

Article

Simultaneous Momentum and Position Measurement and the Instrumental Weyl-Heisenberg Group

Christopher S. Jackson ¹ and Carlton M. Caves ^{2,*}

¹ Independent Researcher, Gold Beach, OR 97444, USA; omgphysics@gmail.com

² Center for Quantum Information and Control, University of New Mexico, Albuquerque, NM 87131, USA

* Correspondence: ccaves@unm.edu

Abstract: The canonical commutation relation, $[Q, P] = i\hbar$, stands at the foundation of quantum theory and the original Hilbert space. The interpretation of P and Q as observables has always relied on the analogies that exist between the unitary transformations of Hilbert space and the canonical (also known as contact) transformations of classical phase space. Now that the theory of quantum measurement is essentially complete (this took a while), it is possible to revisit the canonical commutation relation in a way that sets the foundation of quantum theory not on unitary transformations but on positive transformations. This paper shows how the concept of simultaneous measurement leads to a fundamental differential geometric problem whose solution shows us the following. The simultaneous P and Q measurement (SPQM) defines a universal measuring instrument, which takes the shape of a seven-dimensional manifold, a universal covering group we call the instrumental Weyl-Heisenberg (IWH) group. The group IWH connects the identity to classical phase space in unexpected ways that are significant enough that the positive-operator-valued measure (POVM) offers a complete alternative to energy quantization. Five of the dimensions define processes that can be easily recognized and understood. The other two dimensions, the normalization and phase in the center of the IWH group, are less familiar. The normalization, in particular, requires special handling in order to describe and understand the SPQM instrument.

Keywords: Weyl-Heisenberg group; measuring instrument; Kraus operator; Cartan decomposition; Harish-Chandra decomposition; right-invariant derivative; Maurer-Cartan form; Wiener path integral; diffusion equation; stochastic differential equation



Citation: Jackson, C.S.; Caves, C.M. Simultaneous Momentum and Position Measurement and the Instrumental Weyl-Heisenberg Group. *Entropy* **2023**, *25*, 1221. <https://doi.org/10.3390/e25081221>

Academic Editors: Andrei Khrennikov and Karl Svozil

Received: 1 June 2023

Revised: 26 July 2023

Accepted: 4 August 2023

Published: 16 August 2023



Copyright: © 2023 by the authors. Licensee MDPI, Basel, Switzerland. This article is an open access article distributed under the terms and conditions of the Creative Commons Attribution (CC BY) license (<https://creativecommons.org/licenses/by/4.0/>).

1. Introduction

After World War II, theoretical quantum physics became dominated by the design of quantum field theory. There were three branches of physics that stemmed from this: high energy, condensed matter, and atomic-molecular-optical (AMO) physics. Although incredibly developed as predictive methods, quantum field theory in all three of these branches has left some very basic ideas of quantum observation underdeveloped. That this is indeed still the case is evident in how the coherent-state resolutions of the identity or “overcomplete bases” [1–4] are usually finessed:

$$1_Z^{(\text{amp})} = Z \int_{\mathbb{C}} \frac{d^2\alpha}{\pi} |\sqrt{Z}\alpha\rangle\langle\sqrt{Z}\alpha| \quad (1)$$

for bosonic amplitudes (where Z is an arbitrary complex scalar) [1,5–8] and

$$1_j^{(\text{spin})} = (2j+1) \int_{S^2} \frac{d\mu(\hat{n})}{4\pi} |j, \hat{n}\rangle\langle j, \hat{n}| \quad (2)$$

for fermionic spins (where $2j$ is an integer) [2–4,9]. These overcomplete bases are key both to establishing functional integration [1,10–19] and to understanding the observation

of energy quanta [20]. Yet these overcomplete bases do not fit the most basic idea of a Hermitian eigenbasis and as such are often not considered a serious form of observation. By itself, this inattention to how one could observe in these overcomplete bases leaves a missing piece at the foundation of quantum mechanics and quantum field theory.

Meanwhile, an understanding of the overcomplete bases as a bona fide means of observation has been slowly coming to light. First, the idea of observables matured into the general mathematical theories of instruments and operations [21–32]. In this more contemporary language, the overcomplete bases of Equations (1) and (2) are called *positive-operator-valued measures* (POVMs) and are understood to be informationally complete, meaning the distribution observed with any one of these POVMs is enough to reconstruct the quantum state. With this theoretical technology, it has further been discovered that the overcomplete bases correspond to various forms of continual (or continuous) simultaneous observation. The standard coherent-state POVM given in Equation (1) was discovered to be the effect of simultaneously observing both quadratures of a leaky cavity, a form of observation we will call Goetsch-Graham-Wiseman (GGW) heterodyne detection [33–37]. Then, the spin-coherent POVM of Equation (2) was discovered to be the effect of simultaneously observing the three orthogonal spin components, a form of observation we call the *spin-isotropic measurement* [38–40].

Before proceeding, we caution that this paper uses a mathematical apparatus not familiar to most physicists and quantum scientists. This apparatus is introduced here naturally as it becomes both desirable and necessary. Readers who are made uncomfortable by this apparatus are urged to consult the companion paper [40], which attempts to persuade the reader that the unfamiliar mathematical concepts and techniques are essential tools—a new way of thinking and doing—and then introduces these tools as gently as possible.

GGW heterodyne detection and the spin-isotropic measurement work in a very similar way, but they are different in one very important respect. While GGW heterodyne detection assumes energy-conserving system-meter interactions, the spin-isotropic measurement assumes Hermitian meter displacements, $-iH^{(\text{iso})}dt/\hbar = -\sqrt{\kappa dt} J_k \otimes 2\sigma\partial_q$, where J_k is an orthogonal spin component of the system, q is the meter register, σ is the width of the meter pointer, and κ is the measurement rate. In both cases, the measuring instrument consists of Kraus operators defined by a time-ordered exponential over the duration T of the measurement. For GGW heterodyne detection, the Kraus operators are [37]

$$L^{(\text{GGW})}[dw_{[0,T]}] = \mathcal{T} \exp\left(\int_0^{T-dt} -2a^\dagger a \kappa dt + 2a\sqrt{\kappa} dw_t^*\right), \quad (3)$$

where $a = (Q + iP)/\sqrt{2\hbar}$ is the usual complex-amplitude operator and $dw_t = (dW_t^q + i dW_t^p)/\sqrt{2}$ is the registered complex Wiener path. For spin-isotropic measurement, the Kraus operators are [39,40]

$$L^{(\text{iso})}[d\vec{W}_{[0,T]}] = \mathcal{T} \exp\left(\int_0^{T-dt} -\vec{J}^2 \kappa dt + \vec{J} \cdot \sqrt{\kappa} d\vec{W}_t\right), \quad (4)$$

where $\vec{J} = (J_x, J_y, J_z)$ is the triple of orthogonal spin-component observables and $d\vec{W}_t = (dW_t^x, dW_t^y, dW_t^z)$ is the registered three-vector of Wiener paths. The most striking feature about the instruments defined by Equations (3) and (4) is that they can be integrated *universally*, that is, independently of matrix representation. The difference between the two cases can now be summarized as the following: integrating Equation (4) defines a seven-dimensional manifold that requires the theory of symmetric spaces, whereas integrating Equation (3) defines a three-dimensional manifold that is much more straightforward.

This paper is an analysis of the quadrature analog of the spin-isotropic measurement, a form of observation we call the *simultaneous P and Q measurement* (SPQM). The name SPQM is our homage to Alberto Barchielli, who appears to be the first to have considered and analyzed this problem [41] (*Senatus PopulusQue Romanus* (SPQR), translated as “The

Senate and People of Rome”, is an enduring symbol of ancient Rome [42]. SPQM has been similarly associated with the city of Milan, so we offer it as a tribute to Barchielli and his lifetime of work on continuous measurement theory at the University of Milan). SPQM generates a measuring instrument with Kraus operators [40]

$$L^{(\text{SPQM})}[dw_{[0,T)}] = \mathcal{T} \exp\left(\int_0^{T-dt} -2H_o \kappa dt + P\sqrt{\kappa} dW_t^p + Q\sqrt{\kappa} dW_t^q\right), \quad (5)$$

where $2H_o \equiv P^2 + Q^2$ and P and Q are (dimensionless) canonical momentum and position (or the conserved quadrature components of a harmonic oscillator). This time-ordered exponential defines another fundamental seven-dimensional manifold, the universal covering group, which for SPQM we call the *instrumental Weyl-Heisenberg group*, $G = \text{IWH}$. The universal covering group is defined by a map γ with the universal property that for any Hilbert space \mathcal{H} carrying the paths of Kraus operators $L^{(\text{SPQM})} : \mathbb{C}^{T/dt} \rightarrow \text{GL}(\mathcal{H})$, there exists a unique representation R such that

$$L^{(\text{SPQM})} = R \circ \gamma \quad \text{where} \quad \mathbb{C}^{T/dt} \xrightarrow{\gamma} \text{IWH} \xrightarrow{R} \text{GL}(\mathcal{H}). \quad (6)$$

This universal way of considering $G = \text{IWH}$ essentially amounts to suspending the choice of \mathcal{H} and, therefore, \hbar , but it is very important to appreciate that the measuring instrument is, in fact, fundamentally independent of the Hilbert space and, therefore, \hbar . The same is true for the spin-isotropic measurement, except that different irreducible representations do not amount to choices of \hbar , but rather to choices of the total angular momentum number j .

The universal covering group IWH can be navigated in much the same way as can be done for semisimple Lie groups, with the use of right-invariant vector fields and decompositions similar to those of Cartan and Harish-Chandra. In particular, the sample paths defined by $x(t) = \gamma[dw_{[0,t)}]$ diffuse according to a Fokker-Planck-Kolmogorov equation,

$$\frac{1}{\kappa} \frac{\partial}{\partial t} D_t(x) = \left(2H_o + \frac{1}{2} \overleftarrow{P} \overleftarrow{P} + \frac{1}{2} \overleftarrow{Q} \overleftarrow{Q} \right) [D_t](x), \quad (7)$$

where

$$D_T(x) \equiv \int \mathcal{D}\mu[dw_{[0,T)}] \delta(x, \gamma[dw_{[0,T)}]) \quad (8)$$

is the *Kraus-operator distribution function* of the SPQM instrument with respect to the Haar measure [18,43–45] of the IWH group and $\overleftarrow{H_o}$, \overleftarrow{Q} , and \overleftarrow{P} are right-invariant derivatives tangent to the IWH group. We regard “Kraus-operator distribution function”, “Kraus-operator distribution”, and “Kraus-operator density” as interchangeable, despite subtle differences some might attribute to these usages. We abbreviate Kraus-operator distribution function as KOD to invite the reader to use any of these terms. The KOD can be considered the *universal unraveling* of the total (or unconditional) operation (a completely positive, trace-preserving superoperator),

$$\mathcal{Z}_T^{(\text{SPQM})} \equiv \int \mathcal{D}\mu[dw_{[0,T)}] \mathcal{O}\left(L^{(\text{SPQM})}[dw_{[0,T)}]\right) = \int_{\text{IWH}} d^7\mu(x) D_T(x) \mathcal{O}(R(x)), \quad (9)$$

where $\mathcal{D}\mu[dw_{[0,T)}]$ is the Wiener path measure, $d^7\mu(x)$ is the Haar measure, and $\mathcal{O}(L) \equiv L \odot L^\dagger$. The technology of right-invariant differentiation [46–48] will not be familiar to most quantum physicists and information scientists. Introducing SPQM, the IWH group, the concept of right-invariant motion, and the KOD is the subject of Section 2.

Section 3 translates the coordinate-independent formulation from Section 2 to forms that physicists are more likely to recognize. Indeed, what the aforementioned decompositions do is to coordinate the points of the IWH group [2,3,49–52]. The decomposition of the IWH group, similar to Harish-Chandra [52–54] (also known as “Gauss” [2,51,55]) decomposition, is given by

$$x = e^{a^\dagger \nu} e^{-H_0 r + \Omega z} e^{a \mu^*}, \quad (10)$$

with the purity coordinate $r \in \mathbb{R}$, which we will call the ruler; central coordinate $z = -s + i\psi \in \mathbb{C}$; phase-space coordinates $\nu, \mu \in \mathbb{C}$; and $\Omega = \hbar 1_{\mathcal{H}}$. The decomposition of the IWH group, similar to the Cartan decomposition [50–52,54], is given by

$$x = D_\beta e^{i\Omega\phi} e^{-H_0 r - \Omega\ell} D_\alpha^{-1}, \quad (11)$$

with different central coordinates $\phi, \ell \in \mathbb{R}$, the same ruler r , and phase-space coordinates $\beta, \alpha \in \mathbb{C}$ appearing in the conventional displacement operator $D_\alpha \equiv e^{a^\dagger \alpha - a \alpha^*}$. Introducing these decompositions and using them to transform Equation (5) into standard Itô-form stochastic differential equations (SDEs) [56–60] and to transform Equation (7) into a coordinate Fokker-Planck-Kolmogorov (FPK) diffusion equation [59,60] is the subject of Section 3, which also solves those SDEs and (mostly) solves the FPK diffusion equation.

As a function of the registers of SPQM, we find that the ruler satisfies

$$r_T = 2\kappa T. \quad (12)$$

For the remaining Harish-Chandra coordinates, we find that the phase points follow Ornstein-Uhlenbeck [59–61] and GGW processes,

$$\nu[dw_{[0,T)}] = \int_0^{T_-} \sqrt{\kappa} dw_t e^{-2\kappa(T-t)} \quad \text{and} \quad \mu[dw_{[0,T)}] = \int_0^{T_-} \sqrt{\kappa} dw_t e^{-2\kappa t}, \quad (13)$$

where $T_- \equiv T - dt$, and the center follows a quadratic functional process,

$$z[dw_{[0,T)}] = \frac{1}{2} \int_0^{T_-} \int_0^{T_-} \kappa dw_t^* dw_s e^{-2\kappa|t-s|} (1 + \text{sgn}(t-s)), \quad (14)$$

where the sign function is used: $\text{sgn}(u) = u/|u|$ for $u \neq 0$ and $\text{sgn}(0) = 0$. The Cartan phase-space coordinates follow linear functionals,

$$\beta[dw_{[0,T)}] = \int_0^{T_-} \sqrt{\kappa} dw_t \frac{\cosh 2\kappa t}{\sinh 2\kappa T} \quad \text{and} \quad \alpha[dw_{[0,T)}] = \int_0^{T_-} \sqrt{\kappa} dw_t \frac{\cosh 2\kappa(T-t)}{\sinh 2\kappa T}. \quad (15)$$

The Cartan central coordinates ℓ and ϕ follow from the coordinate transformation between Harish-Chandra and Cartan coordinates, which is given in Equations (A34), (A35), (A38) and (A39), but we do not write those solutions explicitly here.

As for the KOD $D_t(x)$, we will not be able to solve analytically for the distribution over all seven dimensions. Summing over the center

$$Z \equiv \{e^{1z} : z \in \mathbb{C}\} \triangleleft \text{IWH}, \quad (16)$$

however, gives a *reduced SPQM unraveling* of the total operation,

$$\mathcal{Z}_T^{(\text{SPQM})} = \int_{\text{IWH}/Z} d^5\mu(Zx) C_T(Zx) \mathcal{O}(D_\beta e^{-H_0 r} D_\alpha^\dagger), \quad (17)$$

where the integral is over the quotient group IWH/Z , which we call the *reduced instrumental Weyl-Heisenberg* (RIWH) group. The reduced Kraus-operator distribution (RKOD),

$$C_T(Zx) \equiv \int_Z d\phi d\ell D_T(x) e^{-2\ell} = \int \mathcal{D}\mu[dw_{[0,T]}] e^{-2\ell[dw_{[0,T]}]} \delta(Zx, Z\gamma[dw_{[0,T]}]), \quad (18)$$

is a marginal over the center that includes the Cartan center factor $e^{-2\ell}$. We call $C_T(Zx)$ the *Cartan-section reduced distribution function*, and we are able to solve for it from its FPK diffusion equation. The solution is a Gaussian with ill-defined normalization,

$$d^5\mu(Zx) C_T(Zx) = 2 \sinh r dr \delta(r - 2\kappa T) \frac{d^2\beta}{\pi} \frac{1}{\Sigma_T} e^{-|\beta-\alpha|^2/\Sigma_T} \frac{d^2\alpha}{\pi}, \quad (19)$$

where the mean-square distance between the two phase points is given by

$$\Sigma_T = \kappa T - \tanh \kappa T. \quad (20)$$

The normalization factor $2 \sinh 2\kappa T$ is particularly interesting, as the POVM completeness relation for the SPQM instrument boils down to (assuming $\hbar = 1$)

$$2 \sinh 2\kappa T \int \frac{d^2\alpha}{\pi} D_\alpha e^{-H_0 4\kappa T} D_\alpha^\dagger = 1_{\mathcal{H}}, \quad (21)$$

and this can be recognized as equivalent to the result of energy quantization,

$$\text{tr} e^{-H_0 4\kappa T} = \sum_{n=0}^{\infty} e^{-(n+\frac{1}{2})4\kappa T} = \frac{1}{2 \sinh 2\kappa T}. \quad (22)$$

This demonstrates that the KOD can be considered an alternative to energy quantization. Indeed, the operator H_0 does not appear in SPQM as an energy observable, but rather as a term required by the positivity of sampling measurement records. The late-time limit of the completeness relation is of particular interest: when $T \gg 1/\kappa$, $D_\alpha e^{-H_0 4\kappa T} D_\alpha^\dagger = e^{-2\kappa T} |\alpha\rangle\langle\alpha|$, showing that the SPQM POVM elements approach Glauber coherent states; the completeness relation shows that it does so uniformly,

$$1_{\mathcal{H}} = \int \frac{d^2\alpha}{\pi} |\alpha\rangle\langle\alpha|, \quad (23)$$

thus giving the coherent-state resolution of the identity.

Everything in our analysis—the path integrals, the KODs, the FPK diffusion equations, the SDEs—follows from the path integral for the overall quantum operation $\mathcal{Z}_T^{(\text{SPQM})}$, given in Equation (9), which integrates over the sample paths as they appear in the time-ordered Kraus operator $L^{(\text{SPQM})}[dw_{[0,T]}]$. As the analysis develops in Sections 2 and 3, however, there emerges a disconnect between the reduced distribution $C_T(Zx)$, which expresses the completeness relation, and the SDEs and stochastic integrals for the phase-point coordinates: $C_t(Zx)$ has ill-defined normalization, and it is not the weight function whose moments are those of the Cartan phase-point variables, β and α , as they are expressed in the stochastic-integral solutions of Equation (15). The disconnect is all about the real center term, $e^{-\Omega\ell}$, which scales the Kraus operators when they are written in Cartan coordinates.

The three faces of the KOD stochastic trinity—path integrals, diffusion equations, and SDEs—having been sundered in Section 3, are re-united in Section 4. The vehicle for the reunion is the *Harish-Chandra-section reduced distribution function*,

$$B_T(Zx) \equiv \int_Z d\psi ds D_T(x) e^{-2s} = \int \mathcal{D}\mu[dw_{[0,T]}] e^{-2s[dw_{[0,T]}]} \delta(Zx, Z\gamma[dw_{[0,T]}]), \quad (24)$$

where $-s$ is the real part of the Harish-Chandra center coordinate z . The distribution $B_T(Zx)$ can be determined in two equivalent ways: first, from the FPK diffusion equation for $B_T(Zx)$ and, second, from applying the stochastic integrals for the Harish-Chandra phase-plane variables, ν and μ , to the above path-integral expression for $B_T(Zx)$. The stochastic trinity thus restored, one returns to completeness through the relation

$$C_T(Zx) = e^{f(Zx)} B_T(Zx), \quad (25)$$

where $f(Zx)$ is a quadratic function of the phase-plane variables, given in Equations (A30) and (A38), which comes directly out of the coordinate transformation between Cartan and Harish-Chandra coordinates; $e^{2f(Zx)}$ can be regarded as a positive gauge transformation [19].

Section 5 concludes with musings on the stochastic trinity and the Lie-group manifolds that house instrument evolution.

2. The IWH Group and Coordinate-Free Right-Invariant Motion

This section defines the *simultaneous P and Q measurement* (SPQM) process and presents the instrumental Weyl-Heisenberg group, $G = \text{IWH}$, as the universal covering group of SPQM.

Section 2.1 introduces Kraus operators and the concept of observables generating positive transformations instead of unitary transformations. Section 2.2 introduces the SPQM process and the group IWH. Section 2.3 explains how the SPQM instrument is universal and defines a *Kraus-operator distribution (or density)* (KOD) over the IWH group. Section 2.4 explains how the KOD diffuses over time with the introduction of *right-invariant derivatives*, a differential-geometric technology that will be unfamiliar to most physicists and quantum scientists. Section 2.5 explains how the sample paths of the Kraus-operator diffusion are described with the introduction of the *right-invariant one-forms*, which are duals of right-invariant derivatives.

The diffusion equation and stochastic differential equations in Sections 2.4 and 2.5 are solved in Section 3.

2.1. Observables and Infinitesimal Positive Transformation

While observables are often considered to be infinitesimal generators of unitary transformations, they can also generate *positive transformations*. Let X be a Hermitian observable, κ be a real number with units of inverse time, and dW be a standard Wiener increment [10–12,15,16,19,59,60], which has the measure

$$d\mu(dW) = \frac{d(dW)}{\sqrt{2\pi dt}} e^{-dW^2/2dt}. \quad (26)$$

Unitary transformations can be infinitesimally generated either deterministically or stochastically, such as

$$e^{-iX\kappa dt} \quad \text{or} \quad \sqrt{d\mu(dW)} e^{-iX\sqrt{\kappa} dW}. \quad (27)$$

Positive transformations, on the other hand, are fundamentally stochastic. In addition, the infinitesimal generators of positive transformations are not of the canonical form, with a single parameter conjugate to the infinitesimal generator. The positive transformations we will be interested in do not involve jump operators, as in photodetection [37,62], but rather are differential, with infinitesimal generators of the form [40,63–65]

$$\sqrt{d\mu(dW)} L_X(dW) \equiv \sqrt{d\mu(dW)} e^{-X^2\kappa dt + X\sqrt{\kappa} dW}. \quad (28)$$

As a set, these positive transformations define a *measuring instrument*, complete over the Hilbert space, \mathcal{H} , according to the relation,

$$\int_{\mathbb{R}} d\mu(dW) L_X(dW)^\dagger L_X(dW) = 1_{\mathcal{H}}; \quad (29)$$

the operators $L_X(dW)$ are known as *Kraus operators* [29,31], and the set of elements, $\{d\mu(dW) L_X(dW)^\dagger L_X(dW)\}$, is known as a *positive-operator-valued measure* (POVM). We will call the positive transformations of Equation (28) *differential Kraus operators* and their set a *differential instrument*, both in anticipation of the (multi-dimensional) differential geometry coming up and to contrast these Kraus operators with jump operators.

The form of these infinitesimally generated positive transformations comes from the requirement that the total operation be completely positive and trace preserving. In particular, we can define the superoperator $A \odot B$ by

$$A \odot B(\rho) \equiv A\rho B \quad (30)$$

and the (selective, Kraus-rank-one) *operations* by

$$\mathcal{O}(L) \equiv L \odot L^\dagger. \quad (31)$$

Then, the aforementioned Kraus operators define a *total operation*,

$$\int d\mu(dW) \mathcal{O}\left(e^{-X^2 \kappa dt + X \sqrt{\kappa} dW}\right) = e^{-\frac{1}{2} \kappa dt \text{ad}_X^2}, \quad (32)$$

where the *adjoint* superoperator is defined by

$$\text{ad}_X \equiv X \odot 1 - 1 \odot X. \quad (33)$$

In this context, the infinitesimal generator or *Lindbladian*,

$$-\frac{1}{2} \text{ad}_X^2 = X \odot X - \frac{1}{2} (X^2 \odot 1 + 1 \odot X^2), \quad (34)$$

defines X as a *Lindblad operator* [28].

Kraus operators can be interpreted as an *indirect measurement* [29,31,40,66,67] where, in the differential case, a meter with an initial meter wavefunction

$$\sqrt{dq} \langle q|0\rangle = \sqrt{\frac{dq}{\sqrt{2\pi\sigma^2}}} e^{-q^2/2\sigma^2}, \quad (35)$$

is displaced by the system according to the interaction

$$-\frac{i}{\hbar} H dt = \sqrt{\kappa dt} X \otimes \left(-2\sigma \frac{\partial}{\partial q}\right), \quad (36)$$

which, in turn, registers some “position” q , fixing

$$dW = \sqrt{dt} \frac{q}{\sigma}, \quad (37)$$

so that

$$\sqrt{d\mu(dW)} L_X(dW) = \sqrt{dq} \langle q| e^{-iH dt/\hbar} |0\rangle \quad (38)$$

(for further details, see [40]).

An irresistible sidenote, developed more generally and in more detail in [40], is that the stochastic unitary transformations of Equation (27) follow from the same meter interaction, given in Equation (36), but with registration of the meter momentum p instead of its

position q . As such, the stochastic unitary transformations have a total operation identical to Equation (32),

$$\int d\mu(dW) \mathcal{O}(e^{-iX\sqrt{\kappa}dW}) = e^{-\frac{1}{2}\kappa dt \text{ad}_X^2}, \quad (39)$$

This alternative unraveling of the total operation corresponds to a symmetry of the general Lindbladian,

$$\mathcal{L}(A) = A \odot A^\dagger - \frac{1}{2}(A^\dagger A \odot 1 + 1 \odot A^\dagger A), \quad (40)$$

which is

$$\mathcal{L}(-iX) = \mathcal{L}(X), \quad (41)$$

so that $-iX$ is an alternative Lindblad operator of the total operation.

2.2. SPQM and the IWH Group

The subject of this paper is the continual (or continuous) simultaneous observation of the canonical observables P and Q defined by the canonical commutation relations,

$$[Q, P] = i\Omega \quad \text{and} \quad [\Omega, Q] = [\Omega, P] = 0. \quad (42)$$

As a Lie algebra of infinitesimal generators, these observables are usually considered to generate unitary displacement operators,

$$D_\alpha = e^{iQ\alpha_2 - iP\alpha_1}, \quad \text{where} \quad \alpha = \frac{\alpha_1 + i\alpha_2}{\sqrt{2}}, \quad (43)$$

which together define the three-dimensional *unitary Weyl-Heisenberg group*,

$$K = \text{UWH} \equiv \left\{ D_\alpha e^{i\Omega\phi} : \alpha \in \mathbb{C}, \phi \in \mathbb{R} \right\}. \quad (44)$$

If the generators operate irreducibly on the Hilbert space \mathcal{H} , then, by Shur's lemma, $\Omega = \hbar 1_{\mathcal{H}}$ for some $\hbar \in \mathbb{R}$ (because $g\Omega g^{-1} = \Omega$ for all $g \in \text{UWH}$.) Assuming that \hbar is finite—that is, $\hbar \neq 0$ —all such representations are essentially equivalent because the observables can always be rescaled so as to make the choice $\hbar = 1$. Therefore, it is usually assumed that

$$\Omega = 1_{\mathcal{H}}. \quad (45)$$

We shall now assume that $\hbar = 1$ but we will continue nonetheless to use Ω for the infinitesimal generator to emphasize that its existence is not defined by the Hilbert space but instead by the canonical commutation relations. In particular, what is defined by the Hilbert space is the relation $\Omega^2 = \Omega$, which is associative algebraic and not Lie algebraic.

For quantization, the unitary Weyl-Heisenberg group is supplemented with the unitary group generated by

$$H_o \equiv \frac{P^2 + Q^2}{2}, \quad (46)$$

defining the four-dimensional *dynamical Weyl-Heisenberg group*,

$$\text{DWH} \equiv \left\{ e^{-iH_o s} D_\alpha e^{i\Omega\phi} : s \in \mathbb{R}, \alpha \in \mathbb{C}, \phi \in \mathbb{R} \right\}. \quad (47)$$

Here, the use of “dynamical” refers to the analogies between H_0 and the classical Hamiltonian of a simple harmonic oscillator, upon which quantum mechanics was originally founded.

The observables P and Q can be measured simultaneously in the sense that the positive transformations they generate commute infinitesimally (to order dt),

$$L(dw) \equiv L_Q(dW^q) L_P(dW^p) \quad (48)$$

$$= e^{-Q^2 \kappa dt + Q \sqrt{\kappa} dW^q} e^{-P^2 \kappa dt + P \sqrt{\kappa} dW^p} \quad (49)$$

$$= e^{-(Q^2 + P^2) \kappa dt + Q \sqrt{\kappa} dW^q + P \sqrt{\kappa} dW^p + \frac{1}{2} [Q, P] \kappa dW^q dW^p} \quad (50)$$

$$= e^{-(Q^2 + P^2) \kappa dt + Q \sqrt{\kappa} dW^q + P \sqrt{\kappa} dW^p} \quad (51)$$

$$= L_P(dW^p) L_Q(dW^q), \quad (52)$$

so long as the Wiener outcome increments, dW^q and dW^p , are independent, their joint Wiener measure being

$$d\mu(dw) \equiv d\mu(dW^q) d\mu(dW^p) = \frac{d^2(dw)}{\pi dt} e^{-dw^* dw / dt}. \quad (53)$$

Here we switch to using a complex Wiener increment

$$dw \equiv \frac{dW^q + idW^p}{\sqrt{2}}. \quad (54)$$

Continually repeating this simultaneous measurement for a finite amount of time, T , defines the overall Kraus operators,

$$L[dw_{[0,T)}] = \mathcal{T} \prod_{k=0}^{T/dt-1} L(dw_{kdt}), \quad (55)$$

where $\mathcal{T} \prod$ denotes a time-ordered product. It is important to appreciate that, while the Kraus operators commute infinitesimally, they do not commute over finite amounts of time. Finally, these Kraus operators are accompanied by the Wiener path measure,

$$\mathcal{D}\mu[dw_{[0,T)}] = \prod_{k=0}^{T/dt-1} d\mu(dw_{kdt}) = \left(\prod_{k=0}^{T/dt-1} d^2(dw_{kdt}) \right) \left(\frac{1}{\pi dt} \right)^{T/dt} \exp \left(- \int_0^T \frac{|dw_t|^2}{dt} \right), \quad (56)$$

which is written here in terms of the complex Wiener increments. In summary, we have defined a time-dependent instrument with Kraus operators

$$\sqrt{\mathcal{D}\mu[dw_{[0,T)}]} L[dw_{[0,T)}] = \sqrt{\mathcal{D}\mu[dw_{[0,T)}]} \mathcal{T} \exp \left(\int_0^T -H_0 2\kappa dt + Q \sqrt{\kappa} dW_t^q + P \sqrt{\kappa} dW_t^p \right), \quad (57)$$

where $\mathcal{T} \exp \int$ is the time-ordered exponential. This is the *simultaneous P and Q measurement* (SPQM.) It is worth repeating that H_0 appears here due to the form of the differential positive transformations of Equation (28) and in this context is not a Hamiltonian because it is not generating unitary transformations.

While SPQM has been considered as far back as [41], it has not been fully solved before. There are two other ways of measuring P and Q simultaneously that are important to distinguish from SPQM: the Arthurs-Kelly measurement [67,68] and the Goetsch-Graham-Wiseman (GGW) model of heterodyne detection [33–37]. The Arthurs-Kelly measurement has the same system-meter interaction as the SPQM process but is different in that Arthurs and Kelly imagine continually interacting the same two meters with the system until

the measurement is terminated, whereas in the SPQM process, the system interacts with many pairs of meters successively, registering the complex Wiener path $dw_{[0,T]}$. The GGW model of heterodyne detection has the same many-meter model as the SPQM process, but each system-meter interaction is energy-conserving (the so-called leaky cavity), whereas the system-meter interaction of the SPQM process is the meter displacement given in Equation (36).

The total operation of the SPQM process is a Wiener-like path integral,

$$\mathcal{Z}_T \equiv \int \mathcal{D}\mu[dw_{[0,T]}] \mathcal{O}\left(L[dw_{[0,T]}]\right), \quad (58)$$

which is absolutely trivial to solve,

$$\mathcal{Z}_T = \left(\int d\mu(dw) \mathcal{O}\left(e^{-H_0 2\kappa dt + Q\sqrt{\kappa}dW_t^q + P\sqrt{\kappa}dW_t^p}\right) \right)^{\circ T/dt} = e^{-\frac{1}{2}\kappa T(\text{ad}_Q^2 + \text{ad}_P^2)}. \quad (59)$$

The interest of this article, however, is entirely in the manifold diffusion process defined by SPQM, where Equation (57) is understood to define sample paths in a finite-dimensional manifold. The infinitesimal generators of these sample paths are Q , P , and H_0 . By simply considering their Lie brackets to the first order,

$$[Q, P] = i\Omega, \quad (60)$$

$$[H_0, Q] = -iP, \quad (61)$$

$$[H_0, P] = iQ, \quad (62)$$

and second order,

$$[[H_0, Q], Q] = -\Omega, \quad (63)$$

$$[[H_0, P], P] = -\Omega, \quad (64)$$

$$[H_0, [H_0, Q]] = Q, \quad (65)$$

$$[H_0, [H_0, P]] = P, \quad (66)$$

we can see that SPQM defines a seven-dimensional manifold, a representation of what we will call the *instrumental Weyl-Heisenberg group*,

$$G = \text{IWH} \equiv \left\{ e^{-H_0 r} e^{Qq} e^{Pp} D_\alpha e^{i\Omega\phi} e^{-\Omega\ell} : r \in \mathbb{R}, q, p \in \mathbb{R}, \alpha \in \mathbb{C}, \phi, \ell \in \mathbb{R} \right\}. \quad (67)$$

There is a fourth and final Weyl-Heisenberg group worth defining, the six-dimensional *complex Weyl-Heisenberg group*,

$$\text{CWH} \equiv \left\{ e^{Qq} e^{Pp} e^{-\Omega\ell} D_\alpha e^{i\Omega\phi} : q, p \in \mathbb{R}, \alpha \in \mathbb{C}, \phi, \ell \in \mathbb{R} \right\}, \quad (68)$$

which is a maximal normal subgroup of the IWH group called the derived subgroup. In Lie-group terminology, the IWH group is said to be solvable while CWH is said to be nilpotent or unipotent [51,52]. In particular, the derived series of G is

$$G = \text{IWH} \triangleright \text{CWH} \triangleright Z \triangleright 1, \quad (69)$$

where the *center* of G is defined,

$$Z \equiv \left\{ e^{i\Omega\phi} e^{-\Omega\ell} : \phi, \ell \in \mathbb{R} \right\}. \quad (70)$$

Before proceeding, it is useful to introduce the complex-amplitude (also known as annihilation) operator

$$a \equiv \frac{1}{\sqrt{2}}(Q + iP), \quad (71)$$

which has the Lie bracket

$$[a, a^\dagger] = \Omega. \quad (72)$$

We have

$$H_0 = \frac{1}{2}(aa^\dagger + a^\dagger a) = a^\dagger a + \frac{1}{2}\Omega. \quad (73)$$

Note also that the displacement operator,

$$D_\alpha = e^{iQ\alpha_2 - iP\alpha_1} = e^{a^\dagger \alpha - a\alpha^*}, \quad (74)$$

has the ordered forms,

$$D_\alpha = e^{-\frac{1}{2}i\Omega\alpha_1\alpha_2} e^{iQ\alpha_2} e^{-iP\alpha_1} = e^{\frac{1}{2}i\Omega\alpha_1\alpha_2} e^{-iP\alpha_1} e^{iQ\alpha_2} \quad (75)$$

$$= e^{-\frac{1}{2}|\alpha|^2\Omega} e^{a^\dagger \alpha} e^{-a\alpha^*} = e^{\frac{1}{2}|\alpha|^2\Omega} e^{-a\alpha^*} e^{a^\dagger \alpha}, \quad (76)$$

which will prove useful in relating coordinate systems and in evaluating right-invariant derivatives. The first form in Equation (76) is usually called normal ordering, and the second form is called antinormal ordering.

2.3. Haar Measure, Dirac Delta, and Kraus-Operator Distribution Function

Many readers will interpret the groups defined in the previous section as matrix groups, where the observables are quantized in the usual way. The exponentials can be understood more abstractly, however, as generating path-connections, and this way of thinking gives rise to what is called the *universal covering group* [40,52,69–72]. We will now start to consider more seriously the points of the IWH group in this universal fashion and think of the instrument of SPQM as a representation of $G = \text{IWH}$. The map $L : \mathbb{C}^{T/dt} \rightarrow \text{GL}(\mathcal{H})$ from the set of paths, $\mathbb{C}^{T/dt}$, to the operator space, $\text{GL}(\mathcal{H})$, can be factored into two maps,

$$L = R \circ \gamma, \quad (77)$$

where $\gamma : \mathbb{C}^{T/dt} \rightarrow G$ maps Wiener paths to the universal covering group and $R : G \rightarrow \text{GL}(\mathcal{H})$ is the representation, mapping the universal cover to the space of linear operators on \mathcal{H} [52]. We have denoted a sample path by $dw_{[0,T]}$, and we now start denoting elements of the instrumental group by x . To drive the notation home, we note that the instrumental group element and Kraus operator associated with a sample path are denoted by

$$x_T = \gamma[dw_{[0,T]})] \quad \text{and} \quad L[dw_{[0,T]})] = R(x_T) = R(\gamma[dw_{[0,T]})]). \quad (78)$$

The distinction between L and γ emphasizes that the time-ordered exponential of Equation (57) actually defines a diffusion problem on the instrumental group that is independent of the spectral information inherent in the definition of a linear operator. In particular, this means the entire analysis of this article is independent of whether or how H_0 is quantized (remember that H_0 is quantized in the usual way the moment \mathcal{H} is assumed to be irreducible and to have a ground state).

As with every finite-dimensional Lie group, $G = \text{IWH}$ has a right-invariant *Haar measure* [18,43–45],

$$d\mu(xg) = d\mu(x). \quad (79)$$

As is *not* always the case for (solvable) Lie groups [51], it turns out that this right-invariant measure is also equal to the left-invariant measure,

$$d\mu(x) = d\mu(gx), \quad (80)$$

and this left invariance will turn out to be important for the interpretation of the SPQM process as a diffusion, as will be seen in the very next section. Comments about the existence and uniqueness of the Haar measure will be given in Appendices D and E, but first it is worth taking it for granted and appreciating what can be done with it.

With the Haar measure, we can introduce the accompanying singular distributions or “Dirac deltas” defined by the property (sometimes called reproduction [4,73]) that, for any function f of $G = \text{IWH}$,

$$\int_G d\mu(x) \delta(y, x) f(x) = f(y). \quad (81)$$

From the invariance properties of the Haar distribution, the Dirac delta distributions inherit the corresponding covariance properties,

$$\delta(xg, yg) = \delta(x, y) = \delta(gx, gy). \quad (82)$$

With the Haar measure and the Dirac delta, we can define a universal instrument by adding up all of the Wiener paths that end at the same Kraus operator, starting from the origin. This becomes visible by considering the total operation,

$$\mathcal{Z}_T = \int \mathcal{D}\mu[dw_{[0,T]}] \mathcal{O}(L[dw_{[0,T]}]) \quad (83)$$

$$= \int \mathcal{D}\mu[dw_{[0,T]}] \mathcal{O}(R(\gamma[dw_{[0,T]}])) \quad (84)$$

$$= \int \mathcal{D}\mu[dw_{[0,T]}] \mathcal{O}(R(\gamma[dw_{[0,T]}])) \int_G d\mu(x) \delta(x, \gamma[dw_{[0,T]}]) \quad (85)$$

$$= \int_G d\mu(x) \mathcal{O}(R(x)) \int \mathcal{D}\mu[dw_{[0,T]}] \delta(x, \gamma[dw_{[0,T]}]). \quad (86)$$

In summary, the SPQM process defines a *universal instrument*, which unravels the total operation,

$$\mathcal{Z}_T = \int_G d\mu(x) D_T(x) \mathcal{O}(R(x)), \quad (87)$$

according to a Haar-based *Kraus-operator distribution function* (KOD),

$$D_T(x) \equiv \int \mathcal{D}\mu[dw_{[0,T]}] \delta(x, \gamma[dw_{[0,T]}]), \quad (88)$$

which is defined by a Wiener path integral [10–12,15,16,19]. The total operation is a completely positive, trace-preserving superoperator. The trace-preserving property is equivalent to saying that the POVM elements, $d\mu(x) D_T(x) R(x)^\dagger R(x)$, satisfy a completeness relation,

$$1_{\mathcal{H}} = \int_G d\mu(x) D_T(x) R(x)^\dagger R(x). \quad (89)$$

The term “universal” refers to the fact that this description of the instrument is common to every representation and comes from the concepts of a universal covering group and universal enveloping algebra [40,71,74]. It is important to understand that the universal covering group, IWH, is defined purely by the local structure (that is, the Lie algebra) of the observables and the quadratic term, here H_0 , which accompanies the observables due to the nature of differential positive transformations. In particular, this means $G = \text{IWH}$ is not defined by the Hilbert space of states. This ability to describe the measuring instrument without reference to a Hilbert space is so striking that we give it a name: *universal instrument autonomy* [37,40]. SPQM is in a very special class of universal instruments for which the universal covering group is finite-dimensional; we dub such instruments *principal instruments* [40].

2.4. Diffusion Equation in Terms of Right-Invariant Derivatives

The definition of the Kraus-operator distribution function given in Equation (88) can be thought of as a Feynman-Kac formula [15,16,19] for the solution of a Fokker-Planck-Kolmogorov (FPK) diffusion equation [59,60]. This FPK diffusion equation can be obtained easily with the help of the so-called right-invariant derivatives [39,40,46–48],

$$\overleftarrow{X}[f](x) \equiv \lim_{h \rightarrow 0} \frac{f(e^{Xh}x) - f(x)}{h}, \quad (90)$$

which can be seen to have commutators [40]

$$\overleftarrow{X}\overleftarrow{Y} - \overleftarrow{Y}\overleftarrow{X} = -[\overleftarrow{X}, \overleftarrow{Y}]. \quad (91)$$

This negative sign is why left-invariant derivatives are usually considered instead of right-invariant ones. Nevertheless, because the convention is to consider operators as acting to the right, we are more-or-less stuck with having to consider a right-invariant basis for local transformations.

With the right-invariant derivatives, D_{t+dt} can then be expanded about t in the standard way. We start with

$$D_{t+dt}(x) = \int \mathcal{D}\mu[dw_{[0,t+dt]}] \delta(x, \gamma[dw_{[0,t+dt]}]) \quad (92)$$

$$= \int d\mu(dw_t) \int \mathcal{D}\mu[dw_{[0,t]}] \delta(x, \gamma(dw_t) \gamma[dw_{[0,t]}]), \quad (93)$$

where

$$\gamma(dw_t) = e^{\delta_t}, \quad \delta_t \equiv -H_0 2\kappa dt + Q\sqrt{\kappa} dW^q + P\sqrt{\kappa} dW^p = -H_0 2\kappa dt + a\sqrt{\kappa} dw_t^* + a^\dagger \sqrt{\kappa} dw_t \quad (94)$$

is the purely group-theoretic version of the differential Kraus operator $L(dw_t)$ given in Equation (48); that is, $L(dw_t) = R(\gamma(dw_t))$. Here we also define the *forward generator* δ_t : $\gamma(dw_t) = e^{\delta_t}$ is the fundamental differential positive operator for SPQM, and the forward generator

δ_t is thus the core mathematical object for the theory of the SPQM instrument. Continuing with Equation (93), we have

$$D_{t+dt}(x) = \int d\mu(dw_t) \int \mathcal{D}\mu[dw_{[0,t]}] \delta\left(\gamma(dw_t)^{-1}x, \gamma[dw_{[0,t]}}\right) \quad (95)$$

$$= \int d\mu(dw_t) D_t(\gamma(dw_t)^{-1}x) \quad (96)$$

$$= \int d\mu(dw_t) D_t(e^{-\delta_t}x) \quad (97)$$

$$= \int d\mu(dw_t) e^{-\overleftarrow{\delta}_t} [D_t](x) \quad (98)$$

$$= \int d\mu(dw_t) \left(D_t(x) - \overleftarrow{\delta}_t [D_t](x) + \frac{1}{2} \overleftarrow{\delta}_t [\overleftarrow{\delta}_t [D_t]](x) \right) \quad (99)$$

$$= D_t(x) + \kappa dt \Delta [D_t](x), \quad (100)$$

where

$$\Delta \equiv 2H_o + \frac{1}{2} (\overleftarrow{Q} \overleftarrow{Q} + \overleftarrow{P} \overleftarrow{P}) \quad (101)$$

is the FPK forward generator. Equation (95) is where the left invariance of the Haar measure is used. Equation (96) is the analog of a Chapman-Kolmogorov equation for the distribution function [59,60]. Equation (98) moves e^{δ_t} outside the argument of the KOD to become an exponential of the right-invariant derivative

$$\overleftarrow{\delta}_t = -\overleftarrow{H_o} 2\kappa dt + \overleftarrow{Q} \sqrt{\kappa} dW_t^q + \overleftarrow{P} \sqrt{\kappa} dW_t^p, \quad (102)$$

which we call the *vector-valued SPQM increment*. Equation (99) Taylor-expands the distribution function to the second order—that is, to order dt —as required by the Itô rule for the Wiener outcome increments. The remaining step to the FPK forward generator Δ involves averaging over the Wiener distribution $d\mu(dw_t)$; in this averaging, the deterministic term $-2\kappa dt \overleftarrow{H_o}$ in $\overleftarrow{\delta}_t$ contributes a first-derivative term to the FPK forward generator, whereas the Wiener outcome increment terms in $\overleftarrow{\delta}_t$ contribute second-derivative diffusion terms.

In summary, the KOD of SPQM evolves according to the FPK diffusion equation,

$$\frac{1}{\kappa} \frac{\partial}{\partial t} D_t(x) = \Delta [D_t](x), \quad (103)$$

with initial condition

$$D_0(x) = \delta(x, 1), \quad (104)$$

where the group identity is defined

$$1 \equiv e^0, \quad (105)$$

which is the origin of the IWH group. Equation (103) will be (mostly) solved in Section 3.5 after having established two coordinate systems. The subtlest thing about Equation (103) is remembering that the Lie algebra of the three directions apparent in the equation means that the motion beyond the first order is actually seven-dimensional.

Before proceeding to SDEs, we draw attention to an important property of right-invariant derivatives. The reader might already be thinking about this property by wondering why we did not write $\overleftarrow{\delta}_t$ and Δ in terms of the right-invariant derivatives associated with a and a^\dagger . The reason is that the map from operators to right-invariant derivatives, $X \mapsto \overleftarrow{X}$, is only \mathbb{R} -linear and not \mathbb{C} -linear [40]. This means that \overleftarrow{Q} , \overleftarrow{P} , $i\overleftarrow{Q}$, and $-i\overleftarrow{P}$ are

\mathbb{R} -linearly independent, with \underline{Q} and \underline{P} displacing Kraus operators in positive directions in the IWH group and $i\underline{Q}$ and $-i\underline{P}$ displacing in unitary directions. The right-invariant derivatives \underline{a} , \underline{a}^\dagger , $-i\underline{a}$, and $i\underline{a}^\dagger$ each represent a different way of equally combining displacements in the positive and unitary directions. In view of this, it is instructive to note that the vector-valued SPQM increment of Equation (102) has the form

$$\underline{\delta}_t = -\underline{H}_0 2\kappa dt + \frac{1}{2}\sqrt{\kappa} dw_t (\underline{a} + \underline{a}^\dagger + i\underline{a} - i\underline{a}^\dagger) + \frac{1}{2}\sqrt{\kappa} dw_t^* (\underline{a} + \underline{a}^\dagger - i\underline{a} + i\underline{a}^\dagger). \quad (106)$$

This means, in particular, that

$$\Delta \neq 2H_0 + \frac{1}{2}(\underline{a} \underline{a}^\dagger + \underline{a}^\dagger \underline{a}). \quad (107)$$

2.5. Sample-Path SDEs in Terms of Right-Invariant One-Forms

As has been mentioned, the time-ordered exponential of the SPQM process, given in Equation (57), can be interpreted as defining sample paths in the seven-dimensional manifold $G = \text{IWH}$. Sample paths are usually described by stochastic differential equations (SDEs). We finish this section by explaining how such SDEs can be expressed in terms of the right-invariant structure.

The basis of right-invariant derivatives,

$$\{e_\nu\} \equiv \left\{ -\underline{H}_0, \underline{Q}, \underline{P}, -i\underline{P}, i\underline{Q}, -\underline{\Omega}, i\underline{\Omega} \right\}, \quad (108)$$

defines a dual basis of right-invariant one-forms,

$$\theta^\mu(e_\nu) \equiv \delta_\nu^\mu, \quad (109)$$

In terms of the right-invariant one-forms, the Haar measure has a simple expression in terms of wedge products,

$$d\mu(x) = \theta^{-H_0} \wedge \theta^Q \wedge \theta^P \wedge \theta^{iQ} \wedge \theta^{-iP} \wedge \theta^{-\Omega} \wedge \theta^{i\Omega}. \quad (110)$$

Also, in terms of the right-invariant one-forms, the SDEs equivalent to Equation (57) are obtained by reading off the coefficient conjugate to the corresponding generator in the vector-valued SPQM increment of Equation (102),

$$\theta^{i\Omega}(\underline{\delta}_t) = 0, \quad (111)$$

$$\theta^{-\Omega}(\underline{\delta}_t) = 0, \quad (112)$$

$$\theta^{-iP}(\underline{\delta}_t) = 0, \quad (113)$$

$$\theta^{iQ}(\underline{\delta}_t) = 0, \quad (114)$$

$$\theta^Q(\underline{\delta}_t) = \sqrt{\kappa} dW_t^q, \quad (115)$$

$$\theta^P(\underline{\delta}_t) = \sqrt{\kappa} dW_t^p, \quad (116)$$

$$\theta^{-H_0}(\underline{\delta}_t) = 2\kappa dt. \quad (117)$$

These SDEs can be broken into two types: the first-order SDEs,

$$\theta^{-H_0}(\delta_t) = 2\kappa dt, \quad \theta^Q(\delta_t) = \sqrt{\kappa} dW_t^q, \quad \theta^P(\delta_t) = \sqrt{\kappa} dW_t^p, \quad (118)$$

and the Pfaffians,

$$\theta^{-iP}(\delta_t) = \theta^{iQ}(\delta_t) = \theta^{-\Omega}(\delta_t) = \theta^{i\Omega}(\delta_t) = 0. \quad (119)$$

These equations will be solved in Section 3.4 after having established a coordinate system.

The SDEs of Equations (118) and (119) are almost obvious by definition, but there is a subtlety that requires attention. The right-invariant derivatives and one-forms live in the spaces tangent and cotangent to the group manifold G and thus are based on linear transformations that use the chain rule of ordinary calculus. Hence, the stochastic equations that come from the right-invariant one-forms are Stratonovich-form SDEs [59,60]—this means mid-point evaluation of coefficients of stochastic increments—and should be converted to the Itô-form SDEs in which coefficients are evaluated at the beginning of the increment. In the context of the IWH group, the only place this subtlety makes a difference is in the SDEs that come from $\theta^{i\Omega}$ and $\theta^{-\Omega}$. Jackson and Caves [39,40] introduced the *modified Maurer-Cartan stochastic differential* (MMCSd) as a way to get to the Itô-form equations directly. The MMCSd is an example of the Itô correction in SDEs [59,60], specifically, the Itô correction that arises when one transforms between a stochastic variable and the exponential function of that variable, as occurs in the forward generator of Equation (94). In this paper, we get directly to Itô-form SDEs in a different way when we consider the Harish-Chandra decomposition in Section 3.2.

A further subtlety about the right-invariant one-forms is that they have “curl” in the same sense as Gibbs would have defined. In the standard language of forms, this is because the exterior algebra of forms is equivalent to the Lie algebra of derivatives [46],

$$[e_\mu, e_\nu] = -c_{\mu\nu}^\lambda e_\lambda \iff d\theta^\lambda = \frac{1}{2} c_{\mu\nu}^\lambda \theta^\mu \wedge \theta^\nu. \quad (120)$$

This equivalence is standard in modern differential topology, but an introduction is included in Appendix A; though no use will be made of the “curls” $d\theta^\lambda$ in this article, they are given for completeness in Equations (A15)–(A21).

3. The IWH Group and Two Coordinate Systems

Having introduced the instrumental Weyl-Heisenberg group, $G = \text{IWH}$, a coordinate system needs to be established so that we can locate the sample paths of the SPQM process and follow their propagation. If the concept of a universal covering group introduced in the previous section is unclear, seeing how Equation (57) is equivalent to a set of coordinate SDEs should help the reader appreciate that $G = \text{IWH}$ and the SPQM process are independent of matrix representation.

We will use two coordinate systems analogous to what are called Harish-Chandra [52–54] (also known as “Gauss” [2,51,55]) decompositions and Cartan decompositions [50–52,54]. These decompositions were originally designed in the context of semisimple Lie groups [49,51,52,75,76], of which the IWH group is quite the opposite (in the sense of the Levi-Malcev decomposition). In spite of this distinction, it is very useful to think of the IWH group in many ways *as if* it were semisimple. The Harish-Chandra decomposition is easier to prove first and will allow us to make connections between the SPQM process and two processes more familiar to physicists, the Ornstein-Uhlenbeck process and the Goetsch-

Graham-Wiseman (GGW) heterodyne measuring process. The Cartan decomposition is better suited for considering the POVM.

Section 3.1 identifies the analogs of the various elements of semisimple group theory. Section 3.2 proves the Harish-Chandra decomposition of the IWH group in a way that also produces the corresponding Itô-form coordinate SDEs, which are immediately recognized and solved. Section 3.3 introduces the Cartan decomposition and the transformations to Harish-Chandra coordinates and presents the right-invariant derivatives and one-forms in both coordinate systems. Section 3.4 solves the SDEs of the SPQM process in both Cartan and Harish-Chandra coordinates. Section 3.5 solves most of the FPK diffusion equation from the SPQM process in Cartan coordinates, by which we mean introducing the Cartan-section reduced distribution function and solving for it. Section 3.6 explains how the solution of the FPK diffusion equation means that the POVM of the SPQM process offers an alternative perspective on the meaning of energy quantization.

3.1. Usual Elements of Semisimple Lie Group Theory

As introduced in the previous section, the Lie group of interest is the so-called instrumental Weyl–Heisenberg group

$$G = \text{IWH} \equiv \left\{ e^{-H_0 r} e^{Qq} e^{Pp} D_\alpha e^{i\Omega\phi} e^{-\Omega\ell} : r \in \mathbb{R}, q, p \in \mathbb{R}, \alpha \in \mathbb{C}, \phi, \ell \in \mathbb{R} \right\}. \quad (121)$$

Although G is literally solvable, with derived series,

$$G = \text{IWH} \triangleright \text{CWH} \triangleright Z \triangleright 1, \quad (122)$$

and center,

$$Z \equiv \left\{ e^{i\Omega\phi} e^{-\Omega\ell} : \phi, \ell \in \mathbb{R} \right\}, \quad (123)$$

it can be navigated in much the same way as a semisimple group. While some of the terminology [52] will be used here, the theory of semisimple groups will be more-or-less glossed over. The purpose of this section is basically to label the various subgroups that will prove to be both meaningful and useful for navigating the IWH group and, therefore, understanding SPQM. The significance of these subgroups should become apparent as they are applied.

The map from the SPQM instrument to the SPQM POVM,

$$\pi(x) = x^\dagger x, \quad (124)$$

defines the SPQM POVM as similar to a symmetric space, albeit a non-Riemannian one, with *Cartan group involution*

$$x^\dagger \equiv x^{-\dagger} = (x^{-1})^\dagger = (x^\dagger)^{-1}. \quad (125)$$

The subgroup of transformations that are even under the Cartan involution is the usual unitary Weyl–Heisenberg group of Equation (44):

$$K \equiv \left\{ x \in G : x^\dagger = x \right\} = \pi^{-1}(1) = \text{UWH}. \quad (126)$$

On the other hand, the remainder of G displaces from the origin of the *symmetric space*,

$$\mathcal{E} \equiv \pi(G) \cong K \backslash G. \quad (127)$$

Considering the conjugation action of K on \mathcal{E} , almost all of the K -conjugacy classes can be parameterized by the *Cartan subgroup*,

$$A \equiv \left\{ e^{-H_0 r} e^{-\Omega \ell} : r, \ell \in \mathbb{R} \right\}, \quad (128)$$

and regular POVM elements are invariant under the *commutant*,

$$M \equiv \left\{ k \in K : \forall a \in A, k a k^{-1} = a \right\} = \left\{ e^{i\Omega \phi} : \phi \in \mathbb{R} \right\}. \quad (129)$$

Thus the K -conjugacy classes have the topology of the familiar phase space,

$$K/M \cong \mathbb{C}. \quad (130)$$

Indeed, it is the “almost all” feature where G and \mathcal{E} depart from semisimple groups and Riemannian symmetric spaces, since positive transformations of the form e^{Qq+Pp} are characteristically not in the conjugacy classes of A (see Appendix B for an additional perspective). Finally, important also are the *maximal nilpotent (or unipotent) subgroup*,

$$N \equiv \left\{ e^{a\mu^*} : \mu \in \mathbb{C} \right\}, \quad (131)$$

and, perhaps the most important, the *Borel subgroup*,

$$B \equiv A \ltimes N = \left\{ e^{-H_0 r} e^{-\Omega \ell} e^{a\mu^*} : r, \ell \in \mathbb{R}, \mu \in \mathbb{C} \right\}. \quad (132)$$

This group-theoretic context now in hand, we stress that the most important groups for what follows are $G = \text{IWH}$ itself, the center Z of G , and the quotient group $G/Z = \text{IWH}/Z$. The center Z contains a phase, which is of no importance, and a normalization, which is the main source of difficulty in analyzing SPQM. The cosets $Zx \in G/Z$ are parametrized by what we call the ruler, r , and by two complex phase-space parameters. One of these complex phase-space parameters is associated with the POVM, and the other parametrizes a post-measurement displacement operator. We call G/Z the *reduced instrumental Weyl-Heisenberg (RIWH) group*. It is isomorphic to the adjoint group of G , but given the way we will use multipliers on the cosets Zx , we prefer to think of G as a central extension of G/Z .

3.2. Harish-Chandra Decomposition and SDEs as a Proof by Transfinite Induction

The seven-dimensional instrumental Weyl-Heisenberg group G affords a *Harish-Chandra decomposition*, $G = N^\dagger M A N$, where every element can be decomposed into the form

$$x = e^{a^\dagger v} e^{-H_0 r + \Omega z} e^{a\mu^*}, \quad (133)$$

defining seven *Harish-Chandra coordinates* $(v, r, z, \mu) \in \mathbb{C} \times \mathbb{R} \times \mathbb{C} \times \mathbb{C}$. We often break the complex coordinates into real and imaginary parts,

$$v = \frac{v_1 + iv_2}{\sqrt{2}}, \quad \mu = \frac{\mu_1 + i\mu_2}{\sqrt{2}}, \quad z = -s + i\psi. \quad (134)$$

The coordinate r we call the ruler, v and μ are the postmeasurement and POVM phase-plane coordinates, and z is the IWH group center coordinate.

A proof that this decomposition exists for every element of the SPQM process of Equation (57) is not difficult if we allow ourselves to apply transfinite induction: at time $t = dt$ (the first infinitesimal increment), it is easy to see that the decomposition exists because (see Equation (94))

$$x_{dt} = e^\delta = e^{-H_0 2\kappa dt + a^\dagger \sqrt{\kappa} dw_0 + a \sqrt{\kappa} dw_0^*} \quad (135)$$

$$= e^{-H_0 2\kappa dt} e^{a^\dagger \sqrt{\kappa} dw_0 + a \sqrt{\kappa} dw_0^*} \quad (136)$$

$$= e^{a^\dagger \sqrt{\kappa} dw_0} e^{-H_0 2\kappa dt + \Omega \frac{1}{2} \kappa |dw_0|^2} e^{a \sqrt{\kappa} dw_0^*}. \quad (137)$$

Trivially, this also means that the decomposition exists for any finite integer n and infinitesimal time $t = ndt$ simply because the one-parameter subgroups commute to infinitesimal order so long as the Wiener increments are independent. Now for the transfinite step. If we assume that the decomposition holds for a finite time t ,

$$x_t = e^{a^\dagger v_t} e^{-H_0 r_t + \Omega z_t} e^{a \mu_t^*}, \quad (138)$$

then an increment later in the SPQM process, we have

$$x_{t+dt} = e^{\delta_t} x_t \quad (139)$$

$$= e^{a^\dagger \sqrt{\kappa} dw_t} e^{-H_0 2\kappa dt + \Omega \frac{1}{2} \kappa |dw_t|^2} e^{a \sqrt{\kappa} dw_t^*} e^{a^\dagger v_t} e^{-H_0 r_t + \Omega z_t} e^{a \mu_t^*} \quad (140)$$

$$= e^{a^\dagger \sqrt{\kappa} dw_t} e^{-H_0 2\kappa dt} e^{a^\dagger v_t} e^{a \sqrt{\kappa} dw_t^*} e^{-H_0 r_t + \Omega(z_t + \frac{1}{2} \kappa |dw_t|^2 + v_t \sqrt{\kappa} dw_t^*)} e^{a \mu_t^*} \quad (141)$$

$$= e^{a^\dagger (e^{-2\kappa dt} v_t + \sqrt{\kappa} dw_t)} e^{-H_0(r_t + 2\kappa dt) + \Omega(z_t + \frac{1}{2} \kappa |dw_t|^2 + v_t \sqrt{\kappa} dw_t^*)} e^{a(\mu_t + e^{-r_t} \sqrt{\kappa} dw_t)^*}. \quad (142)$$

This concludes the proof of the Harish-Chandra decomposition for the SPQM process and $G = \text{IWH}$.

A consequence of the proof is that the SPQM process in Harish-Chandra coordinates is equivalent to the Itô-form SDEs [56–60],

$$dr_t = 2\kappa dt, \quad (143)$$

$$dv_t = -2\kappa dt v_t + \sqrt{\kappa} dw_t, \quad (144)$$

$$d\mu_t = e^{-r_t} \sqrt{\kappa} dw_t, \quad (145)$$

$$-ds_t + id\psi_t = dz_t = \frac{1}{2} \kappa |dw_t|^2 + v_t \sqrt{\kappa} dw_t^*. \quad (146)$$

Although these are Itô-form SDEs, notice that we did not use the Itô rule in deriving them; in particular, we do not set $|dw_t|^2 = dt$ in the SDE for the center coordinate z .

We now solve these equations for the initial condition $r_0 = v_0 = \mu_0 = z_0 = 0$. These initial values are chosen so that $x_0 = 1$, in agreement with the δ -function initial condition for KOD, which is given in Equation (104). It is straightforward to see that the first three SDEs have as solutions

$$r_T = 2\kappa T, \quad (147)$$

$$v_T = \int_0^{T-} \sqrt{\kappa} dw_t e^{-2\kappa(T-t)}, \quad (148)$$

$$\mu_T = \int_0^{T-} \sqrt{\kappa} dw_t e^{-2\kappa t}, \quad (149)$$

where $T_- \equiv T - dt$. The fourth equation is solved by plugging the solution for v_t into the equation for z_t and integrating, with the result that

$$-s_T + i\psi_T = z_T = \frac{1}{2} \int_0^{T_-} \int_0^{T_-} \kappa dw_t^* dw_s e^{-2\kappa|t-s|} (1 + \text{sgn}(t-s)), \quad (150)$$

where

$$\text{sgn}(u) = \begin{cases} 1, & u > 0, \\ 0, & u = 0, \\ -1, & u < 0, \end{cases} \quad (151)$$

is the sign function. There being subtleties in deriving and interpreting this solution, it is worked out carefully in Appendix C. The solutions for the real and imaginary parts of z are

$$-s_T = \frac{1}{2} \int_0^{T_-} \int_0^{T_-} \kappa dw_t^* dw_s e^{-2\kappa|t-s|}, \quad (152)$$

$$i\psi_T = \frac{1}{2} \int_0^{T_-} \int_0^{T_-} \kappa dw_t^* dw_s e^{-2\kappa|t-s|} \text{sgn}(t-s). \quad (153)$$

The SPQM posterior phase-point variable, v_T from Equation (148), is a linear functional in the measurement record and can be recognized as an Ornstein-Uhlenbeck (OU) process [59–61]. The SPQM prior phase-point variable, μ_T from Equation (149), is a linear functional in the measurement record and can be recognized as the same as the Goetsch-Graham-Wiseman (GGW) heterodyne process. (To be clear, the equations for the remaining variables of the GGW heterodyne process would instead be $r_T^{\text{GGW}} = 2\kappa T$, $v_T^{\text{GGW}} = 0$, and $z_T^{\text{GGW}} = \kappa T$ [37]; in particular, the GGW heterodyne process corresponds to the Borel subgroup B .) The center SPQM variable, z_T from Equation (150), is a quadratic functional in the measurement record; the solution is identical to that obtained by deriving a quantum fluctuation-dissipation theorem for the correlation of v and μ .

It is worth taking a moment to reflect further on Equations (147)–(150). Equation (147) tells us that if H_0 is quantized in the standard way, then the Kraus operators collapse, with an e -folding time $\tau = 1/2\kappa$, to a scaled outer product of coherent states of the form

$$L_{T \gg 1/\kappa} \sim e^{-\kappa T + z_T} e^{a^\dagger v_T} |0\rangle\langle 0| e^{a \mu_T^*} = e^{i\psi_T} e^{\frac{1}{2}(|v_T|^2 + |\mu_T|^2) - s_T - \kappa T} |v_T\rangle\langle \mu_T|. \quad (154)$$

The complementarity in time of the OU and GGW processes (Equations (148) and (149)) is also interesting: whereas the post-measurement variable v_T of Equation (148) depends only on the end of the outcome register, the POVM variable μ_T of Equation (149) depends only on the beginning of the register. It is thus reasonably clear that the POVM of the SPQM process culminates in the usual “measurement in the coherent-state basis”,

$$1_{\mathcal{H}} = \int \frac{d^2\mu}{\pi} |\mu\rangle\langle \mu|, \quad (155)$$

where $d^2\mu = \frac{1}{2} d\mu_1 d\mu_2$. This is just like GGW heterodyne detection, except that the post-measurement state is not vacuum, but instead is scrambled over phase space. That the POVM variable μ_T ceases to evolve after a few e -foldings seems to be an important feature in the interpretation of SPQM as a measuring process [66]. All this said, there remains the elephant in the room that prompts us to label these conclusions as “reasonably clear”: the elephant is the factor that scales the long-time Kraus operator in Equation (154); although Equation (154) makes it absolutely clear that the Kraus operator approaches an outer product of coherent states, what is not clear at all is how the POVM approaches a

uniform distribution of coherent states, as required by the completeness relation given in Equation (155). Figuring this out is, in some sense, what the rest of the paper is about.

3.3. Cartan Decomposition and Various Transformations

As far as what the SPQM process is ultimately doing in time, Section 3.2 in many ways says it all—except for that elephant in the room. We now turn our attention to a detailed understanding of the measuring process at finite times, which requires addressing the elephant. Key to this is the realization that the Harish-Chandra decomposition is not well suited to telling us the behavior of the POVM. Rather, a Cartan decomposition of the instrumental Weyl-Heisenberg group, $G = KAK$, works best for this purpose. A straightforward calculation (left to Appendix B) shows that almost every group element—and every Kraus operator—with a Harish-Chandra decomposition also affords a *Cartan decomposition*,

$$x = \left(D_\beta e^{i\Omega\phi} \right) e^{-H_0 r - \Omega\ell} D_\alpha^{-1}, \quad (156)$$

and, therefore, the POVM elements decompose as

$$x^\dagger x = D_\alpha e^{-H_0 2r - \Omega 2\ell} D_\alpha^{-1}, \quad (157)$$

where the *Cartan coordinates*, $(\beta, \phi, r, \ell, \alpha) \in \mathbb{C} \times \mathbb{R} \times \mathbb{R} \times \mathbb{R} \times \mathbb{C}$, can be obtained by the coordinate transformations (the ruler r is shared by the two coordinate systems),

$$\beta = \frac{e^r v + \mu}{2 \sinh r}, \quad (158)$$

$$\alpha = \frac{e^r \mu + v}{2 \sinh r}, \quad (159)$$

$$\ell = s - f, \quad (160)$$

$$\phi = \psi - \zeta. \quad (161)$$

Here f and ζ are functions of the ruler and the phase-plane coordinates, that is, functions on the RIWH G/Z ,

$$f = \frac{(|v|^2 + |\mu|^2)e^r + v^* \mu + v \mu^*}{4 \sinh r} = \frac{|v + \mu|^2}{4(1 - e^{-r})} + \frac{|v - \mu|^2}{4(1 + e^{-r})} \quad (162)$$

$$\begin{aligned} &= \frac{1}{2} \left(|\beta|^2 + |\alpha|^2 - \beta^* \alpha e^{-r} - \beta \alpha^* e^{-r} \right) = \frac{1 - e^{-r}}{4} |\beta + \alpha|^2 + \frac{1 + e^{-r}}{4} |\beta - \alpha|^2, \\ \zeta &= \frac{v^* \mu - v \mu^*}{4i \sinh r} = \frac{(v - \mu)^*(v + \mu) - (v + \mu)(v - \mu)^*}{8i \sinh r} \\ &= e^{-r} \frac{\beta^* \alpha - \beta \alpha^*}{2i} = e^{-r} \frac{(\beta - \alpha)^*(\beta + \alpha) - (\beta + \alpha)(\beta - \alpha)^*}{4i}. \end{aligned} \quad (163)$$

Notice the singularity in Cartan coordinates at $r = 0$, about which there is further discussion throughout the remainder of the paper and, in particular, in Appendix B.

As we wish to solve Equations (103) and (111)–(117), more relevant to our purpose are the transformations from the right-invariant moving frame to the Cartan coordinate frame. A calculation of these frame transformations can be found in Appendix D, but we include them here for continuity; we also include the transformations to the Harish-Chandra coordinate frame, which are worked out in Appendix E. For the FPK diffusion equation of Equation (103), we require the transformations of the derivatives,

$$\begin{aligned}
\underleftarrow{i\Omega} &= \partial_\phi & &= \partial_\psi \\
-\underleftarrow{\Omega} &= \partial_\ell & &= \partial_s \\
-\underleftarrow{iP} &= \partial_{\beta_1} - \frac{1}{2}\beta_2\partial_\phi & &= \partial_{\nu_1} - e^{-r}\partial_{\mu_1} + \frac{1}{2}\nu_1\partial_s - \frac{1}{2}\nu_2\partial_\psi \\
\underleftarrow{iQ} &= \partial_{\beta_2} + \frac{1}{2}\beta_1\partial_\phi & &= \partial_{\nu_2} - e^{-r}\partial_{\mu_2} + \frac{1}{2}\nu_2\partial_s + \frac{1}{2}\nu_1\partial_\psi \\
\underleftarrow{Q} &= \nabla_1 - \beta_1\partial_\ell + \frac{\beta_2 \cosh r - \alpha_2}{2 \sinh r} \partial_\phi & &= \nabla_1 - \frac{1}{2}\nu_1\partial_s + \frac{1}{2}\nu_2\partial_\psi \\
\underleftarrow{P} &= \nabla_2 - \beta_2\partial_\ell - \frac{\beta_1 \cosh r - \alpha_1}{2 \sinh r} \partial_\phi & &= \nabla_2 - \frac{1}{2}\nu_2\partial_s - \frac{1}{2}\nu_1\partial_\psi \\
-\underleftarrow{H_o} &= \partial_r^C - \beta_1\nabla_1 - \beta_2\nabla_2 + \frac{\beta_1^2 + \beta_2^2}{2}\partial_\ell + \frac{\beta_1\alpha_2 - \beta_2\alpha_1}{2 \sinh r} \partial_\phi & &= \partial_r^{\text{HC}} - \nu_1\partial_{\nu_1} - \nu_2\partial_{\nu_2}
\end{aligned} \quad (164)$$

where

$$\nabla_j \equiv \frac{1}{\sinh r} (\partial_{\alpha_j} + \cosh r \partial_{\beta_j}) = \partial_{\nu_j} + e^{-r} \partial_{\mu_j}. \quad (165)$$

The two coordinate systems share the coordinate r , but the partial derivative with respect to r is, of course, different in the two systems. In the above equation we distinguish ∂_r in the two systems, but we do this nowhere else because it is always clear in which coordinate system we are operating. It is worth recording for working in terms of the complex phase-space variables that

$$\nabla \equiv \frac{1}{\sqrt{2}} (\nabla_1 - i\nabla_2) = \frac{1}{\sinh r} (\partial_\alpha + \cosh r \partial_\beta) = \partial_\nu + e^{-r} \partial_\mu, \quad (166)$$

$$\nabla^* \equiv \frac{1}{\sqrt{2}} (\nabla_1 + i\nabla_2) = \frac{1}{\sinh r} (\partial_{\alpha^*} + \cosh r \partial_{\beta^*}) = \partial_{\nu^*} + e^{-r} \partial_{\mu^*}. \quad (167)$$

For the SDEs of Equations (111)–(117), we require the transformations of the one-forms,

$$\begin{aligned}
\theta^{i\Omega} &= d\phi + \frac{1}{2}(\beta_2 d\beta_1 - \beta_1 d\beta_2) + \frac{1}{2}(\alpha_2 d\alpha_1 - \alpha_1 d\alpha_2) + \cosh r (\beta_1 d\alpha_2 - \beta_2 d\alpha_1) &= d\psi + \frac{1}{2}e^r (\nu_1 d\mu_2 - \nu_2 d\mu_1) \\
\theta^{-\Omega} &= d\ell + \frac{1}{2}(\beta_1^2 + \beta_2^2)dr + \sinh r (\beta_1 d\alpha_1 + \beta_2 d\alpha_2) &= ds + \frac{1}{2}e^r (\nu_1 d\mu_1 + \nu_2 d\mu_2) \\
\theta^{-iP} &= d\beta_1 - \cosh r d\alpha_1 &= \frac{1}{2}(d\nu_1 - e^r d\mu_1 + \nu_1 dr) \\
\theta^{iQ} &= d\beta_2 - \cosh r d\alpha_2 &= \frac{1}{2}(d\nu_2 - e^r d\mu_2 + \nu_2 dr) \\
\theta^Q &= \beta_1 dr + \sinh r d\alpha_1 &= \frac{1}{2}(d\nu_1 + e^r d\mu_1 + \nu_1 dr) \\
\theta^P &= \beta_2 dr + \sinh r d\alpha_2 &= \frac{1}{2}(d\nu_2 + e^r d\mu_2 + \nu_2 dr) \\
\theta^{-H_o} &= dr &= dr
\end{aligned} \quad (168)$$

With the one-form transformations in hand, the Haar measure, written in terms of the one-forms in Equation (110), becomes in the two coordinate systems,

$$d^7\mu(x) = d\phi d\ell \frac{d^2\beta}{\pi} dr \sinh^2 r \frac{d^2\alpha}{\pi} = d\psi ds \frac{d^2\nu}{2\pi} dr e^{2r} \frac{d^2\mu}{2\pi}. \quad (169)$$

Here and elsewhere, complex phase-plane measures are denoted by $d^2\beta = \frac{1}{2}d\beta_1 d\beta_2$. The factors of $1/\pi$ in the Cartan measure are conventional in quantum optics and ultimately come from the coherent-state completeness relation of Equation (155). The factors of $1/2\pi$ in Harish-Chandra coordinates then follow from the transformation from Cartan to Harish-Chandra coordinates. The left invariance of these measures is discussed in Appendices D and E. While we are considering measures, we should record the reduced measure on the five-dimensional group $\text{RIWH} = \text{IWH}/\mathbb{Z}$:

$$d^5\mu(Zx) = \frac{d^7\mu(x)}{d\mu(Z)} = \frac{d^2\beta}{\pi} dr \sinh^2 r \frac{d^2\alpha}{\pi} = \frac{d^2\nu}{2\pi} dr e^{2r} \frac{d^2\mu}{2\pi}. \quad (170)$$

Here $d\mu(Z) = d\phi d\ell = d\psi ds$ is the measure on the center Z .

Accompanying the coordinate Haar measure $d^7\mu(x)$ are the coordinate forms of the conjugate δ -function,

$$\delta(x, x') = \delta(\phi - \phi') \delta(\ell - \ell') \frac{1}{\sinh^2 r} \delta(r - r') \pi \delta^2(\beta - \beta') \pi \delta^2(\alpha - \alpha') \quad (171)$$

$$= \delta(\psi - \psi') \delta(s - s') e^{-2r} \delta(r - r') 2\pi \delta^2(\nu - \nu') 2\pi \delta^2(\mu - \mu'). \quad (172)$$

There are obvious coordinate forms for the δ -function $\delta(Zx, Zx')$ that is conjugate to $d^5\mu(Zx)$:

$$\delta(Zx, Zx') = \frac{1}{\sinh^2 r} \delta(r - r') \pi \delta^2(\beta - \beta') \pi \delta^2(\alpha - \alpha') \quad (173)$$

$$= e^{-2r} \delta(r - r') 2\pi \delta^2(\nu - \nu') 2\pi \delta^2(\mu - \mu'). \quad (174)$$

We are especially interested in $\delta(x, 1)$. As the identity 1 has Harish-Chandra coordinates $\phi = \ell = r = 0$ and $\nu = \mu = 0$, we have

$$\delta(x, 1) = \delta(\phi) \delta(s) \delta(r) 2\pi \delta^2(\nu) 2\pi \delta^2(\mu). \quad (175)$$

The Cartan form of $\delta(x, 1)$ requires more attention because of the coordinate singularity at $r = 0$ (see Appendix B for discussion). The singularity is about more than just the $1/\sinh^2 r$ in the Cartan form of the δ -function, although that singularity is the root of the difficulties that require attention. We provide the necessary attention in Appendix F.

3.4. Solving the SDEs

Section 2.5 left off by showing that SPQM corresponds to the right-invariant Stratonovich-form SDEs of Equations (111)–(117). With the frame transformations from Equation (168) of the previous section at hand, the three first-order SDEs, given in Equation (118), find the following expressions,

$$2\kappa dt = \theta^{-H_0}(\delta_t) = dr, \quad (176)$$

$$\begin{aligned} \sqrt{\kappa} dw_t &= \frac{1}{\sqrt{2}} \left(\theta^Q(\delta_t) + i\theta^P(\delta_t) \right) = \beta dr + \sinh r d\alpha \\ &= \frac{1}{2} (dv + e^r d\mu + v dr), \end{aligned} \quad (177)$$

and the four Pfaffians, displayed in Equation (119), give

$$\begin{aligned} 0 &= \frac{1}{\sqrt{2}} \left(\theta^{-iP}(\delta_t) + i\theta^{iQ}(\delta_t) \right) = d\beta - \cosh r d\alpha \\ &= \frac{1}{2} (dv - e^r d\mu + v dr), \end{aligned} \quad (178)$$

$$\begin{aligned} 0 &= \theta^{-\Omega}(\delta_t) = d\ell + \frac{1}{2} (\beta_1^2 + \beta_2^2) dr + \sinh r (\beta_1 d\alpha_1 + \beta_2 d\alpha_2) \\ &= ds + \frac{1}{2} e^r (v_1 d\mu_1 + v_2 d\mu_2), \end{aligned} \quad (179)$$

$$\begin{aligned} 0 &= \theta^{i\Omega}(\delta_t) = d\phi + \frac{1}{2} (\beta_2 d\beta_1 - \beta_1 d\beta_2) + \frac{1}{2} (\alpha_2 d\alpha_1 - \alpha_1 d\alpha_2) + \cosh r (\beta_1 d\alpha_2 - \beta_2 d\alpha_1) \\ &= d\psi + \frac{1}{2} e^r (v_1 d\mu_2 - v_2 d\mu_1), \end{aligned} \quad (180)$$

where the evaluation of the coordinate one-forms on the vector-valued SPQM increment, δ_t , is no longer denoted.

Equation (176) is the ruler equation given previously in Equation (143); it has the solution $r_t = 2\kappa t$ for the initial condition $r_0 = 0$. Summing and differencing Equations (177) and (178), one finds that

$$\begin{aligned}\sqrt{\kappa} dw_t &= d\beta + \beta dr - e^{-r} d\alpha \\ &= dv + v dr,\end{aligned}\quad (181)$$

$$\begin{aligned}\sqrt{\kappa} dw_t &= -d\beta + \beta dr + e^r d\alpha \\ &= e^r d\mu.\end{aligned}\quad (182)$$

As was discussed in Section 2.5, these are Stratonovich-form SDEs, which means that the coefficients are evaluated at the midpoint, $t + \frac{1}{2} dt$, of the increment—technically, the midpoint has no status in the stochastic calculus, so one can regard midpoint evaluation as $a_{t+dt/2} = \frac{1}{2}(a_t + a_{t+dt}) = a_t + \frac{1}{2}da_t$ —but for these SDEs, that evaluation does not produce an Itô correction (because $dr_t = 2\kappa dt$ has no stochastic term), and the equations can be read as Itô-form SDEs, with the coefficients evaluated at the beginning of the increment. When r and dr are substituted into these SDEs, the equations for the Harish-Chandra complex phase-point coordinates, v and μ , become Equations (144) and (145), with the solutions shown in Equations (148) and (149) for the initial conditions $v_0 = 0$ and $\mu_0 = 0$.

The SDEs for Cartan phase-space points are

$$d\beta = \cosh 2\kappa t d\alpha, \quad (183)$$

$$d\alpha = \operatorname{csch} 2\kappa t (-2\kappa\beta dt + \sqrt{\kappa} dw_t). \quad (184)$$

For integrating, these SDEs are more profitably written as

$$d(v e^{2\kappa t}) = d(\beta e^{2\kappa t} - \alpha) = e^{2\kappa t} \sqrt{\kappa} dw_t, \quad (185)$$

$$d\mu = d(-\beta e^{-2\kappa t} + \alpha) = e^{-2\kappa t} \sqrt{\kappa} dw_t, \quad (186)$$

and the reason is that these are the same as integrating the SDEs for the Harish-Chandra phase-plane coordinates. The solutions, satisfying initial conditions $v_0 = \mu_0 = 0$,

$$v_T e^{2\kappa T} = \beta_T e^{2\kappa T} - \alpha_T = \int_0^{T-} \sqrt{\kappa} dw_t e^{2\kappa t}, \quad (187)$$

$$\mu_T = -\beta_T e^{-2\kappa T} + \alpha_T = \int_0^{T-} \sqrt{\kappa} dw_t e^{-2\kappa t}, \quad (188)$$

are those displayed in Equations (148) and (149) for the Harish-Chandra phase-space points. Summarizing, we have that the solutions for the ruler and the Cartan phase-space coordinates are

$$r_T = 2\kappa T, \quad \beta_T = \int_0^{T-} \sqrt{\kappa} dw_t \frac{\cosh 2\kappa t}{\sinh 2\kappa T}, \quad \alpha_T = \int_0^{T-} \sqrt{\kappa} dw_t \frac{\cosh 2\kappa(T-t)}{\sinh 2\kappa T}. \quad (189)$$

The Stratonovich-form SDE for the Harish-Chandra center coordinate z follows from the SDEs of Equations (179) and (180),

$$dz_t = -ds_t + id\psi_t = e^{r_{t+dt/2}} v_{t+dt/2} d\mu_t^*. \quad (190)$$

Converted to Itô form, this equation is

$$\begin{aligned} dz_t &= (v_t + \frac{1}{2}dv_t) e^{r_t} d\mu_t^* \\ &= (v_t + \frac{1}{2}\sqrt{\kappa}dw_t)\sqrt{\kappa}dw_t^* \\ &= \frac{1}{2}\kappa|dw_t|^2 + v_t dw_t^*, \end{aligned} \quad (191)$$

in agreement with Equation (146). The solution for the initial condition $z_0 = 0$ is carefully worked out in Appendix C and given in Equation (150).

The Cartan normalization and phase-space coordinates have Stratonovich-form SDEs,

$$-d\ell = \frac{1}{2}(\beta_1^2 + \beta_2^2)dr + \sinh r(\beta_1 d\alpha_1 + \beta_2 d\alpha_2), \quad (192)$$

$$d\phi = \frac{1}{2}(\beta_1 d\beta_2 - \beta_2 d\beta_1) + \frac{1}{2}(\alpha_1 d\alpha_2 - \alpha_2 d\alpha_1) + \cosh r(\beta_2 d\alpha_1 - \beta_1 d\alpha_2), \quad (193)$$

where all the coefficients are evaluated at the midpoint $t + \frac{1}{2}dt$. Converted to equivalent Itô-form SDEs, these equations become

$$-d\ell_t = |\beta|^2 dr + \sinh r(\beta d\alpha^* + \beta^* d\alpha) + \frac{1}{2} \sinh r(d\beta d\alpha^* + d\beta^* d\alpha) \quad (194)$$

$$= \coth 2\kappa t \kappa |dw_t|^2 - 2|\beta_t|^2 \kappa dt + \beta_t \sqrt{\kappa} dw_t^* + \beta_t^* \sqrt{\kappa} dw_t, \quad (195)$$

$$id\phi = \frac{1}{2}(\beta^* d\beta - \beta d\beta^* + \alpha^* d\alpha - \alpha d\alpha^*) + \cosh r(\beta d\alpha^* - \beta^* d\alpha) + \frac{1}{2} \cosh r(d\beta d\alpha^* - d\beta^* d\alpha) \quad (196)$$

$$= \operatorname{csch} 2\kappa t (\alpha_t \beta_t^* - \alpha_t^* \beta_t) \kappa dt + \frac{1}{2}(\beta_t \coth 2\kappa t - \alpha_t \operatorname{csch} 2\kappa t) \sqrt{\kappa} dw_t^* - \frac{1}{2}(\beta_t^* \coth 2\kappa t - \alpha_t^* \operatorname{csch} 2\kappa t) \sqrt{\kappa} dw_t. \quad (197)$$

Though more complicated, these equations have a similar character to the SDEs for the Harish-Chandra center coordinate $z = -s + i\psi$. The Itô correction for $d\ell$ is the last term in Equation (194), and it becomes the \coth term at the beginning of Equation (195), whereas the Itô correction for $d\phi$ at the end of Equation (196) vanishes. We could solve these equations directly, but it is both easier and more productive to combine the solution for z with the coordinate transformation to Cartan coordinates, thus giving

$$\ell_T = s_T - f_T \quad \text{and} \quad \phi_T = \psi_T - \zeta_T, \quad (198)$$

where s_T and ψ_T are the Harish-Chandra solutions shown in Equations (152) and (153) and f_T and ζ_T are the functions from Equations (162) and (163) with all the coordinates evaluated at time T .

3.5. Solving Most of the FPK Diffusion Equation

Section 2.4 left off showing that the sample paths of SPQM diffuse according to the KOD $D_t(x)$, which satisfies the FPK equation given in Equation (103), with the initial condition $D_0(x) = \delta(x, 1)$. The crucial mathematical object in the diffusion equation is the FPK forward generator Δ , which is written in terms of right-invariant derivatives in Equation (101).

With the frame transformations of Equation (164) at hand, it is easy to express the three pieces of the FPK forward generator in Cartan coordinates,

$$\overleftarrow{\Delta}_{H_0} = -2\partial_r + 2\left(\beta_1 \nabla_1 + \beta_2 \nabla_2 - \frac{\beta_1^2 + \beta_2^2}{2} \partial_\ell\right) + \partial_\phi B_{H_0}, \quad (199)$$

$$\frac{1}{2}\overleftarrow{Q}\overleftarrow{Q} = \frac{1}{2}\left(\nabla_1 - \beta_1\partial_\ell + \frac{\beta_2\cosh r - \alpha_2}{2\sinh r}\partial_\phi\right)^2 \quad (200)$$

$$= \frac{1}{2}(\nabla_1 - \beta_1\partial_\ell)^2 + \frac{1}{2}\partial_\phi B_Q \quad (201)$$

$$= \frac{1}{2}\left(\nabla_1^2 - (\nabla_1\beta_1 + \beta_1\nabla_1)\partial_\ell + \beta_1^2\partial_\ell^2\right) + \frac{1}{2}\partial_\phi B_Q \quad (202)$$

$$= \frac{1}{2}\left(\nabla_1^2 - (\coth r + 2\beta_1\nabla_1)\partial_\ell + \beta_1^2\partial_\ell^2\right) + \frac{1}{2}\partial_\phi B_Q \quad (203)$$

$$= -\frac{1}{2}\coth r\partial_\ell + \frac{1}{2}\nabla_1^2 - \partial_\ell\left(\beta_1\nabla_1 - \frac{1}{2}\beta_1^2\partial_\ell\right) + \frac{1}{2}\partial_\phi B_Q, \quad (204)$$

$$\frac{1}{2}\overleftarrow{P}\overleftarrow{P} = \frac{1}{2}\left(\nabla_2 - \beta_2\partial_\ell - \frac{\beta_1\cosh r - \alpha_1}{2\sinh r}\partial_\phi\right)^2 \quad (205)$$

$$= \frac{1}{2}(\nabla_2 - \beta_2\partial_\ell)^2 + \frac{1}{2}\partial_\phi B_P \quad (206)$$

$$= -\frac{1}{2}\coth r\partial_\ell + \frac{1}{2}\nabla_2^2 - \partial_\ell\left(\beta_2\nabla_2 - \frac{1}{2}\beta_2^2\partial_\ell\right) + \frac{1}{2}\partial_\phi B_P, \quad (207)$$

where we introduce the quantities

$$B_{H_0} = \frac{\beta_2\alpha_1 - \beta_1\alpha_2}{\sinh r}, \quad (208)$$

$$B_Q = \frac{\beta_2\cosh r - \alpha_2}{\sinh r}(\nabla_1 - \beta_1\partial_\ell) + \left(\frac{\beta_2\cosh r - \alpha_2}{2\sinh r}\right)^2\partial_\phi, \quad (209)$$

$$B_P = -\frac{\beta_1\cosh r - \alpha_1}{\sinh r}(\nabla_2 - \beta_2\partial_\ell) + \left(\frac{\beta_1\cosh r - \alpha_1}{2\sinh r}\right)^2\partial_\phi, \quad (210)$$

which are independent of ϕ and commute with ∂_ϕ . The ϕ -derivative terms quickly disappear from the analysis, ultimately because ϕ is irrelevant to the instrument elements as a consequence of the symmetry $\mathcal{O}(e^{i\Omega\phi}L) = \mathcal{O}(L)$. Putting these expressions together, the FPK forward generator of Equation (101) becomes in Cartan coordinates,

$$\Delta = -2\partial_r - \coth r\partial_\ell + \frac{1}{2}(\nabla_1^2 + \nabla_2^2) + (2 - \partial_\ell)\left(\beta_1\nabla_1 + \beta_2\nabla_2 - \frac{\beta_1^2 + \beta_2^2}{2}\partial_\ell\right) + \partial_\phi B, \quad (211)$$

where $B = B_{H_0} + \frac{1}{2}B_Q + \frac{1}{2}B_P$. It is worth noting that

$$\frac{1}{2}(\nabla_1^2 + \nabla_2^2) = \nabla^*\nabla = \nabla\nabla^*, \quad (212)$$

where ∇ and ∇^* , defined in Equations (166) and (167), are derivatives with respect to complex phase-space coordinates.

Due to the cubic and quartic nature of the last few terms in Equation (211), we do not hope to find a complete analytic solution to Equation (103). However, “5/7-ths” of the distribution can be analyzed quite easily. Remember that we are interested in the instrument, and observe that the instrument elements can be partitioned by reconsidering the total operation \mathcal{Z}_T , written in terms of $D_T(x)$ in Equation (87), as

$$\mathcal{Z}_T = \int_G d^7\mu(x) D_T(x) \mathcal{O}(x) \quad (213)$$

$$= \int_{G/Z} d^5\mu(Zx) \left(\int_Z d\phi d\ell D_T(x) e^{-2\ell} \right) \mathcal{O}(D_\beta e^{-H_0 r} D_\alpha^\dagger) \quad (214)$$

$$= \int_{G/Z} d^5\mu(Zx) C_T(Zx) \mathcal{O}(D_\beta e^{-H_0 r} D_\alpha^\dagger), \quad (215)$$

where we use the coset measure $d^5\mu(Zx)$ introduced in Equation (170) and define the *Cartan-section reduced distribution function*,

$$C_T(Zx) \equiv \int_Z d\phi d\ell D_T(x) e^{-2\ell}. \quad (216)$$

Readers uncomfortable with the coset notation can think that, in this equation, $x = e^{i\Omega\phi} e^{-\Omega\ell} D_\beta e^{-H_0 r} D_\alpha^\dagger$ and $Zx = D_\beta e^{-H_0 r} D_\alpha^\dagger$. Even more prosaically, one can regard D_T as being a function of all seven Cartan coordinates and C_T as being a function of five of them, the ruler r and the complex phase-space coordinates β and α . Our excuse—quite a good excuse, really—for using the coordinate-independent coset notation is that we will elaborate on this distribution function and another one in Section 4, but working there mainly in Harish-Chandra coordinates.

Equation (215) is a new unraveling of the SPQM instrument, which we call the *reduced SPQM instrument*, with instrument elements

$$d^5\mu(Zx) C_T(Zx) \mathcal{O}(D_\beta e^{-H_0 r} D_\alpha^\dagger), \quad (217)$$

in which the Cartan reduced distribution $C_T(Zx)$ is conjugate to the operation $\mathcal{O}(D_\beta e^{-H_0 r} D_\alpha^\dagger)$. If H_0 is quantized in the standard way, we have, at late times,

$$D_\beta e^{-H_0 r} D_\alpha^\dagger|_{T \gg 1/\kappa} \sim e^{-\kappa T} D_{\beta_T} |0\rangle\langle 0| D_{\alpha_T}^\dagger = e^{-\kappa T} |\beta_T\rangle\langle \alpha_T|. \quad (218)$$

This makes clear that $C_T(Zx)$ is the elephant in the room hinted at in Section 3.2: it is the phase-space weighting function that is crucial for the POVM completeness relation, which takes the form

$$1_{\mathcal{H}} = \int_G d^7\mu(x) D_T(x) x^\dagger x = \int_{G/Z} d^5\mu(Zx) C_T(Zx) D_\alpha^\dagger e^{-H_0 2r} D_\alpha. \quad (219)$$

As an elephant, however, $C_t(Zx)$ is not normalized to unity—indeed, its normalization is ill defined. Moreover, $C_t(Zx)$ is not the weight function whose moments are those of the Cartan phase-point variables β and α according to the SDE solutions given in Equation (189). The distribution function that does give these moments is the straight marginal of $D_T(x)$ over the center Z ,

$$D_T(Zx) \equiv \int_Z d\phi d\ell D_T(x). \quad (220)$$

This distribution is normalized and has finite moments, those coming from the SDE solutions of Equation (189). For those very reasons, however, $D_T(Zx)$ cannot possibly give rise to a POVM completeness relation; it will not be seen again in this paper.

To derive an evolution equation for the reduced distribution function $C_t(Zx)$ of Equation (216), one takes its time derivative, substitutes omtp tje omtegra; the FPK equation given in Equation (103), and pushes the FPK forward generator Δ from Equation (211) through the center integrals by integrating by parts. Integration by parts on ϕ gets rid of the derivatives with respect to ϕ , and integration by parts on ℓ translates to substituting $\partial_\ell \rightarrow 2$, resulting in the partial differential equation (PDE),

$$\frac{1}{\kappa} \frac{\partial}{\partial t} C_t(Zx) = \left(-2\partial_r - 2\coth r + \nabla^* \nabla \right) [C_t](Zx). \quad (221)$$

for which one should recall that $\nabla^* \nabla = \frac{1}{2}(\nabla_1^2 + \nabla_2^2)$. This PDE is ballistic in the ruler r —solution proportional to $\delta(r - 2\kappa t)$ —and Gaussian-preserving in the phase-space variables, but as a consequence of the $-2\coth r$ term, the PDE does not preserve normalization [15,16,19].

To solve for $C_t(Zx)$ requires knowing $D_{dt}(x)$, which is, when $dt \rightarrow 0$, the δ -function $\delta(x, 1)$. This can be done fairly easily by evaluating the Cartan-coordinate solutions from Equations (189) and (198) at $T = dt$. We perform that task in Appendix F, where we also identify all the δ -function forms for initial conditions. The result for $D_{dt}(x)$ is

$$D_{dt}(x) = \delta(\phi) \delta\left(\ell + \frac{1}{2}\kappa dt |\beta + \alpha|^2\right) \frac{1}{\sinh^2 r} \delta(r - 2\kappa dt) 2\pi \frac{\kappa dt}{\pi} e^{-\kappa dt |\beta + \alpha|^2} 2\pi \delta^2(\beta - \alpha). \quad (222)$$

The distinctive feature of this distribution is the wide, normalized Gaussian in $\beta + \alpha$, which limits to a uniform distribution in $\beta + \alpha$ as $dt \rightarrow 0$. What the wide Gaussian is about is the fact that the identity is represented by Cartan coordinates $\phi = \ell = r = 0$ and $\beta = \alpha$, with $\beta + \alpha$ free to take on any complex value. With this expression, it is easy to see that

$$C_{dt}(Zx) = \int_Z d\phi d\ell D_{dt}(x) e^{-2\ell} = \frac{2}{r} \delta(r - 2\kappa dt) \pi \delta^2(\beta - \alpha). \quad (223)$$

The integral of this distribution has a zero from the $1/r$ behavior multiplying the $\sinh^2 r$ in the measure $d^5\mu(Zx)$ and an infinity from the uniformity in $\beta + \alpha$; therefore, the normalization is ill defined. This ill-defined normalization is, however, exactly what is needed to give a well-defined POVM. For these reasons, $C_{dt}(Zx)$ does not limit to $\delta(Zx, Z1)$ as $dt \rightarrow 0$, for which see Appendix F.

The initial condition $C_{dt}(Zx)$ is independent of $\beta + \alpha$, and the distribution $C_t(Zx)$ remains so under the PDE of Equation (221). To see the consequences most clearly, it is useful to transform to sum and difference variables,

$$\beta_{\pm} = \beta \pm \alpha, \quad (224)$$

in which the covariant derivative ∇ of Equation (166) becomes

$$\nabla = \coth(r/2) \partial_{\beta_+} + \tanh(r/2) \partial_{\beta_-}. \quad (225)$$

Therefore, the weight function evolves according to the PDE,

$$\frac{\partial}{\partial t} C_t(Zx) = \kappa \left(-2\coth r - 2\partial_r + \tanh^2(r/2) \partial_{\beta_-}^* \partial_{\beta_-} \right) C_t(Zx), \quad (226)$$

with solution

$$C_T(Zx) = \frac{1}{\sinh r} \delta(r - 2\kappa T) 2\pi \frac{1}{\pi \Sigma_T} e^{-|\beta - \alpha|^2 / \Sigma_T}. \quad (227)$$

The width of the difference in phase points, Σ_T , satisfies the differential equation, $d\Sigma_t/dt = \kappa \tanh^2 \kappa t$, with a solution, given the initial condition $\Sigma_0 = 0$,

$$\Sigma_T = \kappa T - \tanh \kappa T. \quad (228)$$

In summary, the SPQM instrument can be considered as the *reduced SPQM instrument unraveling*,

$$\mathcal{Z}_T = \int d^5\mu(Zx) C_T(Zx) \mathcal{O}(D_\beta e^{-H_0 r} D_\alpha^\dagger), \quad (229)$$

with instrument elements

$$d^5\mu(Zx) C_T(Zx) \mathcal{O}(D_\beta e^{-H_0 r} D_\alpha^\dagger) = 2 \sinh 2\kappa T dr \delta(r - 2\kappa T) \frac{d^2\alpha}{\pi} \frac{d^2\beta}{\pi \Sigma_T} e^{-|\beta-\alpha|^2/\Sigma_T} \mathcal{O}(D_\beta e^{-H_0 r} D_\alpha^\dagger), \quad (230)$$

where the width Σ_T of the difference in phase points is given by Equation (228). There are four notable features in the temporal behavior of the reduced SPQM instrument:

1. The ruler r (or purity parameter) is ballistic, which means that $e^{-H_0 r}$ collapses exponentially to $e^{-\kappa T} |0\rangle\langle 0|$ in the standard quantization. More generally, $D_\beta e^{-r H_0} D_\alpha^\dagger$ collapses exponentially at late times to an outer product of coherent states, $e^{-\kappa T} |\beta\rangle\langle\alpha|$;
2. The dependence on the future and past phase-space parameters, β and α , is only in their difference;
3. The distribution of the difference spreads out very slowly for small times as $\Sigma_T \propto T^3$ and then for long times becomes normal diffusion, with $\Sigma_T \propto T$;
4. There is a center normalization, $2 \sinh 2\kappa T$, that increases over time.

This center normalization is the focus of the next section, which finds that the elephant provides an alternative perspective on the quantum.

3.6. POVM as an Alternative Perspective on the Quantum

The center normalization just mentioned is remarkable in that it is equivalent to traditional energy quantization. Specifically, the completeness relation for the SPQM process, shown in Equation (219), is

$$1_{\mathcal{H}} = 2 \sinh 2\kappa T \int \frac{d^2\alpha}{\pi} D_\alpha e^{-H_0 4\kappa T} D_\alpha^\dagger \int \frac{d^2\beta}{\pi \Sigma_T} e^{-|\beta-\alpha|^2/\Sigma_T} = 2 \sinh 2\kappa T \int \frac{d^2\alpha}{\pi} D_\alpha e^{-H_0 4\kappa T} D_\alpha^\dagger. \quad (231)$$

It is important to appreciate that, for late times $T \gg 1/\kappa$, when $e^{-H_0 4\kappa T}$ collapses to $e^{-2\kappa T} |0\rangle\langle 0|$ in the standard quantization, this completeness relation becomes the coherent-state resolution of the identity of Equation (155):

$$1_{\mathcal{H}} = \int \frac{d^2\alpha}{\pi} |\alpha\rangle\langle\alpha|. \quad (232)$$

The completeness relation says much more, however, when considered for arbitrary times T . If \mathcal{H} is an irreducible representation, Schur's lemma says that

$$\int \frac{d^2\alpha}{\pi} D_\alpha e^{-H_0 4\kappa T} D_\alpha^\dagger = 1_{\mathcal{H}} \text{tr}(e^{-H_0 4\kappa T}). \quad (233)$$

The trace, which one recognizes as a partition function, is defined within the representation and is evaluated using traditional energy quantization as

$$\text{tr}(e^{-H_0 4\kappa T}) = \sum_{n=0}^{\infty} e^{-(n+\frac{1}{2})4\kappa T} = \frac{1}{2 \sinh 2\kappa T}. \quad (234)$$

The center normalization in the POVM completeness relation thus evaluates the partition function of $e^{-H_0 4\kappa T}$ without using traditional energy quantization. This is to be expected in view of the Stone-von Neumann theorem, but expected though it is, please appreciate

how different the setting of this paper is from the original ideas of energy quantization and thermal equilibrium. Remember that here the operator H_o comes from the trace-preserving character of the instrument. It is not the energy of the system; indeed we have explicitly eschewed any notion of system energy or of a Hamiltonian. The operator H_o plays the role of a “dissipator”—specifically, a dissipator that damps the POVM to the coherent states—but there is no notion of energy associated with this dissipation. The coherent states and H_o arise within a group structure constructed solely from the measured observables, Q and P . Moreover, the parameter conjugate to H_o , the ruler r , is quite literally time rather than an inverse temperature. It seems remarkable that this result holds, from the completeness of the SPQM POVM, without any assumption of a Hamiltonian, a ground state, or thermal equilibrium.

4. Reduced Distribution Functions and Feynman-Kac Path Integrals

This section further considers reduced distribution functions, their path-integral expressions and diffusion equations, and their relation to SDEs. The path-integral expression,

$$D_T(x) = \int \mathcal{D}\mu[dw_{[0,T]}] \delta(x, \gamma[dw_{[0,T]}]), \quad (235)$$

is generally considered to be a Feynman-Kac formula [15,16,19] for the associated diffusion equation. Our analysis is rooted in the path integral of Equation (58) for the overall SPQM quantum operation,

$$\mathcal{Z}_T = \int \mathcal{D}\mu[dw_{[0,T]}] \mathcal{O}(R(\gamma[dw_{[0,T]}])). \quad (236)$$

Equation (236) expresses the relation between a complex Wiener path $dw_{[0,T]}$, a point in the group manifold IWH, $\gamma[dw_{[0,T]}]$, and, in turn, the overall Kraus operator $R(\gamma[dw_{[0,T]}])$ written as a time-ordered product of incremental Kraus operators. Equation (235) defines the Kraus-operator distribution function (KOD) as the amalgamation of all paths that lead to the same Kraus operator. The KOD inherits the path-integral expression given in Equation (236), and from this path integral, one can derive a diffusion equation for the KOD. The reason the path integral is called a Feynman-Kac formula is that everybody after Kac thinks about going in the opposite direction, starting with the diffusion equation and formulating an equivalent path integral.

Section 4.1 reviews the Cartan-section reduced distribution $C_T(Zx)$, introduces the Harish-Chandra-section reduced distribution $B_T(Zx)$ of Equation (24), and shows that these two are related by a positive gauge transformation. Section 4.2 defines a normalized version of the Harish-Chandra reduced distribution, denoted by $\tilde{B}_T(Zx)$, and finds its path-integral expression in terms of a modified path-integration measure in which the outcome increments are correlated. Section 4.3 formulates the diffusion equations for $B_T(Zx)$ and $\tilde{B}_T(Zx)$, and Section 4.4 solves the path integral for $\tilde{B}_T(Zx)$ using the stochastic integrals for the Harish-Chandra phase-space coordinates.

4.1. Feynman-Kac Formulas

The *ur* KOD of Equation (235),

$$D_T(x) \equiv \int \mathcal{D}\mu[dw_{[0,T]}] \delta(x, \gamma[dw_{[0,T]}]), \quad (237)$$

unravels \mathcal{Z}_T over the universal domain of $G = \text{IWH}$,

$$\mathcal{Z}_T = \int_G d^7\mu(x) D_T(x) \mathcal{O}(x), \quad (238)$$

Reduced distributions are defined on $G/Z = \text{RIWH}$. The first of these reduced distributions, introduced in Equation (216) as the Cartan-section reduced distribution,

$$C_T(Zx) \equiv \int_Z d\phi d\ell D_T(x) e^{-2\ell}, \quad (239)$$

can also be defined by the Feynman-Kac path integral,

$$C_T(Zx) = \int \mathcal{D}\mu[dw_{[0,T)}] e^{-2\ell[dw_{[0,T)}]} \delta\left(Zx, Z\gamma[dw_{[0,T)}]\right), \quad (240)$$

where $\ell[dw_{[0,T)}] = \ell_T$ is the solution of the SDE of Equation (195) for the Cartan center coordinate ℓ ,

$$-\ell[dw_{[0,T)}] = \int_0^{T-} \left((\coth 2\kappa t - 2|\beta_t|^2) \kappa dt + \beta_t \sqrt{\kappa} dw_t^* + \beta_t^* \sqrt{\kappa} dw_t \right), \quad (241)$$

with the notation here emphasizing that this solution is a functional of the sample path of Wiener outcome increments. As was noted in Equation (198), one can use the transformation from Harish-Chandra coordinates to write

$$-\ell[dw_{[0,T)}] = f_T - s[dw_{[0,T)}], \quad (242)$$

where f_T is the function of $G/Z = \text{RIWH}$ given in Equation (162),

$$2f(Zx) = e^{-r/2} \sinh(r/2) |\beta + \alpha|^2 + e^{-r/2} \cosh(r/2) |\beta - \alpha|^2, \quad (243)$$

with the ruler and the phase-plane coordinates evaluated at time T , and

$$-2s_T = -2s[dw_{[0,T)}] = \int_0^{T-} \int_0^{T-} \kappa dw_t^* dw_s e^{-2\kappa|t-s|} \quad (244)$$

is the stochastic integral for the Harish-Chandra center coordinate s , given in Equation (152) and derived in Appendix C. The function $C_T(Zx)$ unravels \mathcal{Z}_T onto $\text{RIWH} = G/Z$ as in Equation (215),

$$\mathcal{Z}_T = \int_{G/Z} d^5\mu(Zx) C_T(Zx) \mathcal{O}\left(D_\beta e^{-H_0 r} D_\alpha^\dagger\right), \quad (245)$$

and $C_T(Zx)$ satisfies the FPK equation displayed in Equation (221),

$$\frac{1}{\kappa} \frac{\partial}{\partial t} C_t(Zx) = \left(-2\partial_r - 2\coth r + \nabla^* \nabla \right) C_t(Zx). \quad (246)$$

There is another natural reduced distribution, conjugate to the Harish-Chandra section,

$$B_T(Zx) \equiv \int_Z d\psi ds D_T(x) e^{-2s} \quad (247)$$

$$= \int \mathcal{D}\mu[dw_{[0,T)}] e^{-2s[dw_{[0,T)}]} \delta\left(Zx, Z\gamma[dw_{[0,T)}]\right). \quad (248)$$

This Harish-Chandra reduced distribution function, $B_T(Zx)$, unravels \mathcal{Z}_T as

$$\mathcal{Z}_T = \int_{G/Z} d^5\mu(Zx) B_T(Zx) \mathcal{O}\left(e^{a^\dagger v} e^{-H_0 r} e^{a\mu^*}\right). \quad (249)$$

The two reduced distribution functions are equivalent to one another through a positive gauge transformation [19],

$$C_T(Zx) = \int_Z d\phi d\ell D_T(x) e^{-2\ell} \quad (250)$$

$$= \int_Z d\psi ds D_T(x) e^{-2[s-f(Zx)]} \quad (251)$$

$$= e^{2f(Zx)} B_T(Zx). \quad (252)$$

In the lingo of Feynman-Kac formulas [19], $2f(Zx)$ is the “convective pressure”.

Section 3.5 introduced the Cartan reduced distribution $C_T(Zx)$ and showed that it is the distribution that expresses POVM completeness. The price for relevance to POVM completeness is that $C_T(Zx)$ has ill-defined normalization and is disconnected from the stochastic-integral solutions for the phase-space variables. The next three sections investigate the Harish-Chandra reduced distribution $B_T(Zx)$. It is clear that $B_T(Zx)$ is not the right distribution for addressing POVM completeness because of the Gaussian gauge function $e^{-2f(Zx)}$, but this Gaussian gauge transformation is just what is needed to get a Gaussian distribution function that can be normalized and whose normalized version can be evaluated from the moments of the phase-space variables, albeit, as we shall see, moments defined relative to a modified path-integration measure.

4.2. Normalized Harish-Chandra Reduced Distribution Function and Modified Path-Integration Measure

The Feynman-Kac formula for $B_T(Zx)$, shown in Equation (248), suggests that we combine $e^{-2s[dw_{[0,T]}]}$ with the Weiner measure of Equation (56),

$$\mathcal{D}\mu[dw_{[0,T]}] e^{-2s[dw_{[0,T]}]} = \left(\prod_{k=0}^{T/dt-1} d^2(dw_{kdt}) \right) \left(\frac{1}{\pi dt} \right)^{T/dt} \exp \left(- \int_0^{T-} \frac{|dw_t|^2}{dt} + \int_0^{T-} \int_0^{T-} \kappa dw_t^* dw_s e^{-2\kappa|t-s|} \right). \quad (253)$$

The quadratic functional in the exponential can be written as

$$- \int_0^{T-} \frac{|dw_t|^2}{dt} + \int_0^{T-} \int_0^{T-} \kappa dw_t^* dw_s e^{-2\kappa|t-s|} = - \frac{1}{dt} \int_{t=0}^{T-} \int_{s=0}^{T-} dw_t^* dw_s \left(\delta_{ts} - \kappa dt e^{-2\kappa|t-s|} \right) \quad (254)$$

$$= - \frac{1}{dt} \sum_{k=0}^{N-1} \sum_{l=0}^{N-1} dw_k^* dw_l \left(\delta_{kl} - \kappa dt e^{-2\kappa dt |k-l|} \right). \quad (255)$$

When converting between stochastic integrals and sums, we use $t_k = k dt$ ($t_N = T = N dt$) and $dw_k = dw_{kdt} = dw_{t_k}$.

Now we define the real, symmetric, and positive $N \times N$ matrix M_T , whose matrix elements are

$$M_{kl} = \delta_{kl} - \kappa dt e^{-2\kappa dt |k-l|}, \quad (256)$$

It is elegant for various formal expressions to introduce a continuous version of M_T , but we do not bother with that here since we work with the sums that the stochastic integrals represent. Notice that M_T is a Toeplitz matrix; that is, $M_{k+j,l+j} = M_{kl}$. Putting this together, we have

$$\mathcal{D}\mu[dw_{[0,T]}] e^{-2s[dw_{[0,T]}]} = \left(\prod_{k=0}^{N-1} d^2(dw_{kdt}) \right) \left(\frac{1}{\pi dt} \right)^N \exp \left(- \frac{1}{dt} \sum_{k=0}^{N-1} \sum_{l=0}^{N-1} dw_k^* M_{kl} dw_l \right), \quad (257)$$

which integrates over the Wiener outcome paths to

$$\int \mathcal{D}\mu[dw_{[0,T)}] e^{-2s[dw_{[0,T)}]} = \left(\frac{1}{\pi dt}\right)^N \frac{\pi^N}{\det(M_T/dt)} = \frac{1}{\det M_T}. \quad (258)$$

This prompts us to define a new (normalized, zero-mean) Gaussian measure on the Wiener outcome paths,

$$\mathcal{D}\mu_M[dw_{[0,T)}] \equiv \det M_T \mathcal{D}\mu[dw_{[0,T)}] e^{-2s[dw_{[0,T)}]} \quad (259)$$

$$= \left(\prod_{k=0}^{N-1} d^2(dw_k)\right) \left(\frac{\det M_T}{(\pi dt)^N}\right) \exp\left(-\frac{1}{dt} \sum_{k=0}^{N-1} \sum_{l=0}^{N-1} dw_k^* M_{kl} dw_l\right). \quad (260)$$

Relative to this measure, the outcome increments are correlated,

$$\langle dw_k^* dw_l \rangle_M = dt (M_T^{-1})_{kl}. \quad (261)$$

Notice that the increment correlation $\langle dw_k^* dw_l \rangle_M$, with k and l held fixed, changes as the total time T changes. The inverse matrix M_T^{-1} matrix is real, symmetric, and positive, all properties inherited from M_T . The inverse does not inherit the Toeplitz property of M_T . That M_T is Toeplitz implies that it is persymmetric, that is, symmetric across the anti-diagonal. M_T^{-1} does inherit the persymmetry, which turns out to have an important consequence in Section 4.4.

Returning to the Harish-Chandra reduced distribution of Equation (247), we see that its normalization can be written in several ways:

$$\mathcal{N}_T \equiv \int d^5\mu(Zx) B_T(Zx) \quad (262)$$

$$= \int d^5\mu(Zx) C_T(Zx) e^{-2f(Zx)} \quad (263)$$

$$= \int_Z d^7\mu(x) D_T(x) e^{-2s} \quad (264)$$

$$= \int \mathcal{D}\mu[dw_{[0,T)}] e^{-2s[dw_{[0,T)}]} = \frac{1}{\det M_T}. \quad (265)$$

The normalized version of the Harish-Chandra reduced distribution is

$$\tilde{B}_T(Zx) \equiv \frac{1}{\mathcal{N}_T} B_T(Zx) = \frac{e^{-2f(Zx)}}{\mathcal{N}_T} C_T(Zx). \quad (266)$$

Section 4.3 finds the diffusion equations satisfied by the reduced Harish-Chandra distributions, and Section 4.4 uses the path integral for the normalized distribution,

$$\tilde{B}_T(Zx) = \int \mathcal{D}\mu_M[dw_{[0,T)}] \delta\left(Zx, Z\gamma[dw_{[0,T)}]\right), \quad (267)$$

to evaluate $\tilde{B}_T(Zx)$ from the SDE solutions for the Harish-Chandra coordinates.

4.3. Diffusion Equation for Harish-Chandra Reduced Distribution Function

The easiest way to get to the diffusion equation for $B_t(Zx)$ is to return to the FPK equation for $D_t(x)$ given in Equation (103), write the FPK forward generator Δ in Harish-Chandra coordinates, and then marginalize over the center to get a PDE for $B_t(Zx)$.

By using the frame transformations shown in Equation (164), it is easy to express the terms of the forward generator in Harish-Chandra coordinates,

$$2H_{\leftarrow} = -2\partial_r + 2(\nu_1\partial_{\nu_1} + \nu_2\partial_{\nu_2}) \quad (268)$$

$$= -2\partial_r - 4 + 2(\partial_{\nu_1}\nu_1 + \partial_{\nu_2}\nu_2), \quad (269)$$

$$\frac{1}{2}\overleftarrow{Q}\overleftarrow{Q} = \frac{1}{2}\left(\nabla_1 - \frac{1}{2}\nu_1\partial_s + \frac{1}{2}\nu_2\partial_\psi\right)^2 \quad (270)$$

$$= \frac{1}{2}(\nabla_1 - \frac{1}{2}\nu_1\partial_s)^2 + \frac{1}{2}\partial_\psi A_Q \quad (271)$$

$$= \frac{1}{2}\left(\nabla_1^2 - \frac{1}{2}\partial_s(\nabla_1\nu_1 + \nu_1\nabla_1) + \frac{1}{4}\nu_1^2\partial_s^2\right) + \frac{1}{2}\partial_\psi A_Q \quad (272)$$

$$= \frac{1}{2}\left(\nabla_1^2 - \frac{1}{2}\partial_s(2\nabla_1\nu_1 - 1) + \frac{1}{4}\nu_1^2\partial_s^2\right) + \frac{1}{2}\partial_\psi A_Q \quad (273)$$

$$= \frac{1}{4}\partial_s\left(1 + \frac{1}{2}\nu_1^2\partial_s\right) + \frac{1}{2}\nabla_1^2 - \frac{1}{2}\partial_s\nabla_1\nu_1 + \frac{1}{2}\partial_\psi A_Q, \quad (274)$$

$$\frac{1}{2}\overleftarrow{P}\overleftarrow{P} = \frac{1}{2}\left(\nabla_2 - \frac{1}{2}\nu_2\partial_s - \frac{1}{2}\nu_1\partial_\psi\right)^2 \quad (275)$$

$$= \frac{1}{2}(\nabla_2 - \frac{1}{2}\nu_2\partial_s)^2 + \frac{1}{2}\partial_\psi A_P \quad (276)$$

$$= \frac{1}{4}\partial_s\left(1 + \frac{1}{2}\nu_2^2\partial_s\right) + \frac{1}{2}\nabla_2^2 - \frac{1}{2}\partial_s\nabla_2\nu_2 + \frac{1}{2}\partial_\psi A_P, \quad (277)$$

where the terms

$$A_Q = \nu_2(\nabla_1 - \frac{1}{2}\nu_1\partial_s) + \frac{1}{4}\nu_2^2\partial_\psi, \quad (278)$$

$$A_P = -\nu_1(\nabla_2 - \frac{1}{2}\nu_2\partial_s) - \frac{1}{4}\nu_1^2\partial_\psi, \quad (279)$$

are independent of ψ and commute with ∂_ψ . Putting these expressions together, the FPK forward generator of Equation (101) becomes in Harish-Chandra coordinates,

$$\Delta = -2\partial_r - 4 + \frac{1}{2}\partial_s\left(1 + \frac{\nu_1^2 + \nu_2^2}{4}\partial_s\right) + \frac{1}{2}(\nabla_1^2 + \nabla_2^2) + 2\partial_{\nu_1}\nu_1 + 2\partial_{\nu_2}\nu_2 - \frac{1}{2}\partial_s(\nabla_1\nu_1 + \nabla_2\nu_2) + \partial_\psi A, \quad (280)$$

where $A = \frac{1}{2}(A_Q + A_P)$.

To derive an evolution equation for $B_t(Zx)$ from Equation (247), one follows the procedure outlined for $C_t(Zx)$ in Equation (221), using here the rules that integration by parts on s and ψ makes the substitutions $\partial_s \rightarrow 2$ and $\partial_\psi \rightarrow 0$. The resulting PDE for $B_t(Zx)$ is

$$\frac{1}{\kappa}\frac{\partial}{\partial t}B_t(Zx) = \left(-2\partial_r - 3 + \frac{\nu_1^2 + \nu_2^2}{2} + \tilde{\Delta}\right)B_t(Zx), \quad (281)$$

where, for brevity, we define a reduced generator for the phase-space-variable derivatives,

$$\tilde{\Delta} \equiv (2\partial_{\nu_1} - \nabla_1)\nu_1 + (2\partial_{\nu_2} - \nabla_2)\nu_2 + \frac{1}{2}(\nabla_1^2 + \nabla_2^2) \quad (282)$$

$$= (2\partial_\nu - \nabla)\nu + (2\partial_{\nu^*} - \nabla^*)\nu^* + \nabla^*\nabla. \quad (283)$$

Notice that

$$2\partial_{v_j} - \nabla_j = \partial_{v_j} - e^{-r} \partial_{\mu_j}, \quad (284)$$

$$2\partial_v - \nabla = \partial_v - e^{-r} \partial_\mu. \quad (285)$$

Converting fully to complex phase-space coordinates puts the PDE in the form

$$\frac{1}{\kappa} \frac{\partial}{\partial t} B_t(Zx) = \left(-2\partial_r - 3 + |v|^2 + \tilde{\Delta} \right) B_t(Zx). \quad (286)$$

This PDE is ballistic in the ruler r —solution proportional to $\delta(r - 2\kappa t)$ —and Gaussian-preserving in the phase-space variables, but as a consequence of the term $-3 + |v|^2$, it does not preserve normalization [15,16,19]. The effect of the positive gauge transformation from $C_t(Zx)$ to $B_t(Zx)$ is twofold: (i) the norm-nonconserving “potential” term changes character, from a ruler-dependent $-2 \coth r$ in the PDE for $C_t(Zx)$ in Equation (246) to a term $-3 + |v|^2$ in the PDE for $B_t(Zx)$ in Equation (286), which depends on the posterior phase-space variable v ; (ii) there are first-derivative, “vector-potential” terms in the PDE for $B_t(Zx)$, corresponding to the Ornstein-Uhlenbeck behavior of v in Equation (148), whereas there are no such terms in the PDE for $C_T(Zx)$.

We turn now to the task of converting the PDE for $B_t(Zx)$ in Equation (286) to a PDE for the normalized distribution $\tilde{B}_t(Zx)$. The initial condition for the PDE in Equation (286) comes from inserting the Harish-Chandra $D_{dt}(x)$ of Equation (A140) into $B_{dt}(Zx)$ as it is expressed in the integral in Equation (247) specialized to $T = dt$. Taking the limit $dt \rightarrow 0$, one finds the expected result that $B_0(Zx)$ is the δ -function on G/Z of Equation (A143):

$$B_0(Zx) = \delta(Zx, Z1) = \delta(r) 2\pi \delta^2(v) 2\pi \delta^2(\mu). \quad (287)$$

This is expected because the identity is represented uniquely in Harish-Chandra coordinates by $\psi = s = r = 0, v = \mu = 0$. The initial condition means that $B_t(Zx)$ is initially normalized to unity; thus the normalization factor \mathcal{N}_t has the initial value $\mathcal{N}_0 = 1$, and the normalized distribution has the same initial condition,

$$\tilde{B}_0(Zx) = B_0(Zx) = \delta(Zx, Z1). \quad (288)$$

The normalization factor,

$$\mathcal{N}_T \equiv \int d^5\mu(Zx) B_T(Zx) = \int e^{2r} dr \frac{d^2v}{2\pi} \frac{d^2\mu}{2\pi} B_T(Zx), \quad (289)$$

satisfies the differential equation,

$$\frac{1}{\kappa} \frac{d\mathcal{N}_t}{dt} = \int e^{2r} dr \frac{d^2v}{2\pi} \frac{d^2\mu}{2\pi} \left(-2\partial_r - 3 + |v|^2 \right) B_t(Zx) \quad (290)$$

$$= \mathcal{N}_t \int d^5\mu(Zx) (1 + |v|^2) \tilde{B}_t(Zx), \quad (291)$$

where the reader should notice that integration by parts on the ruler becomes the rule $\partial_r \rightarrow -2$. This differential equation assumes the form

$$\frac{1}{\kappa} \frac{d \ln \mathcal{N}_t}{dt} = 1 + n_t, \quad (292)$$

where

$$n_t \equiv \langle |v_t|^2 \rangle_M = \int d^5\mu(Zx) |v|^2 \tilde{B}_t(Zx) \quad (293)$$

is the second moment of ν relative to the normalized distribution $\tilde{B}_t(Zx)$. We place a subscript M on this moment because we can use the path-integral expression for $\tilde{B}_T(Zx)$ given in Equation (267) to re-express the moment as

$$n_T = \langle |\nu_T|^2 \rangle_M = \int \mathcal{D}\mu_M[dw_{[0,T]}] |\nu[dw_{[0,T]}]|^2, \quad (294)$$

where $\nu[dw_{[0,T]}] = \nu_T$ is the stochastic-integral solution for ν , given in Equation (148). Once the stochastic integral is plugged into this equation, the correlations of the Wiener increments are evaluated according to the modified measure $\mathcal{D}\mu_M[dw_{[0,T]}]$, that is, as in Equation (261).

The PDE for the normalized reduced distribution $\tilde{B}_t(Zx)$ of Equation (266) now follows as

$$\frac{1}{\kappa} \frac{\partial}{\partial t} \tilde{B}_t(Zx) = \left(-\frac{1}{\kappa} \frac{d \ln \mathcal{N}_t}{dt} - 2\partial_r - 3 + |\nu|^2 + \tilde{\Delta} \right) \tilde{B}_t(Zx). \quad (295)$$

Inserting Equation (292) gives

$$\frac{1}{\kappa} \frac{\partial}{\partial t} \tilde{B}_t(Zx) = \left(-2\partial_r - 4 + |\nu|^2 - n_t + \tilde{\Delta} \right) \tilde{B}_t(Zx). \quad (296)$$

It is easy to see how this equation preserves normalization. The presence of the moment n_t , essential for normalization, makes the equation nonlinear, but it can still be solved easily.

To solve for $\tilde{B}_t(Zx)$, one notes that the PDE is ballistic in the ruler r and Gaussian-preserving in the phase-space variables. Thus, the solution has the form $\tilde{B}_t(Zx) = e^{-2r} \delta(r - 2\kappa t) \Phi_t(Zx)$, where $\Phi_t(Zx)$ is a normalized, zero-mean Gaussian in the phase-space variables ν and μ . The derivatives in the PDE of Equation (296) are invariant under complex conjugation and under simultaneous rephasing of the phase-space variables, that is, $\nu \rightarrow \nu e^{i\chi}$ and $\mu \rightarrow \mu e^{i\chi}$. It is productive to think of the invariance under complex conjugation as invariance under the change of phase-space coordinates $\nu \leftrightarrow \nu^*$ and $\mu \leftrightarrow \mu^*$. It is useful to appreciate that all the diffusion equations in this paper share these invariance properties. The invariance under simultaneous rephasing implies that $\Phi_t(Zx)$ is a zero-mean Gaussian since it starts from a zero-mean δ -function initial condition. It further implies that $\Phi_t(Zx)$ is determined by the three nonzero second moments of the phase-space variables: n_t of Equation (293) and

$$m_t \equiv \langle |\mu_t|^2 \rangle_M = \int d^5\mu(Zx) |\mu|^2 \tilde{B}_t(Zx), \quad (297)$$

$$q_t \equiv \langle \nu_t^* \mu_t \rangle_M = \langle \mu_t^* \nu_t \rangle_M = \int d^5\mu(Zx) \nu^* \mu \tilde{B}_t(Zx), \quad (298)$$

with the invariance under complex conjugation implying that q_t is real. This form of the solution for $\tilde{B}_t(Zx)$ in hand, one derives from the PDE of Equation (296) first-order temporal differential equations for the second moments n_t , m_t , and q_t . Just as the ordinary differential equation (ODE) for the normalization factor involves a second moment, so the equations for the second moments involve fourth moments. The Gaussian form of the solution relates the fourth moments to second moments, thus closing the system of differential equations. The resulting three ODEs, for the derivatives of n_t , m_t , and q_t , have terms that are constant, linear, and quadratic in the moments; the presence of the quadratic terms makes these ODEs (coupled) Riccati equations. The last step is to solve the three Riccati equations with initial conditions $n_0 = m_0 = q_0 = 0$, which are implied by the δ -function initial condition for $\tilde{B}_0(Zx)$. This gives the solution for $\tilde{B}_t(Zx)$. By integrating to find \mathcal{N}_T , using the solution for n_t , one can backtrack to find $B_T(Zx)$ and $C_T(Zx)$ from Equation (266).

We have carried out this procedure of deriving the Riccati equations from the PDE in Equation (296), but we do not present that derivation in this paper, preferring instead to use a different method, which derives the Riccati equations from the path integral for $\tilde{B}_T(Zx)$. Implementing that method is the final task of this paper, carried out in the next section.

4.4. Normalized Harish-Chandra Reduced Distribution Function from Its Path Integral

The path integral for $\tilde{B}_T(Zx)$, displayed in Equation (267), can be written explicitly in terms of the δ -function in Harish-Chandra coordinates,

$$\tilde{B}_T(Zx) = \int \mathcal{D}\mu_M[dw_{[0,T]}] e^{-2r} \delta(r - r_T) 2\pi\delta^2(v - v_T) 2\pi\delta^2(\mu - \mu_T) \quad (299)$$

$$= e^{-4\kappa T} \delta(r - 2\kappa T) \int \mathcal{D}\mu_M[dw_{[0,T]}] 2\pi\delta^2(v - v[dw_{[0,T]}]) 2\pi\delta^2(\mu - \mu[dw_{[0,T]}]) , \quad (300)$$

where $r_T = 2\kappa T$, $v_T = v[dw_{[0,T]}]$, and $\mu_T = \mu[dw_{[0,T]}]$ are the solutions to the SDEs for the ruler and the Harish-Chandra phase-space coordinates, shown in Equations (147)–(149). The measure for the path integral is a (normalized, zero-mean) Gaussian measure in the outcome increments dw_t ; thus, the path integral gives a normalized, zero-mean Gaussian in v and μ , which is determined by the three (real) moments introduced in the previous section:

$$n_T = \langle |v_T|^2 \rangle_M = \int d^5\mu(Zx) |v|^2 \tilde{B}_T(Zx) = \int \mathcal{D}\mu_M[dw_{[0,T]}] |v[dw_{[0,T]}]|^2 , \quad (301)$$

$$m_T = \langle |\mu_T|^2 \rangle_M = \int d^5\mu(Zx) |\mu|^2 \tilde{B}_T(Zx) = \int \mathcal{D}\mu_M[dw_{[0,T]}] |\mu[dw_{[0,T]}]|^2 , \quad (302)$$

$$q_T = \langle v_T^* \mu_T \rangle_M = \langle \mu_T^* v_T \rangle = \int d^5\mu(Zx) v^* \mu \tilde{B}_T(Zx) = \int \mathcal{D}\mu_M[dw_{[0,T]}] v[dw_{[0,T]}]^* \mu[dw_{[0,T]}] . \quad (303)$$

In this context, that the first moments and all the other second moments of the phase-space variables are zero follows from the fact that the measure is invariant under simultaneous rephasing of all the outcome increments; the reality of q_T follows from the fact that the measure is unchanged under the transformation $dw_{[0,T]} \rightarrow dw_{[0,T]}^*$. These properties come from the fact that M_T is real and symmetric.

Plugging in the stochastic-integral solutions for v_T and μ_T puts these moments into the following form:

$$n_T = \langle |v_T|^2 \rangle_M = \kappa \sum_{k,l=0}^{N-1} \langle dw_k^* dw_l \rangle_M e^{-2\kappa(T-t_k)} e^{-2\kappa(T-t_l)} \quad (304)$$

$$= \kappa dt \sum_{k,l=0}^{N-1} e^{-2\kappa dt(N-k)} (M_T^{-1})_{kl} e^{-2\kappa dt(N-l)} , \quad (305)$$

$$m_T = \langle |\mu_T|^2 \rangle_M = \kappa \sum_{k,l=0}^{N-1} \langle dw_k^* dw_l \rangle_M e^{-2\kappa t_k} e^{-2\kappa t_l} \quad (306)$$

$$= \kappa dt \sum_{k,l=0}^{N-1} e^{-2\kappa dt k} (M_T^{-1})_{kl} e^{-2\kappa dt l} , \quad (307)$$

$$q_T = \langle v_T^* \mu_T \rangle_M = \langle \mu_T^* v_T \rangle_M = \kappa \sum_{k,l=0}^{N-1} \langle dw_k^* dw_l \rangle_M e^{-2\kappa(T-t_k)} e^{-2\kappa t_l} \quad (308)$$

$$= \kappa dt \sum_{k,l=0}^{N-1} e^{-2\kappa dt(N-k)} (M_T^{-1})_{kl} e^{-2\kappa dt l} . \quad (309)$$

In the final form of q_T , it is evident that q_T is real. Notice that these expressions satisfy the zero initial conditions.

That M_T is a Toeplitz matrix introduces an additional, quite important symmetry. The inverse matrix does not inherit the Toeplitz property of M , but it does inherit a less restrictive property. That M_T is Toeplitz implies that it is symmetric about the anti-diagonal; that is, $M_{kl} = M_{N-1-l, N-1-k}$. A matrix that is symmetric about the anti-diagonal is called *persymmetric*. It is easy to show that the inverse of a persymmetric matrix is persymmetric, so M_T^{-1} satisfies $(M_T^{-1})_{kl} = (M_T^{-1})_{N-1-l, N-1-k}$. Persymmetry has a major consequence for the three moments, which comes from manipulating m_T :

$$m_T = \kappa dt \sum_{k,l=0}^{N-1} e^{-2\kappa dt (N-1-k)} (M_T^{-1})_{N-1-k, N-1-l} e^{-2\kappa dt (N-1-l)} \quad (310)$$

$$= e^{4\kappa dt} \kappa dt \sum_{k,l=0}^{N-1} e^{-2\kappa dt (N-k)} (M_T^{-1})_{lk} e^{-2\kappa dt (N-l)}. \quad (311)$$

We can set $e^{4\kappa dt} = 1$ and thus conclude that

$$m_T = n_T. \quad (312)$$

The complementarity in time of ν_T and μ_T was discussed in Section 3.2: the post-measurement variable ν_T of Equation (148) depends exponentially on the end of the outcome register, and the POVM variable μ_T of Equation (149) depends exponentially on the beginning of the register. The persymmetry of M_T and M_T^{-1} expresses that the beginning and end of the record look the same statistically, so it is not surprising that the persymmetry implies that $m_T = n_T$.

The next step, deriving Riccati equations for the three moments, involves incrementing the moments from T to $T + dT$. The tedious part of this task is determining how M_T^{-1} increments—that is, finding M_{T+dT}^{-1} —and that can be done using the Schur complement. We relegate this entire task to Appendix G and here skip directly to the coupled Riccati ODEs, taken from Equations (A174)–(A176):

$$\frac{1}{\kappa} \frac{dn_T}{dT} = (1 - n_T)^2, \quad (313)$$

$$\frac{1}{\kappa} \frac{dm_T}{dT} = (q_T + e^{-2\kappa T})^2, \quad (314)$$

$$\frac{1}{\kappa} \frac{dq_T}{dT} = -q_T(1 - n_T) + e^{-2\kappa T}(1 + n_T). \quad (315)$$

With the zero initial conditions, these have the solutions,

$$n_T = m_T = \frac{\kappa T}{1 + \kappa T}, \quad (316)$$

$$q_T = \frac{1}{1 + \kappa T} - e^{-2\kappa T} = -\frac{\kappa T}{1 + \kappa T} + 2e^{-\kappa T} \sinh \kappa T. \quad (317)$$

It is quite instructive to notice that the equality $n_T = m_T$ and the reality of q_T together imply that the sum and difference Harish-Chandra phase-space variables are uncorrelated,

$$\langle (\nu_T \pm \mu_T)^* (\nu_T \mp \mu_T) \rangle_M = \langle |\nu_T|^2 \rangle_M - \langle |\mu_T|^2 \rangle_M \mp \langle \nu_T^* \mu_T \rangle_M \pm \langle \mu_T^* \nu_T \rangle_M = 0, \quad (318)$$

with second moments,

$$\langle |\nu_T \pm \mu_T|^2 \rangle_M = \langle |\nu_T|^2 \rangle_M + \langle |\mu_T|^2 \rangle_M \pm \langle \nu_T^* \mu_T \rangle_M \pm \langle \mu_T^* \nu_T \rangle_M = 2(n_T \pm q_T), \quad (319)$$

where

$$n_T + q_T = 2e^{-\kappa T} \sinh \kappa T, \quad (320)$$

$$n_T - q_T = \frac{-1 + \kappa T}{1 + \kappa T} + e^{-2\kappa T} = 2e^{-\kappa T} \cosh \kappa T \frac{\Sigma_T}{1 + \kappa T}, \quad (321)$$

The width $\Sigma_T = \kappa T - \tanh \kappa T$ was introduced in Equation (228). It is worth noting the early- and late-time behavior of the various moments:

$$n_T = m_T = \begin{cases} \kappa T, & \kappa T \ll 1, \\ 1, & \kappa T \gg 1, \end{cases} \quad (322)$$

$$q_T = \begin{cases} \kappa T, & \kappa T \ll 1, \\ 1/\kappa T, & \kappa T \gg 1, \end{cases} \quad (323)$$

$$n_T + q_T = \begin{cases} 2\kappa T, & \kappa T \ll 1, \\ 1, & \kappa T \gg 1, \end{cases} \quad (324)$$

$$n_T - q_T = \begin{cases} \frac{2}{3}(\kappa T)^3, & \kappa T \ll 1, \\ 1, & \kappa T \gg 1. \end{cases} \quad (325)$$

At early times, ν and μ are tightly correlated, with their sum undergoing standard diffusion; at late times, they become uncorrelated, and their second moments saturate at 1.

The sum and difference variables being uncorrelated, the Gaussian path integral from Equation (299) has the normalized solution,

$$\tilde{B}_T(Zx) = e^{-4\kappa T} \delta(r - 2\kappa T) 4\pi \frac{1}{2\pi(n_T + q_T)} \exp\left(-\frac{|\nu + \mu|^2}{2(n_T + q_T)}\right) 4\pi \frac{1}{2\pi(n_T - q_T)} \exp\left(-\frac{|\nu - \mu|^2}{2(n_T - q_T)}\right). \quad (326)$$

Noting that

$$(n_T + q_T)(n_T - q_T) = \sinh 2\kappa T \frac{2\Sigma_T}{e^{2\kappa T}(1 + \kappa T)}, \quad (327)$$

we can put the normalized solution in the form

$$\tilde{B}_T(Zx) = \frac{1}{\sinh 2\kappa T} \delta(r - 2\kappa T) \frac{2e^{-2\kappa T}(1 + \kappa T)}{\Sigma_T} \exp\left(-|\nu + \mu|^2 \frac{e^{\kappa T}}{4 \sinh \kappa T} - |\nu - \mu|^2 \frac{e^{\kappa T}}{4 \cosh \kappa T} \frac{1 + \kappa T}{\Sigma_T}\right), \quad (328)$$

which satisfies the δ -function initial condition of Equation (288).

To retrieve the unnormalized distribution $B_T(Zx) = \mathcal{N}_T \tilde{B}_T(Zx)$, one needs $\det M_T = 1/\mathcal{N}_T$, which, according to Equation (A169) or Equation (292), satisfies the equation

$$\frac{1}{\kappa} \frac{d \ln \det M_T}{dT} = -(1 + n_T) = -2 + \frac{1}{1 + \kappa T}, \quad (329)$$

with the solution, for the initial condition $\det M_0 = 1$,

$$\det M_T = e^{-2\kappa T}(1 + \kappa T) = \frac{1}{\mathcal{N}_T}. \quad (330)$$

The unnormalized distribution is therefore

$$B_T(Zx) = \frac{1}{\sinh 2\kappa T} \delta(r - 2\kappa T) \frac{2}{\Sigma_T} \exp\left(-|\nu + \mu|^2 \frac{e^{\kappa T}}{4 \sinh \kappa T} - |\nu - \mu|^2 \frac{e^{\kappa T}}{4 \cosh \kappa T} \frac{1 + \kappa T}{\Sigma_T}\right) \quad (331)$$

$$= \frac{1}{\sinh 2\kappa T} \delta(r - 2\kappa T) \frac{2}{\Sigma_T} \exp\left(-|\beta + \alpha|^2 e^{-\kappa T} \sinh \kappa T - |\beta - \alpha|^2 e^{-\kappa T} \cosh \kappa T \frac{1 + \kappa T}{\Sigma_T}\right). \quad (332)$$

The second line transforms to Cartan phase-space coordinates. Unnormalizing changes the Gaussian's prefactor; transforming to Cartan coordinates changes the Gaussian. The final step is to undo the gauge transformation to get back to the Cartan reduced distribution,

$$C_T(Zx) = e^{f(Zx)} B_T(Zx) = \frac{1}{\sinh 2\kappa T} \delta(r - 2\kappa T) \frac{2}{\Sigma_T} \exp\left(-\frac{|\beta - \alpha|^2}{\Sigma_T}\right), \quad (333)$$

which matches the solution in Equation (227), which was obtained from the PDE for $C_t(Zx)$ displayed in Equation (221).

It is fair to ask whether the point of this section is just to provide a different, more complicated route to the solution for $C_T(Zx)$. We think it is more than that, and here is why. It all comes back to the elephant in the room, that is, how one handles the normalization or scaling of the Kraus operators that comes from the center Z . The Cartan reduced distribution, defined in Equation (216) and determined from the diffusion equation in Equation (221), succeeds in representing POVM completeness by marginalizing the ur -distribution $D_T(x)$ over the Cartan-center normalization, $e^{-2l} = e^{2f(Zx)} e^{-2s}$; this gives a distribution uniform in the Cartan sum variable $\beta + \alpha$ and thus spread over all of phase space in a way that gives POVM completeness. As a consequence, however, the moments of $C_T(Zx)$ are not those of the stochastic integrals for the phase-plane variables. The Harish-Chandra reduced distribution $B_T(Zx)$, defined in Equation (247), marginalizes $D_T(x)$ over the Harish-Chandra-center normalization, e^{-2s} . After normalization to unity by the factor \mathcal{N}_T , the normalized distribution $\tilde{B}_T(Zx)$ is determined from the diffusion equation shown in Equation (296) or by applying the path-integral expression of Equation (267), with its modified path measure, to the stochastic integrals for the Harish-Chandra phase-space variables. The route from $\tilde{B}_T(Zx)$ to POVM completeness runs backwards through the normalization factor \mathcal{N}_T and the anti-Gaussian gauge transformation $e^{2f(Zx)}$ and arrives at $C_T(Zx)$. The point of this section, one might say, is to find and reveal these connections among path integrals, diffusion equations, and SDEs; discovering these connections, driven in this paper by a combination of necessity and opportunity, allows us to re-unite the three faces of the stochastic trinity. Accessing the entire stochastic trinity by using a positive gauge transformation that is grounded in a problem's Lie group—this, we hope, might be generally applicable to Feynman-Kac formulas for non-normalization-preserving diffusion equations.

5. Concluding Remarks. The Stochastic Trinity

We set out on the project of analyzing simultaneous measurements of noncommuting observables [39,40] with the goal of showing that such measurements end up with a POVM in the overcomplete coherent-state basis. We now think that we have uncovered something more ambitious, a distinctive new window into the space of quantum dynamics. The formulation of the problem of continual, differential measurements invites one—compels one, really—to think in terms of the paths of Wiener outcome increments. These outcome paths, as sample paths drawn from the Wiener measure, know nothing about the space in which they are wandering. When they are instantiated in the exponents of Kraus operators, however, the time-ordered products of the differential Kraus operators generate a (complex) Lie-group manifold, the instrumental Lie group, in which the Kraus operators are, in the way of groups, both the transformations and the moving points. The motion in the instrumental Lie-group manifold is described by Kraus-operator SDEs or by the

diffusion of the KOD, as embodied in an FPK diffusion equation. The continuous, but not differentiable, paths are handled effortlessly by the Itô calculus of the outcome increments, with its terms of order \sqrt{dt} and dt . This requires getting just beyond the linear structure of vector fields (right-invariant derivatives) and one-forms (right-invariant one-forms) on the Lie-group manifold, as is evident from our discussion of Stratonovich vs. Itô SDEs and in the derivation of the FPK diffusion equation, where right-invariant derivatives end up as diffusive second derivatives. We end up in a very comfortable place, working in all three corners of the stochastic trinity: path integrals, FPK diffusion equations, and SDEs, all three describing motion, equivalently, on the instrumental Lie-group manifold.

Why do others not find the same comfort in all three faces of the trinity? Field theorists, interested in the propagator of closed-system dynamics, have a different way of handling the continuous, but not differentiable paths, coming from a Stratonovich calculus way of dealing with the temporal derivatives in the kinetic terms in a Lagrangian. While they usually have an equivalent Schrödinger-like equation for the propagator, they do not have the analog of SDEs, even though the problem often undergoes a Wick rotation into “Euclidean spacetime”.

Open-systems theorists, both in condensed matter physics and in quantum optics, generally work with master equations or stochastic master equations for the evolution of quantum states and sometimes with diffusion equations for a probability distribution associated with the states. Those who start with diffusion equations can avail themselves of a Feynman-Kac formula for a path integral, but the connection to SDEs has generally not been made. The reason for not using all three faces of the stochastic trinity is, we think, a failure to identify the appropriate Lie-group manifold; this failure is connected to the emphasis on the evolution of quantum states, which obscures nearly entirely the Lie-group manifold on which the open-system dynamics occur.

Suffice it to say that we think we have found something: the home of quantum dynamics, the Lie-group manifold that supports all three faces of the trinity. The exhausted reader who has survived to read this concluding sentence of a very long paper might be pleased—or so we hope—to learn that our ambition is larger than was evident at the beginning.

Author Contributions: Both authors participated in the research and analysis that underlies this article and in the writing of the article. All authors have read and agreed to the published version of the manuscript.

Funding: This work was supported in part by the Center for Quantum Information and Control at the University of New Mexico and in part by the U.S. Department of Energy, Office of Science, Office of Advanced Scientific Computing Research, under the Quantum Computing Application Teams (QCAT) program. Sandia National Laboratories is a multimission laboratory managed and operated by NTESS, LLC., a wholly owned subsidiary of Honeywell International, Inc., for the U.S. DOE’s NNSA under contract DE-NA-0003525. This paper describes objective technical results and analysis. Any subjective views or opinions that might be expressed in the paper do not necessarily represent the views of the U.S. Department of Energy or the United States Government.

Institutional Review Board Statement: Not applicable.

Data Availability Statement: This project generated analysis and theory, not data.

Acknowledgments: This article is part of C.S.J.’s vision for giving an autonomous, universal, transformation-group-theoretic description of measuring instruments. This vision grew out of an appreciation of the differential-geometric treatment of Lie groups along the lines considered by Cartan and others. C.S.J. thanks Mohan Sarovar for all of the helpful discussions and financial support through Sandia National Laboratories and C.M.C. for the incredibly fruitful collaboration. C.M.C. supported himself and is eternally grateful for the opportunity to work with C.S.J.

Conflicts of Interest: The authors declare no conflict of interest.

Appendix A. Exterior Derivatives of the Right-Invariant Basis

The Lie algebra of right-invariant derivatives is equivalent to the exterior algebra of right-invariant one-forms. This can be seen by considering the second exterior differential of a scalar function,

$$0 = dd f \quad (\text{A1})$$

$$= d\left(\underline{e}_\mu[f]\theta^\mu\right) \quad (\text{A2})$$

$$= \underline{e}_\mu\left[\underline{e}_\nu[f]\right]\theta^\mu \wedge \theta^\nu + \underline{e}_\lambda[f]d\theta^\lambda \quad (\text{A3})$$

$$= \frac{1}{2}\left(\underline{e}_\mu\left[\underline{e}_\nu[f]\right] - \underline{e}_\nu\left[\underline{e}_\mu[f]\right]\right)\theta^\mu \wedge \theta^\nu + \underline{e}_\lambda[f]d\theta^\lambda \quad (\text{A4})$$

$$= \frac{1}{2}\left[\underline{e}_\mu, \underline{e}_\nu\right][f]\theta^\mu \wedge \theta^\nu + \underline{e}_\lambda[f]d\theta^\lambda \quad (\text{A5})$$

$$= \underline{e}_\lambda[f]\left(-\frac{1}{2}c_{\mu\nu}{}^\lambda\theta^\mu \wedge \theta^\nu + d\theta^\lambda\right), \quad (\text{A6})$$

which, by linear independence, means that

$$\left[\underline{e}_\mu, \underline{e}_\nu\right] = -c_{\mu\nu}{}^\lambda \underline{e}_\lambda \quad \Longleftrightarrow \quad d\theta^\lambda = \frac{1}{2}c_{\mu\nu}{}^\lambda \theta^\mu \wedge \theta^\nu. \quad (\text{A7})$$

Specifically, if we write out all of the different Lie brackets,

$$H_o = \text{no such brackets}, \quad (\text{A8})$$

$$Q = [H_o, -iP], \quad (\text{A9})$$

$$P = [H_o, iQ], \quad (\text{A10})$$

$$-iP = [H_o, Q], \quad (\text{A11})$$

$$iQ = [H_o, P], \quad (\text{A12})$$

$$-\Omega = [Q, iP] = [iQ, P], \quad (\text{A13})$$

$$i\Omega = [Q, P] = [iQ, -iP], \quad (\text{A14})$$

we can read off all of the one-form differentials or “curls”,

$$d\theta^{i\Omega} = \theta^Q \wedge \theta^P + \theta^{iQ} \wedge \theta^{-iP}, \quad (\text{A15})$$

$$d\theta^{-\Omega} = \theta^{-iP} \wedge \theta^Q + \theta^{iQ} \wedge \theta^P, \quad (\text{A16})$$

$$d\theta^{-iP} = \theta^{H_o} \wedge \theta^Q, \quad (\text{A17})$$

$$d\theta^{iQ} = \theta^{H_o} \wedge \theta^P, \quad (\text{A18})$$

$$d\theta^Q = \theta^{H_o} \wedge \theta^{-iP}, \quad (\text{A19})$$

$$d\theta^P = \theta^{H_o} \wedge \theta^{iQ}, \quad (\text{A20})$$

$$d\theta^{-H_o} = 0, \quad (\text{A21})$$

Appendix B. Global Coordinate Transformation between Cartan and Harish-Chandra Coordinates

The global transformation from Cartan to Harish-Chandra coordinates is most easily accomplished by inserting the displacement-operator normal ordering of Equation (76) into the Cartan decomposition of Equation (156) and then pushing the factors around until they are in Harish-Chandra form:

$$(D_\beta e^{i\Omega\phi})e^{-H_0r-\Omega\ell}D_\alpha^{-1} = e^{-\Omega[\ell+\frac{1}{2}(|\beta|^2+|\alpha|^2)]}e^{i\Omega\phi}e^{a^\dagger\beta}e^{-a\beta^*}e^{-H_0r}e^{-a^\dagger\alpha}e^{a\alpha^*} \quad (\text{A22})$$

$$= e^{-\Omega[\ell+\frac{1}{2}(|\beta|^2+|\alpha|^2)]}e^{i\Omega\phi}e^{a^\dagger\beta}e^{-a\beta^*}e^{-a^\dagger e^{-r}\alpha}e^{-H_0r}e^{a\alpha^*} \quad (\text{A23})$$

$$= e^{-\Omega[\ell+\frac{1}{2}(|\beta|^2+|\alpha|^2)]}e^{i\Omega\phi}e^{\Omega\beta^*\alpha e^{-r}}e^{a^\dagger\beta}e^{-a^\dagger e^{-r}\alpha}e^{-a\beta^*}e^{-H_0r}e^{a\alpha^*} \quad (\text{A24})$$

$$= e^{-\Omega[\ell+\frac{1}{2}(|\beta|^2+|\alpha|^2)]}e^{i\Omega\phi}e^{\Omega\beta^*\alpha e^{-r}}e^{a^\dagger\beta}e^{-a^\dagger e^{-r}\alpha}e^{-H_0r}e^{-ae^{-r}\beta^*}e^{a\alpha^*} \quad (\text{A25})$$

$$= e^{a^\dagger(\beta-e^{-r}\alpha)}e^{-H_0r-\Omega[\ell+\frac{1}{2}(|\beta|^2+\frac{1}{2}|\alpha|^2)-\beta^*\alpha e^{-r}-i\phi]}e^{a(\alpha-e^{-r}\beta)^*}. \quad (\text{A26})$$

Identifying the parameters of the last line gives the coordinate transformation,

$$\nu = \beta - e^{-r}\alpha, \quad (\text{A27})$$

$$\mu = \alpha - e^{-r}\beta, \quad (\text{A28})$$

$$-s + i\psi = z = -\ell - f + i(\phi + \xi), \quad (\text{A29})$$

where

$$f \equiv \frac{1}{2}(|\beta|^2 + |\alpha|^2 - \beta^*\alpha e^{-r} - \beta\alpha^* e^{-r}) = \frac{1-e^{-r}}{4}|\beta + \alpha|^2 + \frac{1+e^{-r}}{4}|\beta - \alpha|^2, \quad (\text{A30})$$

$$\xi \equiv e^{-r}\frac{\beta^*\alpha - \beta\alpha^*}{2i} = e^{-r}\frac{(\beta - \alpha)^*(\beta + \alpha) - (\beta - \alpha)(\beta + \alpha)^*}{4i}, \quad (\text{A31})$$

are functions of the RWI G/Z .

The inverse coordinate transformation is

$$\beta = \frac{e^r\nu + \mu}{2\sinh r}, \quad (\text{A32})$$

$$\alpha = \frac{e^r\mu + \nu}{2\sinh r}, \quad (\text{A33})$$

$$\ell = s - f, \quad (\text{A34})$$

$$\phi = \psi - \xi. \quad (\text{A35})$$

It is quite useful to notice that the sum and difference phase-space coordinates simply rescale under this transformation,

$$\beta + \alpha = \frac{\nu + \mu}{1 - e^{-r}} = \frac{e^{r/2}(\nu + \mu)}{2\sinh(r/2)}, \quad (\text{A36})$$

$$\beta - \alpha = \frac{\nu - \mu}{1 + e^{-r}} = \frac{e^{r/2}(\nu - \mu)}{2\cosh(r/2)}. \quad (\text{A37})$$

We also have

$$f = \frac{|\nu + \mu|^2}{4(1 - e^{-r})} + \frac{|\nu - \mu|^2}{4(1 + e^{-r})} = \frac{(|\nu|^2 + |\mu|^2)e^r + \nu^*\mu + \nu\mu^*}{4\sinh r}, \quad (\text{A38})$$

$$\xi = \frac{(\nu - \mu)^*(\nu + \mu) - (\nu - \mu)(\nu + \mu)^*}{8i\sinh r} = \frac{\nu^*\mu - \nu\mu^*}{4i\sinh r}. \quad (\text{A39})$$

It is crucial to appreciate that the transformation from Harish-Chandra to Cartan coordinates is singular at $r = 0$, with the consequence that positive transformations of the form e^{Qq+Pp} are not represented in Cartan coordinates. At $r = 0$, the Harish-Chandra decomposition becomes

$$x_{r=0} = e^{\Omega z}e^{a^\dagger\nu}e^{a\mu^*} = e^{\Omega(z-\frac{1}{2}\nu\mu^*)}e^{a^\dagger\nu+a\mu^*} = e^{\Omega(z-\frac{1}{2}\nu\mu^*)}D_{(\nu-\mu)/2}e^{a^\dagger\frac{1}{2}(\nu+\mu)+a\frac{1}{2}(\nu+\mu)^*}. \quad (\text{A40})$$

Each displacement operator and each positive transformation of the form e^{Qq+Pp} is represented uniquely in the four-plane of ν and μ coordinates; in particular, the identity operator has coordinates $r = 0$ and $z = \nu = \mu = 0$.

At $r = 0$, the Cartan decomposition reduces to

$$x_{r=0} = e^{i\Omega\phi} e^{-\Omega\ell} D_{\beta} D_{-\alpha} = e^{i\Omega[\phi + (\beta^* \alpha - \beta \alpha^*)/2i]} e^{-\Omega\ell} D_{\beta-\alpha}. \quad (\text{A41})$$

The Cartan coordinate singularity at $r = 0$ is that $x_{r=0}$ does not depend on $\beta + \alpha$. A displacement operator $D_{\beta-\alpha} = D_{\tau}$ is represented by a plane of $\beta + \alpha$ values; specifically, by all the coordinates satisfying $\beta - \alpha = \tau$, $r = 0$, $\ell = 0$, and $\phi = (\beta + \alpha)^* \tau - (\beta + \alpha) \tau^* / 4i$. Most importantly, the identity 1 is represented by the plane of Cartan coordinate values satisfying $\beta = \alpha$ and $r = \phi = \ell = 0$. Positive transformations of the form e^{Qq+Pp} are not represented at all in Cartan coordinates; this is not a problem because these positive operators lie on a boundary that is not accessible to the Kraus operators of SPQM.

Appendix C. Solution for Harish-Chandra Center Coordinate

To solve for the Harish-Chandra center coordinate z from the SDE of Equation (146), it is best to work with the sums over Wiener increments that underlie the Itô stochastic integrals. Thus we begin by writing the solution for the post-measurement Harish-Chandra coordinate ν , given in Equation (148), as

$$\nu_N = \sum_{k=0}^{T/dt-1} \sqrt{\kappa} dw_{kdt} e^{-2\kappa(T-t_k)} = \sum_{k=0}^{N-1} \sqrt{\kappa} dw_k e^{-2\kappa(N-k)dt}, \quad (\text{A42})$$

where $dw_k = dw_{kdt}$ and $t_k = k dt$ ($t_N = T = N dt$). Notice that the initial condition $\nu_0 = 0$ is enforced by having no terms in the sum for $N = 0$. The solution for z is

$$z_N = \frac{1}{2} \sum_{k=0}^{N-1} \kappa |dw_k|^2 + \sum_{k=1}^{N-1} \nu_k \sqrt{\kappa} dw_k^* \quad (\text{A43})$$

$$= \frac{1}{2} \sum_{k=0}^{N-1} \kappa |dw_k|^2 + \sum_{k=1}^{N-1} \sum_{l=0}^{k-1} \kappa dw_k^* dw_l e^{-2\kappa|k-l|dt}, \quad (\text{A44})$$

where we omit the $k = 0$ term on the second line since $\nu_0 = 0$ and where we can insert the absolute value because $k > l$. Notice that, as for ν , the way the initial condition $z_0 = 0$ is enforced is through the fact that there are no terms in the sums when $N = 0$. Now, define

$$y = \sum_{k=1}^{N-1} \sum_{l=0}^{k-1} \kappa dw_k^* dw_l e^{-2\kappa|k-l|dt}, \quad (\text{A45})$$

and manipulate y^* by specifying the summing range in an equivalent way and then switching k and l ,

$$y^* = \sum_{l=0}^{N-2} \sum_{k=l+1}^{N-1} \kappa dw_k dw_l^* e^{-2\kappa|k-l|dt} = \sum_{k=0}^{N-2} \sum_{l=k+1}^{N-1} \kappa dw_k^* dw_l e^{-2\kappa|k-l|dt}. \quad (\text{A46})$$

With this result, we can write our sums as including the entire range of values, $k, l = 1 \dots, N-1$. Converting to the real and imaginary parts of $z = -s + i\psi$, we have

$$-s_N = \text{Re}(z_N) = \frac{1}{2} \left(\sum_{k=0}^{N-1} \kappa |dw_k|^2 + \sum_{k=1}^{N-1} \sum_{l=0}^{k-1} \kappa dw_k^* dw_l e^{-2\kappa|k-l|dt} + \sum_{k=0}^{N-2} \sum_{l=k+1}^{N-1} \kappa dw_k^* dw_l e^{-2\kappa|k-l|dt} \right) \quad (\text{A47})$$

$$= \frac{1}{2} \sum_{k=0}^{N-1} \sum_{l=0}^{N-1} \kappa dw_k^* dw_l e^{-2\kappa|k-l|dt}, \quad (\text{A48})$$

$$i\psi_N = i\text{Im}(z_N) = \frac{1}{2} \left(\sum_{k=1}^{N-1} \sum_{l=0}^{k-1} \kappa dw_k^* dw_l e^{-2\kappa|k-l|dt} - \sum_{k=0}^{N-2} \sum_{l=k+1}^{N-1} \kappa dw_k^* dw_l e^{-2\kappa|k-l|dt} \right) \quad (\text{A49})$$

$$= \frac{1}{2} \sum_{k=0}^{N-1} \sum_{l=0}^{N-1} \kappa dw_k^* dw_l e^{-2\kappa|k-l|dt} \text{sgn}(k-l), \quad (\text{A50})$$

where the sign function is defined as

$$\text{sgn}(u) = \begin{cases} 1, & u > 0, \\ 0, & u = 0, \\ -1, & u < 0. \end{cases} \quad (\text{A51})$$

The center variable of Equation (A43) is therefore

$$z_N = -s_N + i\psi_N = \sum_{k=0}^{N-1} \sum_{l=0}^{N-1} \kappa dw_k^* dw_l e^{-2\kappa|k-l|dt} \frac{1}{2} (1 + \text{sgn}(k-l)), \quad (\text{A52})$$

where one recognizes

$$\frac{1}{2} (1 + \text{sgn}(u)) = H(u) = \begin{cases} 1, & u > 0, \\ \frac{1}{2}, & u = 0, \\ 0 & u < 0. \end{cases} \quad (\text{A53})$$

as the Heaviside step function with the choice $H(0) = \frac{1}{2}$. This appendix is really an exercise in relating the real and imaginary parts of z to its symmetric and antisymmetric parts and, in the process, getting the weighting of the diagonal ($k = l$) term right—equivalently, making the right choice for the $u = 0$ value of the sign and Heaviside functions. The incremental Itô calculus, with its quite explicit diagonal terms $\frac{1}{2}\kappa|dw_k|^2$, leaves no doubt about the right choice.

Converting back to integrals gives

$$-s_T = \frac{1}{2} \int_0^{T-} \int_0^{T-} \kappa dw_t^* dw_s e^{-2\kappa|t-s|}, \quad (\text{A54})$$

$$i\psi_T = \frac{1}{2} \int_0^{T-} \int_0^{T-} \kappa dw_t^* dw_s e^{-2\kappa|t-s|} \text{sgn}(t-s), \quad (\text{A55})$$

which lead to the formula for $z_T = -s_T + i\psi_T$ in Equation (150). To use these integrals, one should return to the incremental sums. If one started with the Stratonovich-form SDE shown in Equation (190) for z , one could convert it to an ordinary temporal differential equation with the Wiener increments replaced by δ -functions. Integrating that equation would give the same results as here, but only after careful attention to how to weight the diagonal ($t = s$) δ -functions in the integral. Thinking in terms of the Stratonovich-form SDE for z , it is clear that solving the equation for z —what we have done in this appendix—is identical to deriving a fluctuation-dissipation theorem for the correlation of v and μ .

Appendix D. Local Transformations between Cartan-Coordinate and Right-Invariant Frames

The Cartan decomposition of a position $x \in \text{IWH}$ is

$$x = \left(D_\beta e^{i\Omega\phi} \right) e^{-H_0 r - \Omega \ell} D_\alpha^{-1}. \quad (\text{A56})$$

Before differentiating x with respect to the Cartan coordinates, it is useful to differentiate the displacement operator with respect to its arguments using the ordered forms in Equation (75):

$$\partial_{\alpha_1} D_\alpha = (-iP + \frac{1}{2}\alpha_2 i\Omega) D_\alpha, \quad (\text{A57})$$

$$\partial_{\alpha_2} D_\alpha = (iQ - \frac{1}{2}\alpha_1 i\Omega) D_\alpha. \quad (\text{A58})$$

We will also need the conjugations

$$\frac{1}{\sqrt{2}} D_\alpha (Q + iP) D_\alpha^\dagger = D_\alpha a D_\alpha^\dagger = a - \alpha = \frac{1}{\sqrt{2}} (Q - \alpha_1 + i(P - \alpha_2)), \quad (\text{A59})$$

and

$$e^{-rH_0} a e^{rH_0} = a e^r, \quad (\text{A60})$$

$$e^{-rH_0} a^\dagger e^{rH_0} = a^\dagger e^{-r}, \quad (\text{A61})$$

$$e^{-rH_0} Q e^{rH_0} = Q \cosh r + iP \sinh r, \quad (\text{A62})$$

$$e^{-rH_0} P e^{rH_0} = P \cosh r - iQ \sinh r. \quad (\text{A63})$$

Differentiating the position x gives

$$\partial_\phi x = i\Omega x, \quad (\text{A64})$$

$$-\partial_\ell x = \Omega x, \quad (\text{A65})$$

$$\partial_{\beta_1} x = -iP x + \frac{1}{2}\beta_2 i\Omega x, \quad (\text{A66})$$

$$\partial_{\beta_2} x = iQ x - \frac{1}{2}\beta_1 i\Omega x, \quad (\text{A67})$$

$$\begin{aligned} -\partial_r x &= D_\beta H_0 D_\beta^\dagger x \\ &= \frac{1}{2}((Q - \beta_1)^2 + (P - \beta_2)^2) x \\ &= H_0 x - \beta_1 Q x - \beta_2 P x + \frac{1}{2}(\beta_1^2 + \beta_2^2) \Omega x, \end{aligned} \quad (\text{A68})$$

$$\begin{aligned} \partial_{\alpha_1} x &= D_\beta e^{-rH_0} (iP + \frac{1}{2}\alpha_2 i\Omega) e^{rH_0} D_\beta^\dagger x \\ &= D_\beta (iP \cosh r + Q \sinh r + \frac{1}{2}\alpha_2 i\Omega) D_\beta^\dagger x \\ &= \cosh r iP x + \sinh r Q x - \beta_1 \sinh r \Omega x - (\beta_2 \cosh r - \frac{1}{2}\alpha_2) i\Omega x, \end{aligned} \quad (\text{A69})$$

$$\begin{aligned} \partial_{\alpha_2} x &= D_\beta e^{-rH_0} (-iQ - \frac{1}{2}\alpha_1 i\Omega) e^{rH_0} D_\beta^\dagger x \\ &= D_\beta (-iQ \cosh r + P \sinh r - \frac{1}{2}\alpha_1 i\Omega) D_\beta^\dagger x \\ &= -\cosh r iQ x + \sinh r P x - \beta_2 \sinh r \Omega x + (\beta_1 \cosh r - \frac{1}{2}\alpha_1) i\Omega x. \end{aligned} \quad (\text{A70})$$

By the chain rule, we have the frame transformation,

$$\begin{aligned}
\partial_\phi &= +i\Omega \\
\partial_\ell &= -\Omega \\
\partial_{\beta_1} &= +\frac{1}{2}\beta_2 i\Omega - iP \\
\partial_{\beta_2} &= -\frac{1}{2}\beta_1 i\Omega + iQ \\
\partial_{\alpha_1} &= -(\beta_2 \cosh r - \frac{1}{2}\alpha_2)i\Omega - \beta_1 \sinh r \Omega + \cosh r iP + \sinh r Q \\
\partial_{\alpha_2} &= +(\beta_1 \cosh r - \frac{1}{2}\alpha_1)i\Omega - \beta_2 \sinh r \Omega - \cosh r iQ + \sinh r P \\
\partial_r &= -\frac{1}{2}(\beta_1^2 + \beta_2^2)\Omega + \beta_1 Q + \beta_2 P - H_o
\end{aligned} \quad . \quad (A71)$$

Inverting the transformation gives

$$i\Omega = \partial_\phi, \quad (A72)$$

$$-\Omega = \partial_\ell, \quad (A73)$$

$$-iP = \partial_{\beta_1} - \frac{1}{2}\beta_2 \partial_\phi, \quad (A74)$$

$$iQ = \partial_{\beta_2} + \frac{1}{2}\beta_1 \partial_\phi, \quad (A75)$$

$$Q = \nabla_1 - \beta_1 \partial_\ell + \frac{\beta_2 \cosh r - \alpha_2}{2 \sinh r} \partial_\phi, \quad (A76)$$

$$P = \nabla_2 - \beta_2 \partial_\ell - \frac{\beta_1 \cosh r - \alpha_1}{2 \sinh r} \partial_\phi, \quad (A77)$$

$$\begin{aligned}
-H_o &= \partial_r - \beta_1 Q - \beta_2 P + \frac{\beta_1^2 + \beta_2^2}{2} \Omega \\
&= \partial_r - \beta_1 \nabla_1 - \beta_2 \nabla_2 + \frac{\beta_1^2 + \beta_2^2}{2} \partial_\ell + \frac{\beta_1 \alpha_2 - \beta_2 \alpha_1}{2 \sinh r} \partial_\phi,
\end{aligned} \quad (A78)$$

where

$$\nabla_j \equiv \frac{1}{\sinh r} (\partial_{\alpha_j} + \cosh r \partial_{\beta_j}). \quad (A79)$$

Transposing the transformation gives

$$\begin{aligned}
\theta^{i\Omega} &= d\phi + \frac{1}{2}\beta_2 d\beta_1 - \frac{1}{2}\beta_1 d\beta_2 - (\beta_2 \cosh r - \frac{1}{2}\alpha_2) d\alpha_1 + (\beta_1 \cosh r - \frac{1}{2}\alpha_1) d\alpha_2 \\
&= d\phi + \frac{1}{2}(\beta_2 d\beta_1 - \beta_1 d\beta_2) + \frac{1}{2}(\alpha_2 d\alpha_1 - \alpha_1 d\alpha_2) + \cosh r (\beta_1 d\alpha_2 - \beta_2 d\alpha_1),
\end{aligned} \quad (A80)$$

$$\theta^{-\Omega} = d\ell + \frac{1}{2}(\beta_1^2 + \beta_2^2) dr + \sinh r (\beta_1 d\alpha_1 + \beta_2 d\alpha_2), \quad (A81)$$

$$\theta^{-iP} = d\beta_1 - \cosh r d\alpha_1, \quad (A82)$$

$$\theta^{iQ} = d\beta_2 - \cosh r d\alpha_2, \quad (A83)$$

$$\theta^Q = \beta_1 dr + \sinh r d\alpha_1, \quad (A84)$$

$$\theta^P = \beta_2 dr + \sinh r d\alpha_2, \quad (A85)$$

$$\theta^{-H_o} = dr. \quad (A86)$$

With the one-form transformations in hand, the Haar measure of Equation (110) becomes in Cartan coordinates,

$$d^7\mu(x) = d\phi d\ell \frac{d^2\beta}{\pi} dr \sinh^2 r \frac{d^2\alpha}{\pi}. \quad (A87)$$

The phase-plane measures are $d^2\beta = \frac{1}{2}d\beta_1 d\beta_2$ and $d^2\alpha = \frac{1}{2}d\alpha_1 d\alpha_2$. The factors of $1/\pi$ are conventional in quantum optics and ultimately come from the coherent-state completeness relation that is given in Equation (155).

It is easy to show that the left-invariant derivatives and one-forms can be obtained from the right-invariant quantities by changing the sign of the quantities associated with

anti-Hermitian operators and transforming the coordinates according to $\phi \leftrightarrow -\phi$ and $\beta \leftrightarrow \alpha$. The sign changes do not change the Haar measure, and the measure is unchanged by the coordinate transformation, which shows that the Haar measure is both right- and left-invariant.

Appendix E. Local Transformations between Harish-Chandra-Coordinate and Right-Invariant Frames

The Harish-Chandra decomposition of a position $x \in \text{IWH}$ is

$$x = e^{a^\dagger \nu} e^{-H_0 r + \Omega z} e^{a \mu^*}. \quad (\text{A88})$$

Break the complex coordinates into real and imaginary parts, as in Equation (134). Differentiating the position gives

$$\partial_{\nu_1} x = \frac{1}{\sqrt{2}} a^\dagger x, \quad (\text{A89})$$

$$\partial_{\nu_2} x = \frac{1}{\sqrt{2}} i a^\dagger x, \quad (\text{A90})$$

$$-\partial_s x = \Omega x, \quad (\text{A91})$$

$$\partial_\psi x = i \Omega x, \quad (\text{A92})$$

$$\begin{aligned} -\partial_r x &= e^{a^\dagger \nu} H_0 e^{-a^\dagger \nu} x \\ &= H_0 x - \nu a^\dagger x \end{aligned} \quad (\text{A93})$$

$$= H_0 x - \nu_1 \frac{1}{\sqrt{2}} a^\dagger x - \nu_2 \frac{1}{\sqrt{2}} i a^\dagger x,$$

$$\begin{aligned} \partial_{\mu_1} x &= \frac{1}{\sqrt{2}} e^{a^\dagger \nu} e^{-H_0 r} a e^{H_0 r} e^{-a^\dagger \nu} x \\ &= \frac{1}{\sqrt{2}} e^{a^\dagger \nu} (e^r a) e^{-a^\dagger \nu} x \\ &= \frac{1}{\sqrt{2}} e^r (a - \nu \Omega) x \end{aligned} \quad (\text{A94})$$

$$= e^r \left(\frac{1}{\sqrt{2}} a x - \frac{1}{2} \nu_1 \Omega x - \frac{1}{2} \nu_2 i \Omega x \right),$$

$$\begin{aligned} \partial_{\mu_2} x &= -i \frac{1}{\sqrt{2}} e^{a^\dagger \nu} e^{-H_0 r} a e^{H_0 r} e^{-a^\dagger \nu} x \\ &= e^r \left(-\frac{1}{\sqrt{2}} i a x - \frac{1}{2} \nu_2 \Omega x + \frac{1}{2} \nu_1 i \Omega x \right). \end{aligned} \quad (\text{A95})$$

By the chain rule, we have the frame transformation,

$$\begin{aligned} \partial_\psi &= i \Omega \\ -\partial_s &= \Omega \\ \partial_{\nu_1} &= \frac{1}{\sqrt{2}} a^\dagger \\ \partial_{\nu_2} &= \frac{1}{\sqrt{2}} i a^\dagger \\ e^{-r} \partial_{\mu_1} &= -\frac{1}{2} \nu_2 i \Omega - \frac{1}{2} \nu_1 \Omega \frac{1}{\sqrt{2}} a \\ e^{-r} \partial_{\mu_2} &= \frac{1}{2} \nu_1 i \Omega - \frac{1}{2} \nu_2 \Omega \frac{1}{\sqrt{2}} i a - \frac{1}{\sqrt{2}} i a \\ -\partial_r &= -\nu_1 \frac{1}{\sqrt{2}} a^\dagger - \nu_2 \frac{1}{\sqrt{2}} i a^\dagger H_0 \end{aligned} \quad (\text{A96})$$

Inverting the transformation gives

$$\overleftarrow{\Omega} = \partial_\psi, \quad (\text{A97})$$

$$-\overleftarrow{\Omega} = \partial_s, \quad (\text{A98})$$

$$\frac{1}{\sqrt{2}}\overleftarrow{a}^\dagger = \partial_{v_1}, \quad (\text{A99})$$

$$\frac{1}{\sqrt{2}}\overleftarrow{ia}^\dagger = \partial_{v_2}, \quad (\text{A100})$$

$$\frac{1}{\sqrt{2}}\overleftarrow{a} = e^{-r}\partial_{\mu_1} - \frac{1}{2}v_1\partial_s + \frac{1}{2}v_2\partial_\psi, \quad (\text{A101})$$

$$-\frac{1}{\sqrt{2}}\overleftarrow{ia} = e^{-r}\partial_{\mu_2} - \frac{1}{2}v_2\partial_s - \frac{1}{2}v_1\partial_\psi, \quad (\text{A102})$$

$$-\overleftarrow{H_o} = \partial_r - v_1\partial_{v_1} - v_2\partial_{v_2}. \quad (\text{A103})$$

Transposing the transformation gives

$$\theta^{i\Omega} = d\psi + \frac{1}{2}e^r(v_1d\mu_2 - v_2d\mu_1), \quad (\text{A104})$$

$$\theta^{-\Omega} = ds + \frac{1}{2}e^r(v_1d\mu_1 + v_2d\mu_2), \quad (\text{A105})$$

$$\sqrt{2}\theta^{a^\dagger} = dv_1 + v_1dr, \quad (\text{A106})$$

$$\sqrt{2}\theta^{ia^\dagger} = dv_2 + v_2dr, \quad (\text{A107})$$

$$\sqrt{2}\theta^a = e^rd\mu_1, \quad (\text{A108})$$

$$\sqrt{2}\theta^{-ia} = e^rd\mu_2, \quad (\text{A109})$$

$$\theta^{-H_o} = dr. \quad (\text{A110})$$

It is instructive to convert the right-invariant derivatives and forms to position and momentum,

$$-i\overleftarrow{P} = \frac{1}{\sqrt{2}}(\overleftarrow{a}^\dagger - \overleftarrow{a}) = \partial_{v_1} - e^{-r}\partial_{\mu_1} + \frac{1}{2}v_1\partial_s - \frac{1}{2}v_2\partial_\psi, \quad (\text{A111})$$

$$i\overleftarrow{Q} = \frac{1}{\sqrt{2}}(\overleftarrow{ia}^\dagger + \overleftarrow{ia}) = \partial_{v_2} - e^{-r}\partial_{\mu_2} + \frac{1}{2}v_2\partial_s + \frac{1}{2}v_1\partial_\psi, \quad (\text{A112})$$

$$\overleftarrow{Q} = \frac{1}{\sqrt{2}}(\overleftarrow{a}^\dagger + \overleftarrow{a}) = \nabla_1 - \frac{1}{2}v_1\partial_s + \frac{1}{2}v_2\partial_\psi, \quad (\text{A113})$$

$$\overleftarrow{P} = \frac{1}{\sqrt{2}}(\overleftarrow{ia}^\dagger - \overleftarrow{ia}) = \nabla_2 - \frac{1}{2}v_2\partial_s - \frac{1}{2}v_1\partial_\psi, \quad (\text{A114})$$

where the derivatives

$$\nabla_j = \partial_{v_j} + e^{-r}\partial_{\mu_j} \quad (\text{A115})$$

are equal to the derivatives of Equation (A79) and

$$\theta^{-iP} = \frac{1}{\sqrt{2}}(\theta^{a^\dagger} - \theta^a) = \frac{1}{2}(dv_1 - e^rd\mu_1 + v_1dr), \quad (\text{A116})$$

$$\theta^{iQ} = \frac{1}{\sqrt{2}}(\theta^{ia^\dagger} + \theta^{ia}) = \frac{1}{2}(dv_2 - e^rd\mu_2 + v_2dr), \quad (\text{A117})$$

$$\theta^Q = \frac{1}{\sqrt{2}}(\theta^{a^\dagger} + \theta^a) = \frac{1}{2}(dv_1 + e^rd\mu_1 + v_1dr), \quad (\text{A118})$$

$$\theta^P = \frac{1}{\sqrt{2}}(\theta^{ia^\dagger} - \theta^{ia}) = \frac{1}{2}(dv_2 + e^rd\mu_2 + v_2dr). \quad (\text{A119})$$

Comparing the expressions for right-invariant derivatives and one-forms in the Harish-Chandra coordinates of this appendix with the expressions in Cartan coordinates from Appendix D, one sees that these expressions relate the coordinate partial derivatives and one-forms in the two coordinate systems. The relations between coordinate partial derivatives and one-forms can, of course, be derived straightforwardly from the global coordinate transformations in Appendix B. We do not record these relations because the expressions in terms of the right-invariant quantities are more useful for our purposes. We remind the reader that, even though the ruler coordinate r is shared between Harish-Chandra and Cartan coordinates, the partial derivative ∂_r is different in the two systems because partial derivatives are defined by holding the other coordinates constant.

The Haar measure, written in terms of right-invariant one-forms in Equation (110), becomes in Harish-Chandra coordinates,

$$d^7\mu(x) = d\psi ds \frac{d^2v}{2\pi} dr e^{2r} \frac{d^2\mu}{2\pi}. \quad (\text{A120})$$

The factors of $1/2\pi$ in the Harish-Chandra phase-plane measures follow from transforming the Cartan Haar measure to Harish-Chandra coordinates. Just as for Cartan coordinates, it is easy to show that the left-invariant one-forms yield the same measure.

Appendix F. Delta Functions and the Singularity in Cartan coordinates

In this appendix, we find the Cartan form of the δ -function $\delta(x, 1)$ in two ways. The first uses the fact that $D_0(x) = \delta(x, 1)$, so we can step slightly away from $T = 0$ to $T = dt$ and find $D_{dt}(x)$ as an expression that limits to $\delta(x, 1)$ as $dt \rightarrow 0$. The second way is perhaps more straightforward: show that the expression for $\delta(x, 1)$ in Cartan coordinates transforms to the known Harish-Chandra-coordinate form of Equation (175). That said, since we needed to know $D_{dt}(x)$ in Section 3.5, we start here with the first way.

Starting from the general path-integral expression for the KOD, given in Equation (88), and setting $T = dt$,

$$D_{dt}(x) = \int d\mu(dw_0) \delta(x, \gamma(dw_0)), \quad (\text{A121})$$

and plugging in the Cartan-coordinate expression for the δ -function, shown in Equation (171), we have

$$D_{dt}(x) = \int d\mu(dw_0) \delta(\phi - \phi_{dt}) \delta(\ell - \ell_{dt}) \frac{1}{\sinh^2 r} \delta(r - r_{dt}) \pi \delta^2(\beta - \beta_{dt}) \pi \delta^2(\alpha - \alpha_{dt}). \quad (\text{A122})$$

It is useful to work in terms of sum and difference coordinates for the phase-space variables, so we note that

$$\pi \delta^2(\beta - \beta_{dt}) \pi \delta^2(\alpha - \alpha_{dt}) = 2\pi \delta^2(\beta + \alpha - (\beta_{dt} + \alpha_{dt})) 2\pi \delta^2(\beta - \alpha - (\beta_{dt} - \alpha_{dt})). \quad (\text{A123})$$

We now specialize the stochastic-integral solutions for the coordinates to the first increment. At $T = dt$, $r_{dt} = 2\kappa dt$ and the phase-space coordinates are

$$v_{dt} = \sqrt{\kappa} dw_0 e^{-2\kappa dt} = \sqrt{\kappa} dw_0, \quad \mu_{dt} = d\mu_0 = \sqrt{\kappa} dw_0, \quad (\text{A124})$$

$$v_{dt} + \mu_{dt} = \sqrt{\kappa} dw_0 (e^{-2\kappa dt} + 1) = 2\sqrt{\kappa} dw_0, \quad v_{dt} - \mu_{dt} = \sqrt{\kappa} dw_0 (e^{-2\kappa dt} - 1) = -2\kappa dt \sqrt{\kappa} dw_0, \quad (\text{A125})$$

$$\beta_{dt} = \sqrt{\kappa} dw_0 \operatorname{csch} 2\kappa dt = \frac{\sqrt{\kappa} dw_0}{2\kappa dt}, \quad \alpha_{dt} = \sqrt{\kappa} dw_0 \coth 2\kappa dt = \frac{\sqrt{\kappa} dw_0}{2\kappa dt}, \quad (\text{A126})$$

$$\beta_{dt} + \alpha_{dt} = \sqrt{\kappa} dw_0 \coth \kappa dt = \frac{\sqrt{\kappa} dw_0}{\kappa dt}, \quad \beta_{dt} - \alpha_{dt} = -\sqrt{\kappa} dw_0 \tanh \kappa dt = -\kappa dt \sqrt{\kappa} dw_0. \quad (\text{A127})$$

In each case, the first expression is exact, and the second is the leading-order contribution. What should give any reader pause are the divergences in the single-increment Cartan phase

variables. These divergences are an expression of the $r = 0$ singularity in Cartan coordinates, and they make the decisive contribution to our determination of $D_{dt}(x)$. It is important to understand that there are subtleties in trying to extract the single-increment behavior directly from the SDEs instead of from the stochastic integrals. For example, the SDEs in Equations (144) and (145) state that $v_{dt} = dv_0 = \sqrt{\kappa} dw_0$ and $\mu_{dt} = d\mu_0 = \sqrt{\kappa} dw_0$, so one misses the leading-order contribution to $v_{dt} - \mu_{dt}$ because that leading-order contribution is zero from the perspective of the SDEs. More tellingly, the $r = 0$ singularity appears directly in Equations (183) and (184) as singularities in $d\beta_0$ and $d\alpha_0$, preventing one from determining the single-increment values β_{dt} and α_{dt} from these equations. The single-increment values can be read off the SDEs in Equations (185) and (186), however, and this is because these are effectively equations for the Harish-Chandra coordinates.

We also need the leading-order single-increment behavior of the several center coordinates:

$$z_{dt} = -s_{dt} + i\psi_{dt} = \frac{1}{2}\kappa|dw_0|^2, \quad (\text{A128})$$

$$-\ell_{dt} = f_{dt} - s_{dt} = f_{dt} = 2\kappa dt \frac{|\beta_{dt} + \alpha_{dt}|^2}{4} = \frac{1}{2\kappa dt} \frac{|\nu_{dt} + \mu_{dt}|^2}{4} = \frac{\kappa|dw_0|^2}{2\kappa dt}, \quad (\text{A129})$$

$$\phi_{dt} = \psi_{dt} - \zeta_{dt} = 0. \quad (\text{A130})$$

Please appreciate that, as indicated, the leading-order contribution to ℓ_{dt} comes entirely from f_{dt} .

Putting all this into Equation (A121), one finds that

$$D_{dt}(x) = \delta(\phi) \int \frac{d^2(\sqrt{\kappa} dw_0)}{\pi \kappa dt} e^{-\kappa|dw_0|^2/\kappa dt} \delta\left(\ell + \frac{\kappa|dw_0|^2}{2\kappa dt}\right) \frac{1}{\sinh^2 r} \delta(r - 2\kappa dt) \times 2\pi\delta^2\left(\beta + \alpha - \frac{\sqrt{\kappa} dw_0}{\kappa dt}\right) 2\pi\delta^2(\beta - \alpha). \quad (\text{A131})$$

Now, writing

$$\delta^2\left(\beta + \alpha - \frac{\sqrt{\kappa} dw_0}{\kappa dt}\right) = (\kappa dt)^2 \delta^2(\sqrt{\kappa} dw_0 - \kappa dt(\beta + \alpha)), \quad (\text{A132})$$

we integrate over the initial Wiener increment dw_0 and find

$$D_{dt}(x) = \delta(\phi) \delta\left(\ell + \frac{1}{2}\kappa dt|\beta + \alpha|^2\right) \frac{1}{\sinh^2 r} \delta(r - 2\kappa dt) 2\pi \frac{\kappa dt}{\pi} e^{-\kappa dt|\beta + \alpha|^2} 2\pi\delta^2(\beta - \alpha). \quad (\text{A133})$$

It is easy to see that this integrates to 1 over the Cartan-coordinate Haar measure in Equation (169).

It is understood that, in defining $\delta(x, 1)$, we take the limit $dt \rightarrow 0$. We might as well take these limits, thus writing

$$\delta(x, 1) = \delta(\phi) \delta(\ell) \frac{1}{\sinh^2 r} \delta(r) \left(2\pi \lim_{dt \rightarrow 0} \frac{\kappa dt}{\pi} e^{-\kappa dt|\beta + \alpha|^2}\right) 2\pi\delta^2(\beta - \alpha) \quad (\text{A134})$$

$$= \delta(\phi) \delta(\ell) \frac{1}{\sinh^2 r} \delta(r) \left(\pi \lim_{dt \rightarrow 0} \frac{4\kappa dt}{\pi} e^{-4\kappa dt|\alpha|^2}\right) \pi\delta^2(\beta - \alpha). \quad (\text{A135})$$

This has singular behavior in the ruler r , but this singularity goes away when integrating against the Haar measure. More important is the Gaussian in $|\beta + \alpha|^2$, which limits to an infinitely wide, normalized Gaussian as $dt \rightarrow 0$. As discussed in Appendix B, the identity 1 is represented by $\phi = \ell = r = 0$ and $\beta = \alpha$, with $\beta + \alpha$ free to take on any complex value. The infinitely wide Gaussian in $\beta + \alpha$ and the δ -function $\delta(\beta - \alpha)$ express this in $\delta(x, 1)$.

Now let us transform Equation (A133) to Harish-Chandra coordinates. With $r = 2\kappa dt$, the sum and difference phase-space variables are related by

$$\beta + \alpha = \frac{\nu + \mu}{2\kappa dt}, \quad (\text{A136})$$

$$\beta - \alpha = \frac{\nu - \mu}{2}, \quad (\text{A137})$$

and the Harish-Chandra center variables are given by, when $\alpha = \beta$,

$$s = \ell + f = \ell + \frac{1}{2}\kappa dt |\beta + \alpha|^2, \quad (\text{A138})$$

$$\psi = \phi + \xi = \phi. \quad (\text{A139})$$

Substituting into Equation (A133) gives

$$D_{dt}(x) = \delta(\phi) \delta(s) \delta(r - 2\kappa dt) 4\pi \frac{1}{\pi 4\kappa dt} e^{-|\nu + \mu|^2 / 4\kappa dt} 4\pi \delta^2(\nu - \mu). \quad (\text{A140})$$

The wide Gaussian in $\beta + \alpha$ becomes a narrow Gaussian in $\nu + \mu$, so when we take the limit $dt \rightarrow 0$, we get

$$\delta(x, 1) = \delta(\phi) \delta(s) \delta(r) 4\pi \delta^2(\nu + \mu) 4\pi \delta^2(\nu - \mu) = \delta(\phi) \delta(s) \delta(r) 2\pi \delta^2(\nu) 2\pi \delta^2(\mu), \quad (\text{A141})$$

in agreement with Equation (175) and consistent with the fact that, in Harish-Chandra coordinates, the identity is represented uniquely by $\phi = s = r = 0$ and $\nu = \mu = 0$.

It is now straightforward to integrate over the center variables in either coordinate system to find the δ -function on G/Z ,

$$\delta(Zx, Z1) = \frac{1}{\sinh^2 r} \delta(r) \left(\pi \lim_{dt \rightarrow 0} \frac{4\kappa dt}{\pi} e^{-4\kappa dt |\alpha|^2} \right) \pi \delta^2(\beta - \alpha) \quad (\text{A142})$$

$$= \delta(r) 2\pi \delta^2(\nu) 2\pi \delta^2(\mu). \quad (\text{A143})$$

Appendix G. Riccati Equations for the Three Moments

In this appendix we derive the Riccati equations for the three second moments, n_T , m_T , and q_T ; the Riccati equations, given in Equations (313)–(315), are derived from the expressions for the second moments in Equations (305), (307) and (309). When we work in terms of the modified path measure $\mathcal{D}\mu_M[dw_{[0,T)}]$ of Equation (260), the correlation of two outcome increments at times t and s within the interval $[0, T)$ depends on the overall time T , as noted in Equation (261). Hence, it seems natural here to reserve t to denote times within the interval $[0, T)$ and to increment the overall time from T to $T + dT$. As a consequence, unlike anywhere else in the paper, we denote the duration of the increments by dT instead of dt .

It is quite convenient to introduce a bra-ket notation and to define vectors

$$|+_T\rangle \equiv \sqrt{dT} \sum_{k=0}^{N-1} e^{-2\kappa dT(N-k)} |k\rangle, \quad (\text{A144})$$

$$|-_T\rangle \equiv \sqrt{dT} \sum_{k=0}^{N-1} e^{-2\kappa dT k} |k\rangle, \quad (\text{A145})$$

where the kets $|k\rangle$ are orthonormal. The matrix elements of M_T are

$$(M_T)_{kl} = \langle k | M_T | l \rangle = \delta_{kl} - \kappa dT e^{-2\kappa dT |k-l|}, \quad k, l = 0, \dots, N-1. \quad (\text{A146})$$

The utility of the bra-ket notation becomes apparent when we write the three moments as matrix elements of M_T^{-1} ,

$$n_T = \kappa \langle +_T | M_T^{-1} | +_T \rangle, \quad (\text{A147})$$

$$m_T = \kappa \langle -_T | M_T^{-1} | -_T \rangle, \quad (\text{A148})$$

$$q_T = \kappa \langle +_T | M_T^{-1} | -_T \rangle = \kappa \langle -_T | M_T^{-1} | +_T \rangle. \quad (\text{A149})$$

Replacing M_T and M_t^{-1} with the unit matrix gives the moments with respect to the original Wiener measure; these moments are the inner products of the kets $|\pm_T\rangle$.

Now increment M_T forward one step in time from T to $T + dT$. What happens is really nothing,

$$\langle k | M_{T+dT} | l \rangle = \delta_{kl} - \kappa dT e^{-2\kappa dT |k-l|}, \quad k, l = 0, \dots, N, \quad (\text{A150})$$

except that an outside row and column, labeled by the index N , are added,

$$k, l = 0, \dots, N-1: \quad \langle k | M_{T+dT} | l \rangle = \langle k | M_T | l \rangle = \delta_{kl} - \kappa dT e^{-2\kappa dT |k-l|}, \quad (\text{A151})$$

$$k = 0, \dots, N-1: \quad \langle k | M_{T+dT} | N \rangle = -\kappa dT e^{-2\kappa dT (N-k)} = -\kappa \sqrt{dT} \langle k | +_T \rangle, \quad (\text{A152})$$

$$l = 0, \dots, N-1: \quad \langle N | M_{T+dT} | l \rangle = -\kappa dT e^{-2\kappa dT (N-k)} = -\kappa \sqrt{dT} \langle +_T | l \rangle, \quad (\text{A153})$$

$$\langle N | M_{T+dT} | N \rangle = 1 - \kappa dT. \quad (\text{A154})$$

The incremented matrix can then be written in block form relative to the outer row and column,

$$M_{T+dT} = M_T - \kappa \sqrt{dT} (|+_T\rangle \langle N| + |N\rangle \langle +_T|) + (1 - \kappa dT) |N\rangle \langle N|. \quad (\text{A155})$$

We also need to increment the kets $|\pm_T\rangle$:

$$|+_T+dT\rangle = \sqrt{dT} \sum_{k=0}^N e^{-2\kappa dT (N+1-k)} |k\rangle = e^{-2\kappa dT} (|+_T\rangle + \sqrt{dT} |N\rangle), \quad (\text{A156})$$

$$|-_T+dT\rangle = \sqrt{dT} \sum_{k=0}^N e^{-2\kappa dT k} |k\rangle = |-_T\rangle + \sqrt{dT} e^{-2\kappa T} |N\rangle. \quad (\text{A157})$$

The harder task, incrementing M_T^{-1} , is done using the Schur complement. Generally, for any matrix written in block form,

$$M = \begin{pmatrix} A & B \\ C & D \end{pmatrix}, \quad (\text{A158})$$

deriving a formal expression for the inverse begins with a block Gaussian elimination whereby the matrix is written as

$$M = \begin{pmatrix} I & 0 \\ CA^{-1} & I \end{pmatrix} \begin{pmatrix} A & 0 \\ 0 & M/A \end{pmatrix} \begin{pmatrix} I & A^{-1}B \\ 0 & I \end{pmatrix}, \quad (\text{A159})$$

where

$$M/A \equiv D - CA^{-1}B \quad (\text{A160})$$

is called the Schur complement. The first thing to notice is that

$$\det M = \det A \det(M/A), \quad (\text{A161})$$

and the second thing is that the inverse of M is

$$M^{-1} = \begin{pmatrix} A^{-1} + A^{-1}B(M/A)^{-1}CA^{-1} & -A^{-1}B(M/A)^{-1} \\ -(M/A)^{-1}CA^{-1} & (M/A)^{-1} \end{pmatrix}. \quad (\text{A162})$$

In the case at hand, the block form of M_{T+dT} relative to the outer row and column, given in Equation (A155), can be written in the matrix form of Equation (A158) by identifying

$$A = M_T, \quad (\text{A163})$$

$$B = -\kappa\sqrt{dT}|+_T\rangle\langle N|, \quad (\text{A164})$$

$$C = -\kappa\sqrt{dT}|N\rangle\langle+_T|, \quad (\text{A165})$$

$$D = (1 - \kappa dT)|N\rangle\langle N|. \quad (\text{A166})$$

The Schur complement is the 1×1 matrix

$$M_{T+dT}/M_T = (1 - \kappa dT - \kappa^2 dT \langle+_T|M_T^{-1}|+_T\rangle)|N\rangle\langle N| = [1 - \kappa dT(1 + n_T)]|N\rangle\langle N|. \quad (\text{A167})$$

One sees immediately how the determinant advances by one increment,

$$\det M_{T+dT} = \det M_T(1 - \kappa dT(1 + n_T)), \quad (\text{A168})$$

and thus obtains the ODE

$$\frac{1}{\kappa} \frac{d \ln \det M_T}{dT} = -(1 + n_T), \quad (\text{A169})$$

which agrees with the differential equation in Equation (292) for the normalization factor $\mathcal{N}_T = 1/\det M_T$. Further, one can work out from Equation (A162) how the inverse increments,

$$\begin{aligned} M_{T+dT}^{-1} &= M_T^{-1} + \kappa^2 dT M_T^{-1}|+_T\rangle\langle+_T|M_T^{-1} \\ &\quad + \kappa\sqrt{dT}(M_T^{-1}|+_T\rangle\langle N| + |N\rangle\langle+_T|M_T^{-1}) \\ &\quad + [1 + \kappa dT(1 + \kappa\langle+_T|M_T^{-1}|+_T\rangle)]|N\rangle\langle N|. \end{aligned} \quad (\text{A170})$$

where we keep only the leading-order term in dT in each block.

Everything is set now to determine how the three moments increment,

$$n_{T+dT} = \kappa\langle+_T+dT|M_{T+dT}^{-1}|+_T+dT\rangle = n_T + \kappa dT(1 - n_T)^2, \quad (\text{A171})$$

$$m_{T+dT} = \kappa\langle-_T+dT|M_{T+dT}^{-1}|-_T+dT\rangle = m_T + \kappa dT(q_T + e^{-2\kappa T})^2, \quad (\text{A172})$$

$$q_{T+dT} = \kappa\langle-_T+dT|M_{T+dT}^{-1}|+_T+dT\rangle = q_T + \kappa dT(-q_T(1 - n_T) + e^{-2\kappa T}(1 + n_T)), \quad (\text{A173})$$

and these lead immediately to the three (coupled) Riccati equations,

$$\frac{1}{\kappa} \frac{dn_T}{dT} = (1 - n_T)^2, \quad (\text{A174})$$

$$\frac{1}{\kappa} \frac{dm_T}{dT} = (q_T + e^{-2\kappa T})^2, \quad (\text{A175})$$

$$\frac{1}{\kappa} \frac{dq_T}{dT} = -q_T(1 - n_T) + e^{-2\kappa T}(1 + n_T), \quad (\text{A176})$$

which are repeated in the main text as Equations (313)–(315).

References

1. Klauder, J.R. The action option and a Feynman quantization of spinor fields in terms of ordinary c-numbers. *Ann. Phys.* **1960**, *11*, 123–168. [\[CrossRef\]](#)
2. Perelomov, A. *Generalized Coherent States and Their Applications*; Texts and Monographs in Physics; Springer: Berlin/Heidelberg, Germany, 1986.
3. Zhang, W.-M.; Fang, D.H.; Gilmore, R. Coherent states: Theory and some applications. *Rev. Mod. Phys.* **1990**, *62*, 867–927. [\[CrossRef\]](#)
4. Brif, C.; Mann, A. Phase-space formulation of quantum mechanics and quantum-state reconstruction for physical systems with Lie-group symmetries. *Phys. Rev. A* **1999**, *59*, 971–987. [\[CrossRef\]](#)
5. Glauber, R.J. Photon correlations. *Phys. Rev. Lett.* **1963**, *10*, 84–86. [\[CrossRef\]](#)
6. Sudarshan, E.C.G. Equivalence of semiclassical and quantum mechanical descriptions of statistical light beams. *Phys. Rev. Lett.* **1963**, *10*, 277–279. [\[CrossRef\]](#)
7. Glauber, R.J. Coherent and incoherent states of the radiation field. *Phys. Rev.* **1963**, *131*, 2766–2788. [\[CrossRef\]](#)
8. Glauber, R.J. Optical coherence and photon statistics. In *Quantum Optics and Electronics*; DeWitt, C., Blandin, A., Cohen-Tannoudji, C., Eds.; Gordon and Breach: Philadelphia, PA, USA, 1965; pp. 65–85.
9. Massar, S.; Popescu, S. Optimal extraction of information from finite quantum ensembles. *Phys. Rev. Lett.* **1995**, *74*, 1259–1263. [\[CrossRef\]](#)
10. Wiener, N. The average of an analytic functional. *Proc. Natl. Acad. Sci. USA* **1921**, *7*, 253–260. [\[CrossRef\]](#)
11. Wiener, N. The average of an analytic functional and the Brownian movement. *Proc. Natl. Acad. Sci. USA* **1921**, *7*, 294–298. [\[CrossRef\]](#)
12. Wiener, N. The average value of a functional. *Proc. Lond. Math. Soc.* **1924**, *s2-22*, 454–467. [\[CrossRef\]](#)
13. Feynman, R.P. *Feynman's PhD Thesis: A New Approach to Quantum Theory*; Brown, L.M., Ed.; World Scientific: Singapore, 2005.
14. Feynman, R.P. Space-time approach to non-relativistic quantum mechanics. *Rev. Mod. Phys.* **1948**, *20*, 367–387. [\[CrossRef\]](#)
15. Kac, M. On distributions of certain Wiener functionals. *Trans. Am. Math. Soc.* **1949**, *65*, 1–13. [\[CrossRef\]](#)
16. Kac, M. *Probability and Related Topics in Physical Sciences*; Lectures in Applied Mathematics; American Mathematical Society: Providence, RI, USA, 1959; Volume I.A.
17. Feynman, R.P.; Hibbs, A.R.; Styer, D.F. *Quantum Mechanics and Path Integrals (Emended Edition)*; Dover: Mineola, NY, USA, 2010.
18. Bourbaki, N. *Elements of Mathematics. Integration II, Chapters 7–9*; Springer: Berlin/Heidelberg, Germany, 2004.
19. Chaichian, M.; Demichev, A. *Path Integrals in Physics. Volume I: Stochastic Processes and Quantum Mechanics*; CRC Press: Boca Raton, FL, USA; Taylor & Francis Group: Abingdon, UK, 2001.
20. Matsubara, T. A new approach to quantum-statistical mechanics. *Prog. Theor. Phys.* **1955**, *14*, 351–378. [\[CrossRef\]](#)
21. Ludwig, G. *Foundations of Quantum Mechanics I*; Texts and Monographs in Physics; Springer: Berlin/Heidelberg, Germany, 1983.
22. Ludwig, G. *Foundations of Quantum Mechanics II*; Texts and Monographs in Physics; Springer: Berlin/Heidelberg, Germany, 1985.
23. Schwinger, J. The algebra of microscopic measurement. *Proc. Natl. Acad. Sci. USA* **1959**, *45*, 1542–1553. [\[CrossRef\]](#)
24. Wigner, E.P. The problem of measurement. *Am. J. Phys.* **1963**, *31*, 6–15. [\[CrossRef\]](#)
25. Jauch, J.M.; Piron, C. Generalized localizability. *Helv. Phys. Acta* **1967**, *40*, 559–570.
26. Davies, E.B.; Lewis, J.T. An operational approach to quantum probability. *Commun. Math. Phys.* **1970**, *17*, 239–260. [\[CrossRef\]](#)
27. Davies, E.B. *Quantum Theory of Open Systems*; Academic Press: Cambridge, MA, USA, 1976.
28. Lindblad, G. On the generators of quantum dynamical semigroups. *Commun. Math. Phys.* **1976**, *48*, 119–130. [\[CrossRef\]](#)
29. Kraus, K. *States, Effects, and Operations: Fundamental Notions of Quantum Theory*; Lecture Notes in Physics; Böhm, A., Dollard, J.D., Wootters, W.K., Eds.; Springer: Berlin/Heidelberg, Germany, 1983; Volume 190.
30. Peres, A. *Quantum Theory: Concepts and Methods*; Kluwer Academic: Amsterdam, The Netherlands, 1993.
31. Nielsen, M.A.; Chuang, I.L. *Quantum Computation and Quantum Information*; Cambridge University Press: Cambridge, UK, 2000.
32. Wiseman, H.M.; Milburn, G.J. *Quantum Measurement and Control*; Cambridge University Press: Cambridge, UK, 2009.
33. Barchielli, A.; Belavkin, V.P. Measurements continuous in time and *a posteriori* states in quantum mechanics. *J. Phys. A Math. Gen.* **1991**, *24*, 1495–1514. [\[CrossRef\]](#)
34. Wiseman, H.M. Quantum Trajectories and Feedback. Ph.D. Thesis, University of Queensland, Brisbane, QLD, Australia, 1994.

35. Goetsch, P.; Graham, R. Linear stochastic wave equations for continuously measured quantum systems. *Phys. Rev. A* **1994**, *50*, 5242–5255. [[CrossRef](#)]
36. Wiseman, H.M. Quantum trajectories and quantum measurement theory. *Quantum Semiclass. Opt. J. Eur. Opt. Soc. Part B* **1996**, *8*, 205. [[CrossRef](#)]
37. Jackson, C.S. The photodetector, the heterodyne instrument, and the principle of instrument autonomy. *arXiv* **2022**, arXiv:2210.11100.
38. Shojaei, E.; Jackson, C.S.; Riofrío, C.A.; Kalev, A.; Deutsch, I.H. Optimal pure-state qubit tomography via sequential weak measurements. *Phys. Rev. Lett.* **2018**, *121*, 130404. [[CrossRef](#)] [[PubMed](#)]
39. Jackson, C.S.; Caves, C.M. How to perform the coherent measurement of a curved phase space by continuous isotropic measurement. I. Spin and the Kraus-operator geometry of $SL(2, \mathbb{C})$. *arXiv* **2021**, arXiv:2107.12396.
40. Jackson, C.S.; Caves, C.M. Simultaneous measurements of noncommuting observables. Positive transformations and instrumental Lie groups. *arXiv* **2023**, arXiv:2306.06167.
41. Barchielli, A.; Lanz, L.; Prosperi, G.M. A model for the macroscopic description and continual observations in quantum mechanics. *Nuovo Cim. B* **1982**, *72*, 79–121. [[CrossRef](#)]
42. Beard, M. *SPQR: A History of Ancient Rome*; Liveright: New York, NY, USA, 2015.
43. Haar, A. Der Massbegriff in der Theorie der kontinuierlichen Gruppen. *Ann. Math. Second Ser.* **1933**, *34*, 147–169. [[CrossRef](#)]
44. von Neumann, J. *Invariant Measures*; American Mathematical Society: Providence, RI, USA, 1999.
45. Nachbin, L. *The Haar Integral*; D. Van Nostrand Company, Inc.: New York, NY, USA, 1965.
46. Knapp, A.W. *Lie Groups, Lie Algebras, and Cohomology*; Mathematical Notes 34; Princeton University Press: Princeton, NJ, USA, 1988.
47. Frankel, T. *The Geometry of Physics: An Introduction*, 3rd ed.; Cambridge University Press: Cambridge, UK, 2012.
48. Kitaev, A. Notes on $\widetilde{SL}(2, \mathbb{R})$ representations. *arXiv* **2018**, arXiv:1711.08169.
49. Gilmore, R. *Lie Groups, Lie Algebras, and Some of Their Applications*; Dover: Mineola, NY, USA, 2002.
50. Helgason, S. *Differential Geometry, Lie Groups, and Symmetric Spaces*; Pure and Applied Mathematics; Academic Press: Cambridge, MA, USA, 1978.
51. Barut, A.O.; Raczka, R. *Theory of Group Representations and Applications*, 2nd revised ed.; World Scientific: Singapore, 1986.
52. Knapp, A.W. *Lie Groups Beyond an Introduction*, 2nd ed.; Progress in Mathematics; Birkhäuser: Cham, Switzerland, 2002; Volume 140.
53. Harish-Chandra. Representations of Semisimple Lie Groups, V. *Am. J. Math.* **1956**, *78*, 1–41. [[CrossRef](#)]
54. Knapp, A.W. *Representation Theory of Semisimple Lie Groups: An Overview Based on Examples*; Princeton Landmarks in Mathematics; Princeton University Press: Princeton, NJ, USA, 1986.
55. Želobenko, D.P. *Compact Lie Groups and Their Representations*; Translations of Mathematical Monographs; American Mathematical Society: Providence, RI, USA, 1973; Volume 40.
56. Itô, K. Stochastic differential equations in a differentiable manifold. *Nagoya Math. J.* **1950**, *1*, 35–47. [[CrossRef](#)]
57. Itô, K.; McKean, H.P., Jr. *Diffusion Processes and Their Sample Paths: Reprint of the 1974 Edition*; Springer: Berlin/Heidelberg, Germany, 1996.
58. Chirikjian, G.S. *Stochastic Models, Information Theory, and Lie Groups, Volume I: Classical Results and Geometric Methods*; Applied and Numerical Harmonic Analysis; Birkhäuser: Cham, Switzerland, 2009.
59. Gardiner, C. *Stochastic Methods: A Handbook for the Natural and Social Sciences*, 4th ed.; Springer Series in Synergetics; Springer: Berlin/Heidelberg, Germany, 2009.
60. Gardiner, C. *Elements of Stochastic Methods*; AIP Publishing: Melville, NY, USA, 2021.
61. Uhlenbeck, G.E.; Ornstein, L.S. On the theory of Brownian motion. *Phys. Rev.* **1930**, *36*, 823–841. [[CrossRef](#)]
62. Srinivas, M.D.; Davies, E.B. Photon counting probabilities in quantum optics. *Opt. Acta* **1981**, *28*, 981–996. [[CrossRef](#)]
63. Jacobs, K.; Knight, P.L. Linear quantum trajectories: Applications to continuous projection measurements. *Phys. Rev. A* **1998**, *57*, 2301–2310. [[CrossRef](#)]
64. Jacobs, K.; Steck, D.A. A straightforward introduction to continuous quantum measurement. *Contemp. Phys.* **2006**, *47*, 279–303. [[CrossRef](#)]
65. Martin, L.; Motzoi, F.; Li, H.; Sarovar, M.; Whaley, K.B. Deterministic generation of remote entanglement with active quantum feedback. *Phys. Rev. A* **2015**, *92*, 062321. [[CrossRef](#)]
66. von Neumann, J. *Mathematical Foundations of Quantum Mechanics: New Edition*; Wheeler, N.A., Ed.; Princeton University Press: Princeton, NJ, USA, 2018.
67. Arthurs, E.; Kelly, J.L., Jr. On the simultaneous measurement of a pair of conjugate observables. *Bell Syst. Tech. J.* **1965**, *44*, 725–729. [[CrossRef](#)]
68. Braunstein, S.L.; Caves, C.M.; Milburn, G.J. Interpretation for a positive P representation. *Phys. Rev. A* **1991**, *43*, 1153–1159. [[CrossRef](#)]
69. Poincaré, H. *Papers on Topology: Analysis Situs and Its Five Supplements (History of Mathematics)*; Stillwell, J., Translator; American Mathematical Society: Providence, RI, USA; The London Mathematical Society: London, UK, 2010; Volume 37.
70. Weyl, H. *The Concept of a Riemann Surface*, 3rd ed.; Dover: Mineola, NY, USA, 2009.
71. Bourbaki, N. *Elements of Mathematics. Lie Groups and Lie Algebras, Chapters 1–3*; Springer: Berlin/Heidelberg, Germany, 1989.
72. Pontrjagin, L. *Topological Groups*; Princeton University Press: Princeton, NJ, USA, 1946.

73. Ali, S.T.; Antoine, J.-P.; Gazeau, J.-P. *Coherent States. Wavelets. and Their Generalizations*, 2nd ed.; Theoretical and Mathematical Physics; Springer: Berlin/Heidelberg, Germany, 2014.
74. Dixmier, J. *Enveloping Algebras*; Graduate Studies in Mathematics; American Mathematical Society: Providence, RI, USA, 1996; Volume 11.
75. Chern, S.-S.; Chevalley, C. Élie Cartan and his mathematical work. *Bull. Am. Math. Soc.* **1952**, *58*, 157–216. [[CrossRef](#)]
76. Howe, R. “Harish-Chandra: 1923–1983”, Biographical Memoirs of the US National Academy of Sciences. 2011. Available online: www.nasonline.org/publications/biographical-memoirs/ (accessed on 1 June 2023).

Disclaimer/Publisher’s Note: The statements, opinions and data contained in all publications are solely those of the individual author(s) and contributor(s) and not of MDPI and/or the editor(s). MDPI and/or the editor(s) disclaim responsibility for any injury to people or property resulting from any ideas, methods, instructions or products referred to in the content.