality appears in Gross [36].

Theorem 2.7 (Logarithmic Sobolev Inequality on W(K)). For all non-zero smooth cylinder functions, $f: W(K) \to \mathbb{R}$,

(2.14)
$$\int_{W(K)} f^2 \log\left(\frac{f^2}{\mu_t(f^2)}\right) d\mu_t \le 2t \int_{W(K)} \|gradf\|^2 d\mu_t,$$

where

$$\mu_t(f^2) := \int_{W(K)} f^2 d\mu_t.$$

For analogous results when K is replaced by a compact Riemannian manifold, see Fang [27] (for the weaker Poincaré inequality) and Hsu [42], Aida and Elworthy [5] and Capitaine, Hsu, and Ledoux [12].

Proof. (Meta-Proof) Since $B : H(K) \to H(\mathfrak{k})$ is an isometry, it preserves all Riemannian geometric structures. Therefore the original logarithmic Sobolev inequality on $W(\mathfrak{k})$ proved by Gross in [35] transfers directly over to W(K) to give Eq. (2.14), Incidentally, in this case the Meta-proof and the real proof are almost the same, see Gross [36].

2.1.1. Generalized Segal – Bargmann – Hall Transform. Another "application" of Meta – Theorem 2.4 and Theorem 1.12 is that they give a simple way to understand the group generalizations of the Segal-Bargmann theorem, see the pioneering work of Hall [39]. Also see Gross and Malliavin [38] and [53, 41, 7, 20].

The classical Segal-Bargmann transform is a unitary isomorphism $S_t^{\mathfrak{k}}$: $L^2(W(\mathfrak{k}), P_t) \to \mathcal{H}L^2(W(\mathfrak{k}_{\mathbb{C}}), P_t^{\mathbb{C}})$ where $\mathcal{H}L^2(W(\mathfrak{k}_{\mathbb{C}}), P_t^{\mathbb{C}})$ denotes the "holomorphic" L^2 – functions on the path space, $W(\mathfrak{k}_{\mathbb{C}})$, into the complexified Lie algebra $\mathfrak{k}_{\mathbb{C}}$. The measure $P_t^{\mathbb{C}}$ is again the law of a certain $\mathfrak{k}_{\mathbb{C}}$ – valued Brownian motion which we do not bother to describe here. What is important is that, on cylinder functions,

$$S_t^{\mathfrak{k}}f = \left(e^{\frac{t}{2}\Delta_{H(\mathfrak{k})}}f\right)_a,$$

where, for a function $F: W(\mathfrak{k}) \to \mathbb{C}$, F_a refers to the analytic continuation of F to $W(\mathfrak{k}_{\mathbb{C}})$. The W(K) – version of this theorem is the existence of a unitary isomorphism

$$S_t^K : L^2(W(K), \mu_t) \to \mathcal{H}L^2(W(K_{\mathbb{C}}), \mu_t^{\mathbb{C}})$$

where $\mathcal{H}L^2(W(K_{\mathbb{C}}), \mu_t^{\mathbb{C}})$ denotes the "holomorphic" L^2 – functions on the path space, $W(K_{\mathbb{C}})$, into the complexified Lie group $K_{\mathbb{C}}$. The measure $\mu_t^{\mathbb{C}}$ is again the law of a certain $K_{\mathbb{C}}$ – valued Brownian motion which we do not bother to describe. On cylinder functions, $f: W(K) \to \mathbb{C}, S_t^K f: W(K_{\mathbb{C}}) \to \mathbb{C}$ is given analogously to $S_t^{\mathfrak{k}} f$ as

$$S_t^K f = \left(e^{\frac{t}{2} \Delta_{H(K)}} f \right)_a$$

Moreover, one has the following commutative diagram

$$\begin{array}{cccc} & L^2(W(\mathfrak{k}), P_t) &\cong & L^2(W(K), \mu_t)) \\ S_t^{\mathfrak{k}} & \downarrow & \downarrow & S_t^K \\ & \mathcal{H}L^2(W(\mathfrak{k}_{\mathbb{C}}), P_t^{\mathbb{C}}) &\cong & \mathcal{H}L^2(W(K_{\mathbb{C}}), \mu_t^{\mathbb{C}}) \end{array}$$

where the horizontal equivalences are determined by Itô maps as in (or similar to) Eq. (2.2).

These results are easily understood in terms of Theorem 1.12. Indeed from Eq. (1.11) we expect,

$$S_t^K f = (e^{\frac{t}{2}\Delta_{H(K)}} f)_a = \left(\left[e^{\frac{t}{2}\Delta_{H(\mathfrak{k})}} (f \circ B^{-1}) \right] \circ B \right)_a = \left(e^{\frac{t}{2}\Delta_{H(\mathfrak{k})}} (f \circ B^{-1}) \right)_a \circ B^{\mathbb{C}},$$

where $B^{\mathbb{C}} : H(K^{\mathbb{C}}) \to H(\mathfrak{k}^{\mathbb{C}})$ is a complex version of *B*. Because of domain issues this is not a rigorous proof, but the spirit is correct. We refer the reader to Driver and Hall [20] for the full details and further generalizations.

3. PINNED WIENER MEASURE RESULTS

We now wish to consider analogous theorems for loop groups where it no longer true that $H_0(K) \subset H(K)$ is a flat Riemannian manifold. However, we can still exploit the fact that $H_0(K)$ is an embedded submanifold of H(K).

3.1. Pinned Wiener Measure.

Definition 3.1 (Pinned Wiener Measure). Pinned Wiener measure μ_t^e is the measured μ_t "conditioned" to live on $\mathcal{L}(K)$. Informally, μ_t^e is given by

$$\mu_t^e(f) = \mathbb{E}(f(k)|k_1 = e) = \frac{\int_{W(K)} f(k)\delta_e(k_1)d\mu_t(k)}{\int_{W(K)} \delta_e(k_1)d\mu_t(k)} = \frac{\int_{W(K)} f(k)\delta_e(k_1)d\mu_t(k)}{p_t(e)}.$$

A rigorous definition of μ_t^e may be given using Doob's construction, namely, μ_t^e is the unique measure on $\mathcal{L}(K)$ such that if f is a bounded measurable function depending on $\sigma|_{[0,\alpha]}$ for some $\alpha \in (0, 1)$, then

$$\mu_t^e(f) \equiv \frac{1}{p_t^K(e)} \mu_t(f p_{t(1-\alpha)}^K(\pi_\alpha)).$$

For more details, see [8], [50], [37], and [19] for example.

We now want to develop the pinned analogues of Meta-Theorem 2.4. For this we will make formal use of the co-area formula described in the next theorem.

Theorem 3.2 (Co-area Formula). Suppose that $\pi : M \to N$ is a C^{∞} – map of smooth finite dimensional Riemannian manifolds and that $n \in N$ is a regular value of π . (Recall that the implicit function theorem then implies $M_n := \pi^{-1}(\{n\})$ is a smooth submanifold of M.) Let λ_M and λ_{M_n} denote the Riemann volume measure on M and M_n respectively. Then

$$\int_{M} f(x)\delta_n(\pi(x))d\lambda_M(x) = \int_{M_n} f(x)\frac{1}{\sqrt{\det(\pi_{*x}\pi_{*x}^{\mathrm{tr}})}}d\lambda_{M_n}(x)$$

where $\pi_{*x}: T_x M \to T_{\pi(x)} N$ is the differential of π and π_{*x}^{tr} is the adjoint of π_{*x} relative to the Riemannian structures on $T_x M$ and $T_{\pi(x)} N$.

Proof. I will not attempt to give a general proof of this fact here but only indicate the main ideas. First off, it suffices to prove the result for functions which have support in a small neighborhood of a point $m \in M_n$. By choosing appropriate coordinates on M and N, we may assume that both M and N are vector spaces and m = 0, and n = 0. We may also assume that π is linear. Let $X = Nul(\pi) = \pi^{-1}(\{0\}) = M_n$ and $Y = Nul(\pi)^{\perp} = Ran(\pi^*)$. Then

$$\begin{split} \int_{M} f \delta_{n}(\pi) d\lambda_{M} &= \int_{X \times Y} f(x+y) \delta(\pi(y)) dx dy \\ &= \int_{X \times Y} f(x) \delta(\pi(y)) dx dy \\ &= \int_{X} f(x) dx \int_{Y} \delta(\pi(y)) dy. \end{split}$$

Now letting $z = \pi(y)$, we find that $dz = |\det(\pi)| dy$, so that

$$\int_{Y} \delta(\pi(y)) dy = \frac{1}{|\det(\pi)|} \int_{N} \delta(z) dy = \frac{1}{|\det(\pi)|}$$

In order to compute the Jacobian factor $det(\pi)$, choose an orthogonal transformation $U: N \to Y$, then

$$|\det(\pi)| = |\det(\pi U)| = \det(\pi U (\pi U)^{\mathrm{tr}})^{1/2} = \det(\pi \pi^{\mathrm{tr}})^{1/2}$$

We will now formally apply Theorem 3.2 with M = H(K), N = K and $\pi : H(K) \to K$ the projection map $\pi(x) := x_1$.

Meta-Theorem 3.3 (Informal Description of Pinned Wiener). *Pinned Wiener* measure, μ_t^e , on $\mathcal{L}(K)$ is given informally by the expression

(3.1)
$$d\mu_t^e(x) = \frac{1}{p_t^K(e)}\rho_t(x)\mathcal{D}_0x.$$

where $\mathcal{D}_0 x$ is the Riemann volume "measure" on $H_0(K)$.

Proof. (Meta-Proof) By Theorem 3.2 we have (formally),

$$d\mu_t^e(x) = \frac{1}{p_t^K(e)} \rho_t(x) \delta_e(x_1) \mathcal{D}x = \frac{1}{p_t^K(e)} \frac{1}{\sqrt{\det(\pi_{*x} \pi_{*x}^{\mathrm{tr}})}} \rho_t(x) \mathcal{D}_0 x.$$

So to finish the proof it suffices to show, $det(\pi_{*x}\pi_{*x}^{tr}) = 1$.

Let $h \in H(\mathfrak{k})$, then $L_{x*}h \in T_xH(K)$ and

$$\pi_* L_{x*} h = \frac{d}{dt} |_0 \pi(x e^{th}) = \frac{d}{dt} |_0 (x(1) e^{th(1)}) = L_{x(1)*} h(1).$$

So if $\xi \in \mathfrak{k}$,

$$\langle \pi_* L_{x*}h, L_{x(1)*}\xi \rangle = \langle h(1), \xi \rangle = (h, s \to s\xi)_{H(\mathfrak{k})}$$

which implies $\pi_{*x}^{\text{tr}} L_{x(1)*} \xi = L_{x*}(s \to s\xi)$. Therefore

$$\pi_{*x}\pi_{*x}^{\mathrm{tr}}L_{x(1)*}\xi = \pi_{*x}L_{x*}(s \to s\xi) = L_{x(1)*}\xi,$$

i.e. $\pi_{*x}\pi_{*x}^{\text{tr}} = id_{T_{x(1)}K}$ and thus $\det(\pi_{*x}\pi_{*x}^{\text{tr}}) = 1$.

This Meta-Theorem suggests the following quasi – invariance theorem.

Theorem 3.4 (Malliavin & Malliavin [50]). Suppose that $k \in H_0(K)$, then

$$d\mu_t^e(xk) = J_k(x)d\mu_t^e(x)$$

where

$$J_{k} = \exp\left(-\frac{1}{t} \int_{0}^{1} \langle k'(s)k^{-1}(s), db_{s} \rangle - \frac{1}{2t} \int_{0}^{1} \left|k'(s)k^{-1}(s)\right|^{2} ds\right)$$
$$= \exp\left(-\frac{1}{t} \int_{0}^{1} \langle k'(s)k^{-1}(s), db_{s} \rangle - \frac{1}{2t} E(k)\right)$$

and $b_s(x) := \int_0^s x^{-1}(r) \circ dx(r)$. Here $\circ dx(r)$ is used to denote the Fisk – Stratonovich differential of x.

Proof. (Meta-Proof) Formally using Eq. (3.1) and the "invariance" of $\mathcal{D}_0 x$,

$$d\mu_t^e(xk) = \frac{1}{p_t^K(e)} \rho_t(xk) \mathcal{D}_0 x = \frac{\rho_t(xk)}{\rho_t(x)} d\mu_t^e(x) =: J_k(x) d\mu_t^e(x)$$

where

$$J_{k}(x) = \frac{\rho_{t}(xk)}{\rho_{t}(x)} = \frac{e^{-\frac{1}{2t}E_{K}(xk)}}{e^{-\frac{1}{2t}E_{K}(x)}} = \exp\left(-\frac{1}{2t}\left[E_{K}(xk) - E_{K}(x)\right]\right)$$
$$= \exp\left(-\frac{1}{2t}\left[E_{\mathfrak{k}}(b(xk)) - E_{\mathfrak{k}}(b(x))\right]\right).$$

Using the Ad_K invariance of $\langle \cdot, \cdot \rangle_{\mathfrak{k}}$,

$$\begin{split} E_K(xk) &= \int_0^1 \left| (x(s)k(s))^{-1} \frac{d}{ds} (x(s)k(s)) \right|_{\mathfrak{k}}^2 ds \\ &= \int_0^1 \left| k(s)^{-1} x(s)^{-1} \left[x'(s)k(s) + x(s)k'(s) \right] \right|_{\mathfrak{k}}^2 ds \\ &= \int_0^1 \left| Ad_{k(s)^{-1}} \left[x(s)^{-1} x'(s) \right] + k(s)^{-1}k'(s) \right|_{\mathfrak{k}}^2 ds \\ &= E_K(x) + 2 \int_0^1 \langle Ad_{k(s)^{-1}} \left[x(s)^{-1} x'(s) \right], k(s)^{-1}k'(s) \rangle_{\mathfrak{k}} ds + E(k) \\ &= E_K(x) + 2 \int_0^1 \langle b'_s(x), Ad_{k(s)} \left[k(s)^{-1}k'(s) \right] \rangle_{\mathfrak{k}} ds + E(k). \end{split}$$

Therefore,

$$J_k(x) = \exp\left(-\frac{1}{2t} \left[2\int_0^1 \langle b'_s(x), k'(s)k(s)^{-1} \rangle_{\mathfrak{k}} \, ds + E(k)\right]\right) \\ = \exp\left(-\frac{1}{t}\int_0^1 \langle k'(s)k^{-1}(s), db_s(x) \rangle - \frac{1}{2t}E(k)\right).$$

For a rigorous proof see Malliavin [50], Shigekawa [55, 54], Driver [15, 16], Hsu [43], Enchev and Stroock [25, 26] and Lyons and Qian[49]. Also see Albeverio and Hoegh-Krohn [8] for the path group analogue. ■

We now wish to find the "heat equation" solved by μ_t^e .

Meta-Theorem 3.5. The "function" $\tilde{\mu}_t^e(x) := \frac{1}{p_t^K(e)} \rho_t(x)$ on $H_0(K)$ solves the heat equation with potential,

(3.2)
$$\partial_t \tilde{\mu}_t^e(x) = \left(\frac{1}{2} \triangle_{\mathcal{L}(K)} + V_t\right) \tilde{\mu}_t^e \text{ with } \tilde{\mu}_t^e \to \delta_{\mathbf{e}} \text{ as } t \downarrow 0$$

where

(3.3)
$$V_t := \frac{1}{2t^2} |b_1|_{\mathfrak{k}}^2 - \left(\frac{\dim \mathfrak{k}}{2t} + \partial_t \log p_t^K(e)\right)$$

and b is defined in Eq. (1.7).

Proof. (Meta-Proof.) Let $\beta = \{\xi_i : i = 1, ..., \dim \mathfrak{k}\}$ be an orthonormal basis for \mathfrak{k} and let $h_i(s) = s\xi_i$. It is easily seen that $S_0 \cup \{h_i : i = 1, ..., \dim \mathfrak{k}\}$ is an orthonormal basis for $H(\mathfrak{k})$ and therefore,

$$\Delta_{H(K)} = \Delta_{H_0(K)} + \sum_{i=1}^{\dim(\mathfrak{k})} \tilde{h}_i^2.$$

Using these remarks and Eq. (2.6),

$$(\partial_t - \Delta_{H_0(K)}/2) \tilde{\mu}_t^e = (\partial_t - \Delta_{H(K)}/2) \tilde{\mu}_t^e + \frac{1}{2} \sum_{i=1}^{\dim(\mathfrak{k})} \tilde{h}_i^2 \tilde{\mu}_t^e$$

$$= \left(\partial_t \frac{1}{p_t^K(e)}\right) \rho_t + \frac{1}{p_t^K(e)} \left(\partial_t - \Delta_{H(K)}/2\right) \rho_t + \frac{1}{p_t^K(e)} \frac{1}{2} \sum_{i=1}^{\dim(\mathfrak{k})} \tilde{h}_i^2 \rho_t$$

$$= -\partial_t \ln(p_t^K(e)) \cdot \tilde{\mu}_t^e + \frac{1}{p_t^K(e)} \frac{1}{2} \sum_{i=1}^{\dim(\mathfrak{k})} \tilde{h}_i^2 \rho_t.$$

Moreover,

$$\begin{split} \tilde{h}_i^2 \rho_t &= \tilde{h}_i \left(-\rho_t \frac{1}{2t} \tilde{h}_i E_K \right) \\ &= \left[\frac{1}{4t^2} \left(\tilde{h}_i E_K \right)^2 - \frac{1}{2t} \tilde{h}_i^2 E_K \right] \rho_t, \end{split}$$

and therefore

$$\left(\partial_t - \Delta_{H_0(K)}/2\right)\tilde{\mu}_t^e = V_t\tilde{\mu}_t^e,$$

where

(3.4)
$$V_t = \frac{1}{2} \sum_{i=1}^{\dim(\mathfrak{k})} \left[\frac{1}{4t^2} \left(\tilde{h}_i E_K \right)^2 - \frac{1}{2t} \tilde{h}_i^2 E_K \right] - \partial_t \ln(p_t^K(e)).$$

Let us now proceed to work on Eq. (3.4). Making use of Eqs. (1.1), (1.7) and (1.10) we find

$$\tilde{h}_{i}E_{K}(x) = \tilde{h}_{i}E_{\mathfrak{k}}(b(x)) = \tilde{h}_{i}\int_{0}^{1}|b_{s}'(x)|^{2} ds$$

$$= 2\int_{0}^{1}\langle b_{s}'(x), \tilde{h}_{i}b_{s}'(x)\rangle ds = 2\int_{0}^{1}\langle b_{s}'(x), h_{i}'(s) + [b_{s}'(x), h_{i}(s)]\rangle ds$$

$$= 2\int_{0}^{1}\langle b_{s}'(x), h_{i}'(s)\rangle ds = 2\langle b_{1}(x), \xi_{i}\rangle$$
(3.5)

and

(3.6)
$$\tilde{h}_{i}^{2}E_{K}(x) = 2\int_{0}^{1} \langle h_{i}'(s) + [b_{s}'(x), h_{i}(s)], h_{i}'(s) \rangle ds = 2 |\xi_{i}|^{2} = 2$$

wherein we have used the Ad_K – invariance of $\langle \cdot, \cdot \rangle$ to conclude

$$\langle [b'_s(x), h_i(s)], h'_i(s) \rangle = s \langle [b'_s(x), \xi_i], \xi_i \rangle = -s \langle b'_s(x), [\xi_i, \xi_i] \rangle = 0.$$

Combining Eqs. (3.4), (3.5) and (3.6) shows V_t may be written as in Eq. (3.3).

The following theorem and corollary is a rigorous version of Meta – Theorem 3.5.

Theorem 3.6 (Airault & Malliavin, [6]). For any smooth cylindrical function $f : \mathcal{L}(K) \to \mathbb{R}$,

(3.7)
$$\partial_t \mu_t^e(f) = \mu_t^e \left[\left(\frac{1}{2} \Delta_{\mathcal{L}(K)} + V_t \right) f \right],$$

where $\triangle_{\mathcal{L}(K)} = \sum_{h \in S_0} \tilde{h}_i^2$,

(3.8)
$$V_t := \frac{1}{2t^2} \left| b_1 \right|_{\mathfrak{k}}^2 - \left(\frac{\dim \mathfrak{k}}{2t} + \partial_t \log p_t^K(e) \right) \text{ with }$$

(3.9)
$$b_1(x) := \int_0^1 x(s)^{-1} \circ dx(s) \, .$$

Proof. See Airault and Malliavin[6] and also Driver and Srimurthy [23] for a simplified proof. \blacksquare

Corollary 3.7. Let

(3.10)
$$F(t,x) := \int_{\mathcal{L}(K)} f(xy^{-1}) d\mu_t^e(y),$$

then

(3.11)
$$\partial_t F(t,x) = \frac{1}{2} \triangle_{\mathcal{L}(K)} F(t,x) + \int_{\mathcal{L}(K)} V_t(y) f(xy^{-1}) d\mu_t^e(y).$$

where V_t is defined in Eq. (3.8).

Proof. (Meta-proof. A rigorous proof of this theorem may be found in Driver and Srimurthy [23].) Writing (formally)

$$F(t,x) = \int_{H_0(K)} f(xy^{-1})\tilde{\mu}_t^e(y)\mathcal{D}_0 y_t$$

we find, as in the proof of Theorem 2.6, that

$$\left(\partial_t - \frac{1}{2}\Delta_{H_0(K)}\right)F(t,x) = \int_{H_0(K)} f(xy^{-1})\left(\partial_t - \frac{1}{2}\Delta_{H_0(K)}\right)\tilde{\mu}_t^e(y)\mathcal{D}_0 y$$

$$= \int_{H_0(K)} f(xy^{-1})V_t(y)\tilde{\mu}_t^e(y)\mathcal{D}_0 y$$

$$= \int_{\mathcal{L}(K)} V_t(y)f(xy^{-1})d\mu_t^e(y)$$

where in the second equality we have used Eq. (3.2).

3.2. Gross' Pioneering Results. Much of the work described above was an outgrowth of the following two fundamental results of Leonard Gross.

Theorem 3.8 (Gross' Ergodicity Theorem, [37]). Suppose $\int_{\mathcal{L}(K)} \|gradf\|^2 d\mu_t^e = 0$ then f is constant μ_t^e – a.e. on each homotopy class of $\mathcal{L}(K)$.

See Gross [37], Aida [1, 3], Leandre [48] and the article by Brian Hall [40] in this volume for more on Gross' ergodicity theorem.

Theorem 3.9 (Gross' Logarithmic Sobolev Inequality, [36]). For all $\lambda > 0$ there are constant C and B such that (3.12)

$$\int_{\mathcal{L}(K)} f^2 \log\left(\frac{f^2}{\mu_t^e(f^2)}\right) d\mu_t^e \le C \int_{\mathcal{L}(K)} \left\{ \|grad_0 f\|^2 + \left(\lambda |b_1|^2 + B\right) f^2 \right\} d\mu_t^e,$$

for all non-zero smooth cylinder functions, $f : \mathcal{L}(K) \to \mathbb{R}$.

Proof. The reader is referred to [36] for Gross' original proof. His proof is very geometrical and fits nicely into the picture presented in this article. The reader should also consult Getzler [31, 32], Aida [2], and Gong and Ma [33]. ■

Open Question: Is is possible to take $\lambda = 0$ in Eq. (3.12)?

Recent work of Eberle [24] has given examples of compact simply connected Riemannian manifolds where the analogue of Eq. (3.12) does **not** hold with $\lambda = 0$. These manifolds have a "dumb bell" shape unlike the very symmetric Lie group case. It is still not known if there exists any compact Riemannian manifolds so that Eq. (3.12) holds with $\lambda = 0$. However, it is known from [22] that Eq. (3.12) does hold when $d\mu_t^e$ is replaced by the heat kernel measure ν_t^e on $\mathcal{L}(K)$ (see Section 4 below). It also known by the work of Gong, Röckner, and Wu [34] that by modifying pinned Wiener measure by certain positive densities, the weaker Poincaré inequality can be made to hold on the loop space of any compact Riemannian manifold. These results make the question of whether Eq. (3.12) holds with $\lambda = 0$ all the more intriguing.

4. Comparing Heat Kernel and Pinned Wiener Measures

Definition 4.1 (Heat Kernel Measure). The heat kernel measures on $\mathcal{L}(K)$ are formed by the one parameter family of probability measures $\{\nu_t^e : t \ge 0\}$ satisfying $\nu_0^e = \lim_{t \ge 0} \nu_t^e = \delta_e$ and

$$\partial_t \nu_t^e(f) = \frac{1}{2} \nu_t^e \left(\triangle_{H_0(K)} f \right)$$

for all smooth cylinder functions $f : \mathcal{L}(K) \to \mathbb{R}$.

The reader is referred to Malliavin [51], Driver and Lohrenz [22], and Driver and [17] for the existence of ν_t^e . Our description of the results in this section will be rather brief compared to the previous sections. This is because to understand the heat kernel measure, ν_t^e , one must go to the path space of $\mathcal{L}(K)$ which is beyond the scope of this article. Nevertheless, I would like to include some basic properties of ν_t^e and its relationship to μ_t^e .

Theorem 4.2 (Heat Kernel Logarithmic Sobolev Theorem, [22]). There is a constant $C < \infty$ such that

(4.1)
$$\int_{\mathcal{L}(K)} f^2 \log \frac{f^2}{\nu_t^e(f^2)} d\nu_t^e \le C \int_{\mathcal{L}(K)} \|grad_0 f\|^2 d\nu_t^e$$

for all smooth cylinder functions $f : \mathcal{L}(K) \to \mathbb{R}$.

Proof. See Driver and Lohrenz [22], Carson [13, 14] and Fang [29] ■

Corollary 4.3 (Quasi-invariance for heat kernel measure). For each $k \in H_0(G)$ which is null homotopic, ν_t^e quasi-invariant under the right and left translations by k.

Proof. See Driver [17, 18] and Fang [28, 29]. The free loop space version of these results was carried out by Trevor Carson in [13, 14]. \blacksquare

Theorem 4.4 (V. Srimurthy,[57]). Let \mathcal{F}_s denote the σ – algebra on W(K) generated by the coordinate functions $x \in W(K) \to x(r) \in K$ for $r \leq s$. Then for any s < 1, the measures ν_t^e and μ_t^e are mutually absolutely continuous on \mathcal{F}_s .

Proof. This is proved in Srimurthy [57] using two parameter stochastic calculus along the lines developed by J. Norris [52]. \blacksquare

More generally we have the following two theorems.

Theorem 4.5 (Driver & Srimurthy [23]). Heat kernel measure ν_t^e is absolutely continuous relative to pinned Wiener measure μ_t^e for all t > 0. Moreover, $d\nu_t^e/d\mu_t^e$ is a bounded function on $\mathcal{L}(K)$.

Proof. (Sketch of the proof.) By standard heat kernel asymptotics, one shows $c_t := \frac{\dim \mathfrak{k}}{2t} + \partial_t \log p_t^K(e) = O(|\ln t|)$ as $t \downarrow 0$ so that

(4.2)
$$C_t := \int_0^t c_\tau d\tau = \int_0^t \left(\frac{\dim \mathfrak{k}}{2\tau} + \partial_\tau \log p_\tau^K(e)\right) d\tau$$

is a continuous function for $t \ge 0$. Let $f : \mathcal{L}(K) \to [0, \infty)$ be a cylinder function and define

$$F(t,x) := e^{C_t} \int_{\mathcal{L}(K)} f(xy^{-1}) d\mu_t^e(y)$$

and

$$G(t,x) := \int_{\mathcal{L}(K)} f(xy^{-1}) d\nu_t^e(y).$$

Then, by Corollary 3.7 and the definition of ν_t^e , F and G satisfy

$$\left(\partial_t - \frac{1}{2}\Delta_{H_0(K)}\right)F(t,x) = \int_{\mathcal{L}(K)} \left[V_t(y) + \dot{C}_t\right]f(xy^{-1})d\mu_t^e(y)$$
$$= \int_{\mathcal{L}(K)} \frac{1}{2t^2}|b_1(y)|^2f(xy^{-1})d\mu_t^e(y) \ge 0$$

and

$$\left(\partial_t - \frac{1}{2}\Delta_{H_0(K)}\right)G(t, x) = 0.$$

Since both G(0, x) = F(0, x), we may apply the Maximum principle (or Du Hamel's principle) to conclude that

$$F(t,x) \ge G(t,x)$$

for all (t, x). Setting $x = \mathbf{e}$ gives

(4.3)
$$e^{C_t} \int_{\mathcal{L}(K)} f(y) d\mu_t^e(y) \ge \int_{\mathcal{L}(K)} f(y) d\nu_t^e(y)$$

for all positive cylinder functions on $\mathcal{L}(K)$. By standard measure theoretic arguments, it then follows that Eq. (4.3) holds for all bounded measurable f from which the result easily follows.

Theorem 4.6 (Aida & Driver [4]). For each $h \in H_0(K)$, let $\nu_t^e(h, A) := \nu_t^e(h^{-1}A) - heat$ kernel measure on $\mathcal{L}(K)$ starting at h. Let $\Pi \subset H_0(K)$ be chosen so that to each homotopy class in $\mathcal{L}(K)$, there is a unique representative in Π . Then μ_t^e is absolutely continuous relative to $\sum_{h \in \Pi} \nu_t^e(h, \cdot)$. In particular, if K is simply connected, then μ_t^e and ν_t^e are mutually absolutely continuous.

Proof. I will only give a sketch of the proof when K is simply connected. By Theorem 4.5, $Z_t(x) = \frac{d\nu_t^e}{d\mu_t^e}(x)$ exists. Let $N := \{x \in \mathcal{L}(K) : Z_t(x) = 0\}$. Using the quasi-invariance of the pinned measure μ_t^e (Theorem 3.4) and of ν_t^e (Corollary 4.3) under left translation by $H_0(K)$, one shows that $h \cdot N = N$ modulo sets of μ_t^e –

measure zero for all $h \in H_0(K)$. Using Gross' ergodicity Theorem 3.8, it follows that either $\mu_t^e(N) = 0$ or that $\mu_t^e(\mathcal{L}(K) \setminus N) = 0$. Since ν_t^e is not the zero measure, we conclude that $\mu_t^e(N) = 0$, i.e. that Z(x) > 0 on $\mathcal{L}(K)$ for μ_t^e - a.e. x and hence μ_t^e is absolutely continuous relative to ν_t^e .

4.1. Feynman-Kac Formula for Pinned Wiener Measure. At this point one should expect when K is simply connected that $d\mu_t^e/d\nu_t^e$ can be expressed in terms of a Feynman – Kac type formula. This is technically problematic to prove owing to the singular behavior of the potential V_t as $t \downarrow 0$. Xiang - Dong Li and the author have been able to find partial results in this direction. The statement of these formula require the notion of Brownian motion on the $\mathcal{L}(K)$ and hence will be omitted. Let me mention though that these formula do begin to give more quantitative information about $d\mu_t^e/d\nu_t^e$ and can be used to give another proof of Theorem 4.5.

APPENDIX A. A WONG - ZAKAI APPROXIMATION THEOREM

A.1. Basic notation and the Theorem. Let $n \in \mathbb{N}$ and \mathbb{F} be either \mathbb{R} or \mathbb{C} and $T \in (0, \infty)$. Let gl_n denote the $n \times n$ matrices with entries from \mathbb{F} and GL_n – denote the invertible matrices in gl_n . For $A, B \in gl_n$ let $(A, B) = tr(A^*B)$ and $|A| := \sqrt{(A, A)}$. We will also use ||A|| to denote the operator norm of A of any bounded operator on an inner product space.

Remark A.1. We will use, without further comment, the following basic properties of the Hilbert Schmidt norm |A|,

- (1) If $I \in gl_n$ is the identity matrix, then $|I| = \sqrt{n}$.
- (2) $||A|| \le |A|$ and $|A^*| = |A|$.
- (3) If $A, B \in gl_n$ then $|AB| \le ||A|| ||B|| \le |A|| ||B||$ and $|AB| \le |A|| ||B|| \le |A|| ||B||$ and for all $p \ge 0$,

$$\left|Ae^B\right|^p \leq |A|^p \left\|e^B\right\|^p \leq |A|^p \, e^{p\|B\|} \leq |A|^p \, e^{p|B|}.$$

Now suppose that W_t is a gl_n – valued Brownian motion, i.e. $W_t = \sum_{A \in \beta} W_t^A A$ where $\beta \subset gl_n$ is an \mathbb{R} – linearly independent subset of gl_n and $\{W_t^A\}_{A \in \beta}$ is a collection of independent \mathbb{R} – valued Brownian motions.

A partition, π , of [0, T] is a set of the form

$$\pi = \{ 0 = t_0 < t_1 < t_2 < \dots < t_N = T \}.$$

For $t \in J_n := [t_n, t_{n+1})$ we will write t_- for t_n, t_+ for $t_{n+1}, \Delta_n := t_{n+1} - t_n, \Delta_n W := W(t_{n+1}) - W(t_n)$ and

$$W^{\pi}(t) := W(t_{-}) + (t - t_{-}) \left(\frac{W(t_{+}) - W(t_{-})}{t_{+} - t_{-}} \right).$$

We will also write $\Delta_t W := W(t_+) - W(t_-)$ and $\Delta_t = t_+ - t_-$. Notice that $\dot{W}^{\pi}(t) = \Delta_t W / \Delta_t$ for $t \notin \pi$. As usual, define

$$|\pi| := \max \{ \Delta_n : n = 0, 1, \dots, N - 1 \}$$

Definition A.2 (SDE). Let g_t be the GL_n – valued process solving the stochastic differential equation,

(A.1)
$$dg = g \circ dW \text{ with } g_0 = I$$
$$= gdW + \frac{1}{2}g(dW)^2 = gdW + \frac{1}{2}g\eta dt$$

where $\eta := \sum_{A \in \beta} A^2$.

Definition A.3 (Wong – Zakai Approximation). The Wong – Zakai approximation associated to g and a partition π is $g_{\pi}(t) = y(t)$ where y solves the ordinary differential equation,

$$\dot{y}(t) = y(t)\dot{W}^{\pi}(t)$$
 with $y(0) = I$.

Explicitly if $t \in J_n$ then

$$g_{\pi}(t) = e^{\Delta_0 W} e^{\Delta_1 W} \dots e^{\Delta_{n-1} W} e^{\frac{(t-t_n)}{\Delta_n} \Delta_n W}.$$

Theorem A.4 (Wong – Zakai Approximation Theorem). For any $p \ge 2$ and there is a constant $C < \infty$ such that

(A.2)
$$\sup_{t \in \pi} \mathbb{E} |g_{\pi}(t) - g_t|^p \le \mathbb{E} \sup_{t \in \pi} |g_{\pi}(t) - g_t|^p \le C |\pi|^{p/2}$$

and

(A.3)
$$\mathbb{E} \sup_{t \le T} |g_{\pi}(t) - g_t|^p \le C |\pi|^{(p/2)-1}$$

Remark A.5. A similar proof as below shows Theorem A.4 also holds if we assume g and g_{π} solve:

$$dg = dW \circ g$$
 with $g_0 = I$

and $g_{\pi} = y$ where y solves

$$\dot{y}(t) = \dot{W}^{\pi}(t)y(t)$$
 with $y(0) = I$.

In this case

$$g_{\pi}(t) = e^{\Delta_{n-1}W} e^{\frac{(t-t_n)}{\Delta_n}\Delta_n W} \dots e^{\Delta_1 W} e^{\Delta_0 W}$$

Theorems of this type have a long history starting with Wong and Zakai [59], also see Bismut [10], Ikeda and Watanabe [45], Kloeden and Platen [46] and Hu [44] to name a few references. Despite all of these references, it seems worthwhile to sketch a proof here in the Lie group case for two reasons. Firstly, the main ideas are more easily seen in this setting and secondly, it is still hard to find results of this nature when the coefficients of the stochastic differential equations are not bounded. This later case is needed when working on the complexified Lie groups as in Brian Hall's article. In a forthcoming paper [21], Hu and the author will give a more general version of the Theorem A.4 which will include a better convergence estimate than the one described in Eq. (A.3). The rest of this appendix will be devoted to the proof of Theorem A.4. Throughout the proof, the letter C will be used to denote a constant **not** depending on π . This constant may vary from line to line.

A.2. Bounds on the approximate solution.

Proposition A.6. Let $y = g_{\pi}$ as above. There exists $C < \infty$ such that

(A.4)
$$\sup_{t \in \pi} \mathbb{E} |y(t)|^p \le n^{p/2} e^{Cp^2 t/2} \ \forall \ 2 \le p \le |\pi|^{-1/2}$$

Moreover, for each $p\in [2,\infty)$ there is a constant $C_p<\infty$ such that

(A.5)
$$\mathbb{E}\left[\sup_{t\in J_n} |y(t) - y(t_n)|^p\right] \le C_p \Delta_n^{p/2} \le C_p |\pi|^{p/2} \quad \forall n = 0, 1, \dots, N-1.$$

Proof. Taylor theorem with integral remainder states,

(A.6)
$$f(t_+) - f(t_-) = \dot{f}(t_-)(t_+ - t_-) + \frac{1}{2}(t_+ - t_-)^2 \int_{t_-}^{t_+} \ddot{f}(t) d\nu(t)$$

where $d\nu(t) = 2(t_+ - t)/(t_+ - t_-)^2 dt$ – a probability measure on $[t_-, t_+]$. Let $f(t) = |y(t)|^p$, then elementary computations and estimates show,

$$\dot{f}(t) = p |y(t)|^{p-2} \operatorname{Re}(y(t), y(t)\dot{W}^{\pi}(t)),$$

$$\ddot{f}(t) = p(p-2) |y(t)|^{p-4} \left[\operatorname{Re}(y(t), y(t) \dot{W}^{\pi}(t)) \right]^2 + p |y(t)|^{p-2} \operatorname{Re}\left[(y(t) \dot{W}^{\pi}(t), y(t) \dot{W}^{\pi}(t)) + (y(t), y(t) \dot{W}^{\pi}(t) \dot{W}^{\pi}(t)) \right]$$

and

(A.7)
$$\left| \ddot{f}(t) \right| \leq (p(p-2)+2p) \left| y(t) \right|^p \left| \dot{W}^{\pi}(t) \right|^2 = p^2 \left| y(t) \right|^p \left| \dot{W}^{\pi}(t) \right|^2.$$

Therefore by Eq. (A.6)),

(A.8)
$$|y(t_+)|^p - |y(t_-)|^p = p |y(t_-)|^{p-2} \operatorname{Re}(y(t_-), y(t_-)\dot{W}^{\pi}(t_-)) + R$$

where the remainder,

$$R = \frac{1}{2}(t_{+} - t_{-})^{2} \int_{t_{-}}^{t_{+}} \frac{d^{2}}{dt^{2}} \left| y(t) \right|^{p} d\nu(t),$$

satisfies (by Eq. (A.7),

$$|R| \le \frac{1}{2}(t_{+} - t_{-})^{2} \int_{t_{-}}^{t_{+}} p^{2} |y(t)|^{p} \left| \dot{W}^{\pi}(t) \right|^{2} d\nu(t) = \frac{p^{2}}{2} \left| \Delta_{t} W \right|^{2} \int_{t_{-}}^{t_{+}} |y(t)|^{p} d\nu(t).$$

Taking expectations of Eq. (A.8) gives,

(A.9)
$$|\mathbb{E}|y(t_{+})|^{p} - \mathbb{E}|y(t_{-})|^{p}| = |\mathbb{E}R| \le \frac{p^{2}}{2} \int_{t_{-}}^{t_{+}} \mathbb{E}\left[|y(t)|^{p} |\Delta_{t}W|^{2}\right] d\nu(t)$$

Since

(A.10)
$$y(t) = y(t_{-}) \exp\left((t - t_{-}) \dot{W}^{\pi}(t)\right) = y(t_{-}) \exp\left(\frac{t - t_{-}}{t_{+} - t_{-}} \Delta_{t} W\right),$$

Remark A.1 implies

$$|y(t)|^{p} = |y(t_{-})|^{p} \exp\left(p\frac{t-t_{-}}{t_{+}-t_{-}} |\Delta_{t}W|\right) \le |y(t_{-})|^{p} e^{p|\Delta_{t}W|}$$

and combining this equation with Eq. (A.9) gives

$$\begin{aligned} |\mathbb{E} |y(t_{+})|^{p} - \mathbb{E} |y(t_{-})|^{p}| &\leq \frac{p^{2}}{2} \int_{t_{-}}^{t_{+}} \mathbb{E} \left[|y(t_{-})|^{p} e^{p|\Delta_{t}W|} |\Delta_{t}W|^{2} \right] d\nu(t) \\ &= \frac{p^{2}}{2} \mathbb{E} \left[|y(t_{-})|^{p} \right] \mathbb{E} \left[e^{p|\Delta_{t}W|} |\Delta_{t}W|^{2} \right]. \end{aligned}$$

Using Taylor's theorem with remainder and the fact that $s \to \mathbb{E}\left[e^{s|W_1|} |W_1|^2\right]$ is a smooth function,

$$\mathbb{E}\left[e^{p|\Delta_t W|} |\Delta_t W|^2\right] = \Delta_t \cdot \mathbb{E}\left[e^{p\sqrt{\Delta_t}|W_1|} |W_1|^2\right] = \Delta_t \left(C_0 + O\left(p\sqrt{\Delta_t}\right)\right) \le C\Delta_t$$

where the last inequality holds provided $0 \le p\sqrt{\Delta_t} \le 1$. Putting this all together implies,

$$\mathbb{E} |y(t_{+})|^{p} \leq \mathbb{E} |y(t_{-})|^{p} + C \frac{p^{2}}{2} \mathbb{E} [|y(t_{-})|^{p}] (t_{+} - t_{-})$$
$$= \mathbb{E} |y(t_{-})|^{p} [1 + C \frac{p^{2}}{2} (t_{+} - t_{-})] \leq \mathbb{E} |y(t_{-})|^{p} e^{Cp^{2} (t_{+} - t_{-})/2}$$

which upon iteration (using $y(0) = y(t_0) = I$) proves Eq. (A.4).

For Eq. (A.5), the fundamental theorem of calculus shows

$$\begin{aligned} |y(t) - y(t_{-})| &= \left| \int_{t_{-}}^{t} y(\tau) \dot{W}_{\pi}(\tau) d\tau \right| = \left| \int_{t_{-}}^{t} y(t_{-}) \exp\left((\tau - t_{-}) \dot{W}^{\pi}(\tau) \right) \right) \dot{W}_{\pi}(\tau) d\tau \\ &\leq \int_{t_{-}}^{t} |y(t_{-})| \left| \dot{W}_{\pi}(\tau) \right| \exp\left((\tau - t_{-}) \left| \dot{W}^{\pi}(\tau) \right) \right| \right) d\tau \\ &\leq |y(t_{-})| \left| \Delta_{t} W \right| e^{|\Delta_{t} W|}. \end{aligned}$$

and therefore

$$\mathbb{E}\left[\sup_{t\in J_n} |y(t) - y(t_n)|^p\right] \le \mathbb{E}\left|y(t_n)\right|^p \cdot \mathbb{E}\left(\left|\Delta_n W\right|^p e^{p|\Delta_n W|}\right) \le C_p \Delta_n^{p/2} \le C \left|\pi\right|^{p/2}.$$

A.3. Bounds on the solution.

Proposition A.7. If $p \in [2, \infty)$, there exists a constant $C_p < \infty$ such that

(A.11)
$$\mathbb{E} |g_t|^p \le \mathbb{E} \sup_{t \le T} |g_t|^p \le C_p e^{C_p t} \le C_p e^{C_p T} \text{ for all } t \in [0, T]$$

and

(A.12)
$$\mathbb{E} \sup_{t \in J_n} |g_t - g_{t_n}|^p \le C_p \Delta_n^{p/2} \le C_p |\pi|^{p/2} \quad \forall n = 0, 1, \dots, N-1.$$

for all n.

Proof. The estimate in Eq. (A.11) is very standard and follows from the Picard iteration method for solving Eq. (A.1). Alternatively one may prove this using Burkholder's and Gronwall's inequalities, see [58] for example. Let $q = \frac{p/2}{p/2-1} = \frac{p}{p-2}$ and $g_T^* := \sup_{t \leq T} |g_t|$. By the definition of g in Eq. (A.1),

$$g_t - g_{t_-} = \int_{t_-}^t g dW + \frac{1}{2} \int_{t_-}^t g \eta dt$$

and so by Burkholder's inequality

$$\mathbb{E} \sup_{t \in J_n} |g_t - g_{t_n}|^p \le C_p \left(\mathbb{E} \left[\int_{t_n}^{t_{n+1}} \sum_{A \in \beta} |g_t A|^2 dt \right]^{p/2} + \mathbb{E} \left[\int_{t_n}^{t_{n+1}} |g_t \eta| dt \right]^p \right) \\ \le C_p \left(\mathbb{E} \left(g_T^* \right)^p |t_{n+1} - t_n|^{p/2} + \mathbb{E} \left(g_T^* \right)^p |t_{n+1} - t_n|^p \right) \\ \le C_p \Delta_n^{p/2} \le C_p |\pi|^{p/2}$$

which proves Eq. (A.12). \blacksquare

A.4. Matrix exponential properties and estimates.

Notation A.8. For $A, B \in gl_n$ let

$$\partial_A e^B := \frac{d}{ds}|_0 e^{B+sA}$$
 and $\Lambda_B(A) := e^{-B} \partial_A e^B$

and

$$\partial_A^2 e^B := \frac{d^2}{ds^2}|_0 e^{B+sA} \text{ and } \Gamma_B(A,A) := e^{-B}\partial_A^2 e^B.$$

Proposition A.9. Using the notation above,

(A.13)
$$\Lambda_B(A) = \int_0^1 e^{-\tau a d_B} A d\tau$$

and

(A.14)
$$\Gamma_B(A,A) = [\Lambda_B(A)]^2 - \int_0^1 \int_0^1 \tau \left([e^{-s\tau a d_B} A, e^{-\tau a d_B} A] \right) d\tau ds$$

Moreover there exists $C \in (0, \infty)$ such that

(A.15)
$$\|\Lambda_B - I\| := \sup_{|A|=1} |\Lambda_B(A) - A| \le C |B| e^{C|B|}$$

and

(A.16)
$$\left| \Gamma_B(A,A) - A^2 \right| \le C \left(|B| + |B|^2 \right) e^{C|B|} |A|^2$$

Proof. Integrating the identity,

$$\frac{d}{dt} \left[\left(\partial_A e^{tB} \right) e^{-tB} \right] = \left[\partial_A \left(e^{tB} B \right) e^{-tB} - \left(\partial_A e^{tB} \right) B e^{-tB} \right]$$
$$= \left[\left(\partial_A e^{tB} \right) B e^{-tB} + e^{tB} A e^{-tB} - \left(\partial_A e^{tB} \right) B e^{-tB} \right]$$
$$= e^{tB} A e^{-tB},$$

shows

$$\partial_A e^B = \left[\int_0^1 e^{\tau B} A e^{-\tau B} d\tau \right] e^B = \int_0^1 e^{\tau B} A e^{(1-\tau)B} d\tau = \int_0^1 e^{(1-\tau)B} A e^{\tau B} d\tau$$
(A.17)
$$= e^B \int_0^1 e^{-\tau B} A e^{\tau B} d\tau = e^B \Lambda_B(A)$$

which proves Eq. (A.13). Differentiating Eq. (A.17) gives

(A.18)
$$\partial_A^2 e^B = e^B \left[\Lambda_B(A)\right]^2 + e^B \int_0^1 \left(\partial_A e^{-\tau a d_B}\right) A d\tau,$$

where $ad_AB := [A, B]$. Applying Eq. (A.13) to $e^{-\tau ad_B}$ and using $\partial_A(-\tau ad_B) = -\tau ad_A$ we learn

(A.19)
$$\partial_A e^{-\tau a d_B} = -\tau \int_0^1 e^{(s-1)\tau a d_B} a d_A e^{-s\tau a d_B} ds.$$

Eq. (A.14) follows from Eqs. (A.18) and (A.19) after a bit of algebra.

Now for the estimates in Eqs. (A.15) and (A.16). By the fundamental theorem of calculus

$$\|e^{-ad_B} - I\| = \left\| \int_0^1 -ad_B e^{-sad_B} ds \right\| \le \int_0^1 ds \, \|ad_B e^{-sad_B}\|$$
(A.20)
$$\le \int_0^1 ds \, \|ad_B\| \cdot \|e^{-sad_B}\| \le \|ad_B\| \cdot e^{\|ad_B\|} \le C \, |B| \, e^{C|B|}$$

and hence

$$\|\Lambda_B - I\| = \left\| \int_0^1 \left(e^{-\tau a d_B} - I \right) d\tau \right\| \le \int_0^1 \left\| e^{-\tau a d_B} - I \right\| d\tau \le C |B| e^{C|B|}$$

which proves Eq. (A.15).

By Eq. (A.20),

$$R_B(A) := A - e^{-ad_B}A$$

satisfies $||R_B|| \leq C |B| e^{C|B|}$. Because [A, A] = 0,

$$[e^{-s\tau ad_B}A, e^{-\tau ad_B}A] = [R_{-s\tau B}(A), A] + [A, R_{-\tau B}(A)] + [R_{-s\tau B}(A), R_{-\tau B}(A)]$$

and thus

$$\left| \int_{0}^{1} \int_{0}^{1} \tau \left([e^{-s\tau a d_{B}} A, e^{-\tau a d_{B}} A] \right) d\tau ds \right| \leq C \left(|B| + |B|^{2} \right) e^{C|B|} |A|^{2}.$$

Similarly,

$$\left| [\Lambda_B(A)]^2 - A^2 \right| = \left| [A + R_B(A)]^2 - A^2 \right|$$

= $\left| AR_B(A) + R_B(A)A + [R_B(A)]^2 \right|$
 $\leq C \left(|B| + |B|^2 \right) e^{C|B|} |A|^2.$

Combining these estimates with the definition of Γ_B in Eq. (A.14) proves Eq. (A.16).

A.5. **Proof of Theorem A.4. Proof.** We are now ready to complete the proof of Theorem A.4. Let us begin by introducing a third process:

$$z_t = z(t) = e^{\Delta_0 W} e^{\Delta_1 W} \dots e^{\Delta_{n-1} W} e^{W_t - W_{t_n}}$$
for $t \in J_n$.

Notice that z is adapted whereas $y = g_{\pi}$ is not. However z = y on π and $z_t = g_{\pi \cup \{t\}}(t)$.

Itô's formula (for $t \ge t_n$) implies

$$de^{W_t - W_{t_n}} = e^{W_t - W_{t_n}} \Lambda_{W_t - W_{t_n}} (dW_t) + \frac{1}{2} e^{W_t - W_{t_n}} \Gamma_{W_t - W_{t_n}} (dW_t, dW_t)$$

from which it follows that

(A.21)
$$dz_t = z_t \left[\Lambda_{W_t - W_{t_-}} (dW_t) + \frac{1}{2} \Gamma_{W_t - W_{t_-}} (dW_t, dW_t) \right].$$

We are now going to compare z and g. In order to do this we will need the SDE solved by $g^{-1}:$

(A.22)
$$dg_t^{-1} = -dW_t \circ g_t^{-1} = -dW_t g_t^{-1} + \frac{1}{2} (dW_t)^2 g_t^{-1} \text{ with } g_0 = I.$$

Letting $Q := zg^{-1}$ and using Eqs. (A.21) and (A.22) we find

$$\begin{split} dQ_t &= z_t \left[\Lambda_{W_t - W_{t_-}} (dW_t) + \frac{1}{2} \Gamma_{W_t - W_{t_-}} (dW_t, dW_t) \right] g_t^{-1} \\ &+ z_t \left[-dW_t + \frac{1}{2} (dW_t)^2 \right] g_t^{-1} - z_t \left[\Lambda_{W_t - W_{t_-}} (dW_t) \right] dW_t g_t^{-1} \\ &= z_t \left[\Lambda_{W_t - W_{t_-}} (dW_t) - dW_t \right] g_t^{-1} \\ &+ z_t \left[\frac{1}{2} \left\{ \Gamma_{W_t - W_{t_-}} (dW_t, dW_t) + (dW_t)^2 \right\} - \left[\Lambda_{W_t - W_{t_-}} (dW_t) \right] dW_t \right] g_t^{-1} \\ &= z_t \left[\Lambda_{W_t - W_{t_-}} (dW_t) - dW_t \right] g_t^{-1} \\ &+ \frac{1}{2} \sum_{A \in \beta} z_t \left[\frac{1}{2} \left\{ \Gamma_{W_t - W_{t_-}} (A, A) + A^2 \right\} - \left[\Lambda_{W_t - W_{t_-}} (A) \right] A \right] g_t^{-1} dt. \end{split}$$

From this equation and Burkholder's inequality,

$$\mathbb{E} \sup_{0 \le t \le \tau} |Q_t - I|^p \le C \mathbb{E} \left[\int_0^{\tau} \sum_A \left| z_t \left[\Lambda_{W_t - W_{t_-}}(A) - A \right] g_t^{-1} \right|^2 dt \right]^{p/2} \\ + C \mathbb{E} \left[\int_0^{\tau} \left| \sum_A z_t \left[\frac{1}{2} \left(\Gamma_{W_t - W_{t_-}}(A, A) + A^2 \right) - \left[\Lambda_{W_t - W_{t_-}}(A) \right] A \right] g_t^{-1} \right| dt \right]^p \\ \le C \tau^{p/2} \mathbb{E} \int_0^{\tau} \sum_A |z_t|^p \left| g_t^{-1} \right|^p \left| \Lambda_{W_t - W_{t_-}}(A) - A^2 \right|^p \frac{dt}{\tau}$$
(4.90)

(A.23)

$$+ C\tau^{p} \mathbb{E} \int_{0}^{\tau} \sum_{A} |z_{t}|^{p} |g_{t}^{-1}|^{p} \left| \frac{1}{2} \left(\Gamma_{W_{t}-W_{t_{-}}}(A,A) + A^{2} \right) - \left[\Lambda_{W_{t}-W_{t_{-}}}(A) \right] A \right|^{p} \frac{dt}{\tau}$$

From our estimates in Eq. (A.15) and (A.16),

$$\left| \Lambda_{W_t - W_{t_-}}(A) - A \right| \le C \left| W_t - W_{t_-} \right| e^{C \left| W_t - W_{t_-} \right|}$$
 and

$$\begin{aligned} \left| \frac{1}{2} \left(\Gamma_{W_t - W_{t_-}}(A, A) + A^2 \right) - \left[\Lambda_{W_t - W_{t_-}}(A) \right] A \right| \\ & \leq C \left(\left| W_t - W_{t_-} \right| + \left| W_t - W_{t_-} \right|^2 \right) e^{C |W_t - W_{t_-}|} \end{aligned}$$

which combined with Eq. (A.23) leads to

$$\mathbb{E} \sup_{0 \le t \le \tau} |Q_t - I|^p \le C\tau^{p/2-1} \mathbb{E} \int_0^\tau |z_t|^p |g_t^{-1}|^p |W_t - W_{t_-}|^p e^{pC|W_t - W_{t_-}|} dt + C\tau^{p-1} \mathbb{E} \int_0^\tau |z_t|^p |g_t^{-1}|^p |C \left(|W_t - W_{t_-}| + |W_t - W_{t_-}|^2 \right) e^{C|W_t - W_{t_-}|} \Big|_p^p dt \le CT^{p/2-1} \mathbb{E} \int_0^T |z_t|^p |g_t^{-1}|^p |W_t - W_{t_-}|^p e^{pC|W_t - W_{t_-}|} dt$$

(A.24)

$$+ CT^{p-1}\mathbb{E}\int_{0}^{T} |z_{t}|^{p} |g_{t}^{-1}|^{p} \left(|W_{t} - W_{t_{-}}| + |W_{t} - W_{t_{-}}|^{2} \right)^{p} e^{pC|W_{t} - W_{t_{-}}|} dt.$$

Let $\delta := t - t_{-} \leq \Delta_t$, then

$$\mathbb{E}\left[\left|W_{t}-W_{t-}\right|^{\rho}e^{C\left|W_{t}-W_{t-}\right|}\right] = \mathbb{E}\left[\left|\sqrt{\delta}W_{1}\right|^{\rho}e^{C\left|\sqrt{\delta}W_{1}\right|}\right]$$
$$= \delta^{\rho/2}\mathbb{E}\left[|W_{1}|^{\rho}e^{C\sqrt{\delta}|W_{1}|}\right] \le C\Delta_{t}^{\rho/2}.$$

So using this estimate, Propositions A.6 and A.7 along with Hölder's inequality implies there exists a constant C such that

$$\mathbb{E}[|z_t|^p |g_t^{-1}|^p |W_t - W_{t_-}|^p e^{pC|W_t - W_{t_-}|}]$$

$$\leq \left(\mathbb{E} |z_t|^{4p}\right)^{1/4} \left(\mathbb{E} |g_t^{-1}|^{4p}\right)^{1/4} \left(\mathbb{E} \left[|W_t - W_{t_-}|^{2p} e^{2pC|W_t - W_{t_-}|}\right]\right)^{1/2}$$

$$\leq C\Delta_t^{p/2} \leq C |\pi|^{p/2} .$$

Using this estimate in Eq. (A.24) shows there exists a constant $C < \infty$ such that

$$\sup_{t \in [0,T]} \mathbb{E} |Q_t - I|^p \le \mathbb{E} \sup_{t \in [0,T]} |Q_t - I|^p \le C |\pi|^{p/2}.$$

One more application of Proposition A.7 and Hölder's inequality shows

$$\mathbb{E} \sup_{t \in [0,T]} |z_t - g_t|^p = \mathbb{E} \sup_{t \in [0,T]} \left| \left(z_t g_t^{-1} - I \right) g_t \right|^p \\ \leq \left(\mathbb{E} \sup_{t \in [0,T]} |Q_t - I|^{2p} \mathbb{E} \sup_{t \in [0,T]} |g_t|^{2p} \right)^{1/2} \leq C |\pi|^{p/2}$$

which proves Eq. (A.2).

For $t \in J_n$,

$$\begin{aligned} |y_t - g_t| &\leq |y_t - y_{t_n}| + |y_{t_n} - g_{t_n}| + |g_t - g_{t_n}| \\ &\leq \sup_{\tau \in J_n} |y_\tau - y_{\tau_n}| + \sup_{\tau \in J_n} |g_\tau - g_{\tau_n}| + \sup_{\tau \in \pi} |y_\tau - g_\tau| \end{aligned}$$

and therefore,

$$\mathbb{E} \sup_{t \leq T} |y_t - g_t|^p$$

$$\leq 3^{1/q} \mathbb{E} \left[\max_n \sup_{\tau \in J_n} |y_\tau - y_{\tau_n}|^p + \max_n \sup_{\tau \in J_n} |g_\tau - g_{\tau_n}|^p + \sup_{\tau \in \pi} |y_\tau - g_\tau|^p \right]$$

$$\leq 3^{1/q} \mathbb{E} \left[\sum_{n=0}^{N-1} \sup_{\tau \in J_n} |y_\tau - y_{\tau_n}|^p + \sum_{n=0}^{N-1} \sup_{\tau \in J_n} |g_\tau - g_{\tau_n}|^p + \sup_{\tau \in \pi} |y_\tau - g_\tau|^p \right]$$

$$\leq C_p \left(\sum_{n=0}^{N-1} |t_{n+1} - t_n|^{p/2} + |\pi|^{p/2} \right) \leq C_p \left(|\pi|^{(p/2)-1} + |\pi|^{p/2} \right).$$

This proves (A.3). \blacksquare

References

- S. Aida, On the irreducibility of certain dirichlet forms on loop spaces over compact homogeneous spaces, New Trends in Stochastic Analysis (New Jersey) (K. D. Elworthy, S. Kusuoka, and I. Shigekawa, eds.), Proceedings of the 1994 Taniguchi Symposium, World Scientific, 1997, pp. 3–42.
- Shigeki Aida, Logarithmic Sobolev inequalities on loop spaces over compact Riemannian manifolds, Stochastic analysis and applications (Powys, 1995), World Sci. Publishing, River Edge, NJ, 1996, pp. 1–19.
- _____, Differential calculus on path and loop spaces. II. irreducibility of the dirichlet forms on loop spaces, preprint (1997), 1–17.
- Shigeki Aida and Bruce K. Driver, Equivalence of heat kernel measure and pinned Wiener measure on loop groups, C. R. Acad. Sci. Paris Sér. I Math. 331 (2000), no. 9, 709–712. MR 1 797 756
- Shigeki Aida and David Elworthy, Differential calculus on path and loop spaces. I. Logarithmic Sobolev inequalities on path spaces, C. R. Acad. Sci. Paris Sér. I Math. **321** (1995), no. 1, 97–102.
- H. Airault and P. Malliavin, Integration on loop groups. II. Heat equation for the Wiener measure, J. Funct. Anal. 104 (1992), no. 1, 71–109.
- Sergio Albeverio, Brian C. Hall, and Ambar N. Sengupta, *The Segal-Bargmann transform for two-dimensional Euclidean quantum Yang-Mills*, Infin. Dimens. Anal. Quantum Probab. Relat. Top. 2 (1999), no. 1, 27–49. MR 1 805 834
- Sergio Albeverio and Raphael Høegh-Krohn, The energy representation of Sobolev-Lie groups, Compositio Math. 36 (1978), no. 1, 37–51.
- Lars Andersson and Bruce K. Driver, Finite-dimensional approximations to Wiener measure and path integral formulas on manifolds, J. Funct. Anal. 165 (1999), no. 2, 430–498. MR 2000j:58059
- Jean-Michel Bismut, Mecanique aleatoire. (french) [random mechanics], Springer-Verlag, Berlin-New York, 1981, Lecture Notes in Mathematics, Vol. 866.
- Theodor Bröcker and Tammo tom Dieck, Representations of compact Lie groups, Springer-Verlag, New York, 1995, Translated from the German manuscript, Corrected reprint of the 1985 translation. MR 97i:22005
- Mireille Capitaine, Elton P. Hsu, and Michel Ledoux, Martingale representation and a simple proof of logarithmic Sobolev inequalities on path spaces, Electron. Comm. Probab. 2 (1997), 71–81 (electronic).
- Trevor R. Carson, Logarithmic sobolev inequalities for free loop groups, University of California at San Diego Ph.D. thesis. This may be retrieved at http://math.ucsd.edu/ driver/driver/thesis.htm, 1997.
- _____, A logarithmic Sobolev inequality for the free loop group, C. R. Acad. Sci. Paris Sér. I Math. 326 (1998), no. 2, 223–228. MR 99g:60108
- Bruce K. Driver, A Cameron-Martin type quasi-invariance theorem for Brownian motion on a compact Riemannian manifold, J. Funct. Anal. 110 (1992), no. 2, 272–376.

BRUCE K. DRIVER

- 16. _____, A Cameron-Martin type quasi-invariance theorem for pinned Brownian motion on a compact Riemannian manifold, Trans. Amer. Math. Soc. **342** (1994), no. 1, 375–395.
- Integration by parts and quasi-invariance for heat kernel measures on loop groups, J. Funct. Anal. 149 (1997), no. 2, 470–547.
- Integration by parts for heat kernel measures revisited, J. Math. Pures Appl. (9) 76 (1997), no. 8, 703–737.
- Bruce K. Driver and Brian C. Hall, Yang-Mills theory and the Segal-Bargmann transform, Comm. Math. Phys. 201 (1999), no. 2, 249–290. MR 2000c:58064
- _____, The energy representation has no non-zero fixed vectors, Stochastic processes, physics and geometry: new interplays, II (Leipzig, 1999), Amer. Math. Soc., Providence, RI, 2000, pp. 143–155. MR 1 803 410
- 21. Bruce K. Driver and Yaozhong Hu, Convergence rate estimate for wong-zakai approximations of stochastic differential equations, (2002).
- Bruce K. Driver and Terry Lohrenz, Logarithmic Sobolev inequalities for pinned loop groups, J. Funct. Anal. 140 (1996), no. 2, 381–448.
- Bruce K. Driver and Vikram K. Srimurthy, Absolute continuity of heat kernel measure with pinned wiener measure on loop groups, Ann. Probab. 29 (2001), no. 2, 691–723.
- Andreas Eberle, Spectral gaps on loop spaces: a counterexample, C. R. Acad. Sci. Paris Sér. I Math. 330 (2000), no. 3, 237–242. MR 2000k:60111
- O. Enchev and D. W. Stroock, Towards a Riemannian geometry on the path space over a Riemannian manifold, J. Funct. Anal. 134 (1995), no. 2, 392–416.
- Ognian Enchev and Daniel W. Stroock, Pinned Brownian motion and its perturbations, Adv. Math. 119 (1996), no. 2, 127–154.
- 27. Shi Zan Fang, Inégalité du type de Poincaré sur l'espace des chemins riemanniens, C. R. Acad. Sci. Paris Sér. I Math. 318 (1994), no. 3, 257–260.
- Shizan Fang, Integration by parts for heat measures over loop groups, J. Math. Pures Appl. (9) 78 (1999), no. 9, 877–894. MR 1 725 745
- Integration by parts formula and logarithmic Sobolev inequality on the path space over loop groups, Ann. Probab. 27 (1999), no. 2, 664–683. MR 1 698 951
- Shizan Fang and Jacques Franchi, A differentiable isomorphism between Wiener space and path group, Séminaire de Probabilités, XXXI, Lecture Notes in Math., vol. 1655, Springer, Berlin, 1997, pp. 54–61.
- 31. Ezra Getzler, Dirichlet forms on loop space, Bull. Sci. Math. (2) 113 (1989), no. 2, 151–174.
- An extension of Gross's log-Sobolev inequality for the loop space of a compact Lie group, Probability models in mathematical physics (Colorado Springs, CO, 1990), World Sci. Publishing, Teaneck, NJ, 1991, pp. 73–97.
- Fu-Zhou Gong and Zhi-Ming Ma, The log-Sobolev inequality on loop space over a compact Riemannian manifold, J. Funct. Anal. 157 (1998), no. 2, 599–623. MR 99f:58222
- Fuzhou Gong, Michael Röckner, and Liming Wu, Poincaré inequality for weighted first order Sobolev spaces on loop spaces, J. Funct. Anal. 185 (2001), no. 2, 527–563. MR 1 856 276
- Leonard Gross, Logarithmic Sobolev inequalities, Amer. J. Math. 97 (1975), no. 4, 1061–1083. MR 54 8263
- 36. _____, Logarithmic Sobolev inequalities on loop groups, J. Funct. Anal. 102 (1991), no. 2, 268–313.
- 37. _____, Uniqueness of ground states for Schrödinger operators over loop groups, J. Funct. Anal. 112 (1993), no. 2, 373–441.
- Leonard Gross and Paul Malliavin, Hall's transform and the Segal-Bargmann map, Itô's stochastic calculus and probability theory, Springer, Tokyo, 1996, pp. 73–116.
- Brian C. Hall, The Segal-Bargmann "coherent state" transform for compact Lie groups, J. Funct. Anal. 122 (1994), no. 1, 103–151. MR 95e:22020
- 40. _____, The Segal-Bargmann transform and the Gross ergodicity theorem, To appear in Conference Proceeding in Honor of L. Gross (2002).
- Richard Holley and Daniel Stroock, Logarithmic Sobolev inequalities and stochastic Ising models, J. Statist. Phys. 46 (1987), no. 5-6, 1159–1194. MR 89e:82013
- Elton P. Hsu, Inégalités de Sobolev logarithmiques sur un espace de chemins, C. R. Acad. Sci. Paris Sér. I Math. **320** (1995), no. 8, 1009–1012.
- 43. _____, Quasi-invariance of the Wiener measure on the path space over a compact Riemannian manifold, J. Funct. Anal. **134** (1995), no. 2, 417–450.

- 44. Yaozhong Hu, Strong and weak order of time discretization schemes of stochastic differential equations, Séminaire de Probabilités, XXX, Springer, Berlin, 1996, pp. 218–227. MR 98j:60081
- Nobuyuki Ikeda and Shinzo Watanabe, Stochastic differential equations and diffusion processes, second ed., North-Holland Publishing Co., Amsterdam, 1989. MR 90m:60069
- Peter E. Kloeden and Eckhard Platen, Numerical solution of stochastic differential equations, Springer-Verlag, Berlin, 1992. MR 94b:60069
- Hui Hsiung Kuo, Gaussian measures in Banach spaces, Springer-Verlag, Berlin, 1975, Lecture Notes in Mathematics, Vol. 463.
- R. Léandre, Cohomologie de Bismut-Nualart-Pardoux et cohomologie de Hochschild entière, Séminaire de Probabilités, XXX, Lecture Notes in Math., vol. 1626, Springer, Berlin, 1996, pp. 68–99.
- Terry Lyons and Zhongmin Qian, Flow equations on spaces of rough paths, J. Funct. Anal. 149 (1997), no. 1, 135–159.
- Marie-Paule Malliavin and Paul Malliavin, Integration on loop groups. I. Quasi invariant measures, J. Funct. Anal. 93 (1990), no. 1, 207–237.
- Paul Malliavin, Hypoellipticity in infinite dimensions, Diffusion processes and related problems in analysis, Vol. I (Evanston, IL, 1989), Progr. Probab., vol. 22, Birkhäuser Boston, Boston, MA, 1990, pp. 17–31.
- 52. J. R. Norris, Twisted sheets, J. Funct. Anal. 132 (1995), no. 2, 273-334.
- Gaku Sadasue, Equivalence-singularity dichotomy for the Wiener measures on path groups and loop groups, J. Math. Kyoto Univ. 35 (1995), no. 4, 653–662. MR 98d:60108
- Ichiro Shigekawa, Transformations of the Brownian motion on a Riemannian symmetric space, Z. Wahrsch. Verw. Gebiete 65 (1984), no. 4, 493–522.
- 55. _____, Transformations of the Brownian motion on the Lie group, Stochastic analysis (Katata/Kyoto, 1982), North-Holland Math. Library, vol. 32, North-Holland, Amsterdam, 1984, pp. 409–422.
- _____, Differential calculus on a based loop group, New trends in stochastic analysis (Charingworth, 1994), World Sci. Publishing, River Edge, NJ, 1997, pp. 375–398. MR 99k:60146
- Vikram K. Srimurthy, On the equivalence of measures on loop space, Probab. Theory Related Fields 118 (2000), no. 4, 522–546. MR 2001m:58022
- Daniel W. Stroock, Lectures on stochastic analysis: diffusion theory, Cambridge University Press, Cambridge, 1987. MR 88d:60203
- Eugene Wong and Moshe Zakai, On the relation between ordinary and stochastic differential equations, Internat. J. Engrg. Sci. 3 (1965), 213–229. MR 32 #505