# Optimal Control With Noisy Time

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*Abstract*—This paper examines stochastic optimal control problems in which the state is perfectly known, but the controller's measure of time is a stochastic process derived from a strictly increasing Lévy process. We provide dynamic programming results for continuous-time finite-horizon control and specialize these results to solve a noisy-time variant of the linear quadratic regulator problem and a portfolio optimization problem with random trade activity rates. For the linear quadratic case, the optimal controller is linear and can be computed from a generalization of the classical Riccati differential equation.

*Index Terms*—Optimal control, stochastic optimal control, stochastic systems, uncertain time.

### I. INTRODUCTION

**E** FFECTIVE feedback control often requires accurate timekeeping. For example, finite-horizon optimal control problems generally result in policies that are time-varying functions of the state. However, chronometry is imperfect and thus feedback laws are inevitably applied at incorrect times. Little appears to be known about the consequences of imperfect timing on control [1]–[3]. This paper addresses optimal control with temporal uncertainty.

A stochastic process can be *time-changed* by replacing its time index by a monotonically increasing stochastic process [4]. Time-changed stochastic processes arise in finance, since changing the time index to a measure of economically relevant events, such as trades, can improve modeling [5]–[7]. This new time index is, however, stochastic with respect to "calendar" time.

We suspect that similar notions of stochastic time changing may facilitate the study of time estimation and movement control in the nervous system. Biological timing is subject to noise and environmental perturbation [8]. Furthermore, humans rationally exploit the statistics of their temporal noise during simple timed movements, such as button pushing [9] and pointing [10]. To analyze more complex movements, a theory of feedback control that compensates for temporal noise seems desirable.

Within control, the most closely related work to the present paper deals with analysis and synthesis of systems with uncertain sampling times. The study of uncertain sampling times has

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a long history in control [11], and is often motivated by problems of clock jitter [12], [13] or network delays [14]. In these works, control inputs are sampled at known times and held over unknown intervals. To derive the dynamic programming principle in this paper, system behavior is analyzed for control inputs held over random intervals, bearing similarity to optimal control with random sampling [15]. Fundamentally, however, studies of sampling uncertainty assumes that an accurate clock can measure the sample times; the present work relaxes this assumption.

Other aspects of imperfect timing have been addressed in control research to a more limited extent. For example, the importance of synchronizing clocks in distributed systems seems clear [16], [17], but more work is needed to understand the implications of asynchronous clock behavior on common control issues, such as stability [18] and optimal performance [19].

This paper focuses on continuous-time stochastic optimal control with perfect state information, but a stochastically timechanged control process. Dynamic programming principles for general nonlinear stochastic control problems are derived, based on extensions of the classical Hamilton-Jacobi-Bellman equation. The results apply to a wide class of stochastic time changes given by strictly increasing Lévy processes. The dynamic programming principles are then specialized to give explicit solutions to a time-changed version of the finite-horizon linear quadratic regulator and a portfolio optimization problem.

Section II defines the notation used in the paper, states the necessary facts about Lévy, and defines the class of noisy clock models used. The main results on time-changed diffusions and optimal control are given in Section III. The results are proved in Section IV, with supplementary arguments given in the Appendices. Section V discusses variations on the problem and potential applications, while Section VI concludes.

#### **II. PRELIMINARIES**

After establishing notation and reviewing Lévy processes, this section culminates in the construction of Lévy-processbased clock models upon which the remainder of the theory of this paper is built.

#### A. Notation

The norm symbol,  $\|\cdot\|$ , is used to denote the Euclidean norm for vectors and the Frobenius norm for matrices.

For a set S, its closure is denoted by S.

The spectrum of matrix A is denoted by  $\operatorname{spec}(A)$ .

The Kronecker product is denoted by  $\otimes$ , while the Kronecker sum is denoted by  $\oplus$ :  $A \oplus B = A \otimes I + I \otimes B$ .

The vectorization operation of stacking the columns of a matrix is denoted by vec.

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A function  $h : \mathbb{R} \times \mathbb{R}^n \to \mathbb{R}$  is in  $\mathcal{C}^{1,2}$  if h(s, x) is continuously differentiable in s, twice continuously differentiable in x. The function h is said to satisfy a *polynomial growth condition*, if in addition, there are constants K and q such that

$$\max\left\{ \left|h(s,x)\right|, \left|\frac{\partial h(s,x)}{\partial s}\right|, \left|\frac{\partial h(s,x)}{\partial x_i}\right|, \left|\frac{\partial^2 h(s,x)}{\partial x_i \partial x_j}\right| \right\} \le K \left(1 + \|x\|^q\right)$$

for  $i, j = 1, \ldots n$ , and all  $x \in \mathbb{R}^n$ . In this case,  $h \in \mathcal{C}_p^{1,2}$  is written.

Stochastic processes will be denoted as  $\zeta_t$ ,  $X_s$ , etc., with time indices as subscripts. Occasionally, processes with nested subscripts will be written with parentheses, e.g.  $\zeta_{\tau_s} = \zeta(\tau_s)$ . Similarly, the elements of a stochastic vector will be denoted as  $X_1(s)$ .

The left limit of a process  $X_t$  will be denoted by  $X_{t^-} = \lim_{s \uparrow t} X_s$ .

Functions that are right-continuous with left-limits will be called *càdlàg*, while functions that are left-continuous with right-limits will be called *càglàd*.

#### B. Background on Lévy Processes

Basic notions from Lévy processes required to define the general class of clock models are now reviewed. The definitions and results can be found in [20].

A real-valued stochastic process  $Z_s$  is called a *Lévy process* if

- $Z_0 = 0$  almost surely (a.s.).
- $Z_s$  has independent, stationary increments: If  $0 \le r \le s$ , then  $Z_r$  and  $Z_s - Z_r$  are independent and  $Z_s - Z_r$  has the same distribution as  $Z_{s-r}$ .
- Z<sub>s</sub> is stochastically continuous: For all a > 0 and all s ≥ 0, lim<sub>r→s</sub> P(|Z<sub>s</sub> Z<sub>r</sub>| > a) = 0.

It will be assumed that Lévy processes in this paper are right-continuous with left-sided limits, i.e., they are càdlàg. No generality is lost since, for every Lévy process,  $Z_t$ , there is a càdlàg Lévy process,  $\tilde{Z}_t$ , such that  $Z_t = \tilde{Z}_t$  for almost all t.

Some of the technical arguments rely on Poisson random measures, which will now be defined. Let  $\mathcal{B}$  be the Borel subsets of  $\mathbb{R}$  and let  $(\Omega, \Sigma, \mathbb{P})$  be a probability space. A *Poisson random measure* is a function  $N : [0, \infty) \times \mathcal{B} \times \Omega \to \mathbb{N} \cup \{\infty\}$ , such that

- For all  $s \ge 0$  and  $\omega \in \Omega$ ,  $N(s, \cdot, \omega)$  is a measure.
- For all disjoint Borel subsets  $A, B \in \mathcal{B}$  such that  $0 \notin \overline{A}$  and  $0 \notin \overline{B}$ ,  $N(\cdot, A, \cdot)$  and  $N(\cdot, B, \cdot)$  are independent Poisson processes.

Typically, the  $\omega$  argument will be dropped, and it will be implicitly understood that N(s, A) denotes a measure-valued stochastic process.

The following relationship between Lévy processes and Poisson random measures will be used in several arguments. For a Lévy process,  $Z_s$ , with jumps denoted by  $\Delta Z_s$ , there is a Poisson random measure that counts the number of jumps into each Borel set A with  $0 \notin \overline{A}$ 

$$N(s, A) = |\{\Delta Z_r \in A : 0 \le r \le s\}|.$$

Here, |S| denotes the cardinality of a finite set, S.

Subordinators: A monotonically increasing Lévy process,  $\tau_s$ , is called a *subordinator*. The following properties of subordinators will be used throughout the paper.

• Laplace Exponent: There is function,  $\psi$ , called the *Laplace exponent*, defined by

$$\psi(z) = bz + \int_{0}^{\infty} (1 - e^{-zt})\lambda(dt) \tag{1}$$

such that

$$\mathbb{E}[e^{-z\tau_s}] = e^{-s\psi(z)} \quad \text{for all} \quad z \ge 0.$$
 (2)

Here  $b \ge 0$  and the measure satisfies  $\int_0^\infty \min\{t, 1\}\lambda(dt) < \infty$ . The measure  $\lambda$  is called a *Lévy measure*. The pair  $(b, \lambda)$  is called the *characteristics* of  $\tau_s$ .

• Lévy-Itô Decomposition: There is a Poisson random measure N such that

$$\tau_s = bs + \int_0^\infty t N(s, dt)$$

Furthermore, if  $A \subset (0, \infty)$  is a Borel set such that  $0 \notin \overline{A}$ , then  $\mathbb{E}[N(1, A)] = \lambda(A)$ .

The function,  $\psi$ , is called the Laplace exponent because (2) is the Laplace transform of the distribution of  $\tau_s$ .

For control problems, simpler formulas will often result from replacing  $\psi$  with the function  $\beta(z) = -\psi(-z)$ . Note then, that  $\beta$  has the form

$$\beta(z) = bz + \int_{0}^{\infty} (e^{zt} - 1)\lambda(dt).$$
(3)

Define  $r_{\max}$  by

$$r_{\max} = \sup\left\{r: \int_{1}^{\infty} e^{rt}\lambda(dt) < \infty\right\}$$

and define the domain of  $\beta$  as

$$\operatorname{dom}(\beta) = \{ z \in \mathbb{C} : \operatorname{Re} z < r_{\max} \}.$$

Note that  $\int_{1}^{\infty} \lambda(dt) < \infty$  implies that  $r_{\max} \in [0, \infty]$ .

The function  $\beta$  is used to construct optimal solutions for the linear quadratic problem, as well as the portfolio problem below. The main properties are given in the following lemma, which is proved in Appendix B.

Lemma 1: For all  $z \in dom(\beta)$ , the function  $\beta$  is analytic at z, and

$$\mathbb{E}[e^{z\tau_s}] = e^{s\beta(z)}.$$
(4)

Furthermore, if A is a square matrix with spec $(A) \subset \text{dom}\beta$ , then

$$\beta(A) = bA + \int_{0}^{\infty} (e^{At} - I)\lambda(dt)$$
(5)

is well defined and

$$\mathbb{E}[e^{A\tau_s}] = e^{s\beta(A)}.$$
(6)

Since  $\beta$  is analytic, several methods exist for numerically computing the matrices  $\beta(A)$  [21]. In special cases, as discussed below,  $\beta(A)$  may be computed using well-known matrix computation methods.

*Example 1:* The simplest non-trivial subordinator is the homogeneous Poisson process  $N_t$ , which is characterized by

$$\mathbb{P}(N_t = k) = e^{-\gamma t} \frac{(\gamma t)^k}{k!}$$

where  $\gamma > 0$  is called the rate constant. Its Laplace exponent is given by  $\psi(z) = \gamma - \gamma e^{-z}$ , which is found by computing the expected value directly. The characteristics are  $(0, \gamma \delta(t-1))$ . In this case, dom $(\beta) = \mathbb{C}$ , and  $\beta(A) = \gamma e^A - \gamma I$ , which can be computed from the matrix exponential.

*Example 2:* The gamma subordinator, which is often used to model "business time" in finance [22], [23], has increments distributed as gamma random variables. It has Laplace exponent  $\psi(z) = \delta \log(1 + z/\gamma)$  with characteristics b = 0 and  $\lambda(dt) = \delta e^{-\gamma t} t^{-1} dt$ . Thus  $\beta(z) = -\delta \log(1 - z/\gamma)$ ,  $\operatorname{dom}(\beta) = \{z \in \mathbb{C} : \operatorname{Re} z < \gamma\}$ , and matrix function  $\beta(A) = -\delta \log(I - \gamma^{-1}A)$  may be computed from the matrix logarithm.<sup>1</sup>

Why Lévy Processes?: In the next subsection, the clock model in this paper will be constructed from a subordinator  $\tau_s$ . The motivation for using Lévy processes will be explained. Consider a continuous-time noisy clock,  $c_s$  which is sampled with period  $\delta$ . A natural model might take the form

$$c_{\delta(k+1)} = c_{\delta k} + \delta + n(k,\delta) \tag{7}$$

where  $n(j, \delta)$  are random variables. In this case, the clock increments consist of a deterministic step of magnitude  $\delta$  plus a random term.

If  $c_s$  is a Lévy process, then by definition, all of the increments  $c_{\delta(k+1)} - c_{\delta k}$  are independent and identically distributed. Thus, the decomposition in (7) holds with  $n(k, \delta) = c_{\delta(k+1)} - c_{\delta k} - \delta$ . If  $c_s$  were not a Lévy process, then (7) may hold for some particular  $\delta$ , but there might be another period,  $\delta' < \delta$ , for which the decomposition fails. The Lévy process assumption will guarantee that the clocks are well-behaved when taking continuous time limits (i.e.,  $\delta \downarrow 0$ ).

## C. Clock Models

Throughout the paper, t will denote the time index of the plant dynamics, while s will denote the value of clock available to the controller. Often, t and s will be called *plant time* and *controller time*, respectively. The interpretation of s and t varies depending on context. In biological motor control, t would denote real time, since the limbs obey Newtonian mechanics with respect to real-time, while s would denote the internal representation of time. For the portfolio problem studied in Section III-B, an opposite interpretation holds. Here, the controller (an investor) can accurately measure calendar time, but price dynamics are simpler with respect a different index, "business time," which represents the progression of economic events [5]–[7]. Thus, s would denote calendar time, while t would denote business time, which might not be observable.



Fig. 1. (a) The inverse Gaussian subordinator,  $\tau_s$ , with  $\gamma = \delta = 2$ . The process was simulated by generating independent inverse Gaussians using the method from [25]. (b) The inverse process,  $\zeta_t$ . Note that the graph of  $\zeta_t$  can be found from the graph of  $\tau_s$  by simply switching the axes.

The relationship between s and t will be described stochastically. Let  $\tau_s$  be a strictly increasing subordinator. In other words, if s < s' then  $\tau_s < \tau_{s'}$  a.s. (Note that any subordinator can be made to be strictly increasing by adding a drift term bswith b > 0.) The process  $\tau_s$  will be interpreted as the amount of plant time that has passed when the controller has measured s units of time. The process  $\zeta_t$  will be an inverse process that describes how much time the controller measures over t units of plant time. Formally,  $\zeta_t$  is defined by

$$\zeta_t = \inf\{\sigma : \tau_\sigma \ge t\}. \tag{8}$$

We claim that  $\zeta(\tau_s) = s$  a.s. Indeed,  $\zeta(\tau_s) = \inf \{ \sigma : \tau_\sigma = \tau_s \}$ , by definition. Since  $\tau_s$  is right continuous and strictly increasing, a.s., the claim follows.

*Example 3:* The case of no temporal uncertainty corresponds to  $\tau_s = s$  and  $\zeta_t = t$ . The Laplace exponent of  $\tau_s$  is computed directly as  $\psi(z) = z$  and the characteristics are (1, 0). Here dom $(\beta) = \mathbb{C}$ .

*Example 4:* A more interesting temporal noise model, also used as a "business time" model [24], is the inverse Gaussian subordinator. Fix  $\gamma > 0$  and  $\delta > 0$ . Let  $C_t = \gamma t + W_t$ , where  $W_t$  is a standard Brownian motion. The inverse Gaussian subordinator is given by

$$\tau_s = \inf\{t : C_t = \delta s\}$$

with Laplace exponent  $\psi(z) = \delta(\sqrt{\gamma^2 + 2z} - \gamma)$ . Here b = 0 and  $\lambda$  is given by

$$\lambda(dt) = \frac{\delta}{\sqrt{2}\Gamma\left(\frac{1}{2}\right)} e^{-\frac{1}{2}\gamma^2 t} t^{-\frac{3}{2}} dt$$

where  $\Gamma$  is the gamma function. Here, dom $(\beta)$  corresponds to Re  $z < \gamma^2/2$  and  $\beta(A) = \delta(\gamma I - \sqrt{\gamma^2 I - 2A})$ , which can be computed from the matrix square root. It can be shown that the inverse process is given by

$$\zeta_t = \sup\{\delta^{-1}C_\sigma : 0 \le \sigma \le t\}.$$

See Fig. 1.

In the preceding example, the process  $\tau_s$  has jumps, but the inverse,  $\zeta_t$ , is continuous. The next proposition generalizes this observation for any strictly increasing subordinator,  $\tau_s$ .

*Proposition 1:* The process  $\zeta_t$  is continuous almost surely.

*Proof:* Fix  $\epsilon > 0$  and  $t \ge 0$ . Set  $s = \zeta_t$ . Strict monotonicity of  $\tau_s$  implies that  $[\tau_{\max\{s-\epsilon,0\}}, \tau_{s+\epsilon}]$  is a nonempty interval, a.s. The inverse property of  $\zeta_t$  implies (almost surely) that  $t \in [\tau_{\max\{s-\epsilon,0\}}, \tau_{s+\epsilon}]$  and  $\zeta_{t'} \in [\max\{s-\epsilon,0\}, s+\epsilon]$  for all  $t' \in [\tau_{\max\{s-\epsilon,0\}}, \tau_{s+\epsilon}]$ .

<sup>&</sup>lt;sup>1</sup>When spec(M)  $\subset \{z \in \mathbb{C} : \operatorname{Re} z > 0\}$ ,  $\log(M)$  refers to the principle logarithm: the unique matrix such that  $e^{\log(M)} = M$  and  $\operatorname{spec}(\log(M)) \subset \{z \in \mathbb{C} : -\pi < \operatorname{Im} z < \pi\}$ . See [21].

#### **III. MAIN RESULTS**

This section presents the main results of the paper. First, given an Itô process,  $Y_t$ , a representation of the time-changed process  $X_s = Y(\tau_s)$  as a semimartingale with respect to controller time, s, is derived. This representation is then used to derive a general dynamic programming principle for finitehorizon control problems with noisy clocks. The time horizon is taken with respect to the controller's clock. As an example, the dynamic programming principle is used to solve a simple portfolio optimization problem under random trade activity rates. Finally, the dynamic programming method is used to solve a noisy-time variant of the linear quadratic regulator problem. All proofs are given in Section IV.

#### A. Time-Changed Stochastic Processes

The results of this subsection shows how to reparametrize the plant dynamics from a process over plant time t to a process over controller time s. This reparametrization will be used to pose optimal control problems with respect to controller time s. The results are proved in Section IV-A.

Let  $W_t$  be a Brownian motion with  $\mathbb{E}[W_t W_t^{\mathsf{T}}] = tI$ . Let Y be a stochastic process defined by

$$dY_t = F_t dt + G_t dW_t \tag{9}$$

where  $F_t$  and  $G_t$  are  $\mathcal{F}_t^W$  predictable processes, where  $(\mathcal{F}_t^W)_{t\geq 0}$  is the  $\sigma$ -algebra generated by  $W_t$ . Furthermore, assume that  $F_t$  and  $G_t$  are left-continuous with right-sided limits.

For a strictly increasing subordinator,  $\tau_s$ , let  $\mathcal{F}^{\tau,W}$  =  $(\mathcal{F}^{\tau,W}_s)_{s\geq 0}$  be the smallest filtration such that for all  $r\in[0,s]$ and all  $t \in [0, \tau_s]$  both  $\tau_r$  and  $W_t$  are measurable.

*Theorem 1:* Let  $\tau_s$  be a subordinator characterized by  $(b, \lambda)$ . If the terms of (9) satisfy

- ∫<sub>0</sub><sup>τ<sub>S</sub></sup> ||F<sub>t</sub>||dt < ∞ almost surely;</li>
   𝔼[∫<sub>0</sub><sup>τ<sub>S</sub></sup> ||G<sub>t</sub>||<sup>2</sup>dt] < ∞;</li>

then the time-changed process  $X_s = Y(\tau_s)$  is an  $\mathcal{F}^{\tau,W}$  semimartingale given by

$$X_{s} = X_{0} + b \int_{0}^{s} F(\tau_{r^{-}}) dr + \sqrt{b} \int_{0}^{s} G(\tau_{r^{-}}) d\tilde{W}_{r} + \sum_{0 \le r \le s} \left( \int_{\tau_{r^{-}}}^{\tau_{r}} F_{t} dt + \int_{\tau_{r^{-}}}^{\tau_{r}} G_{t} dW_{t} \right).$$
(10)

Here  $\tilde{W}_s$  is an  $\mathcal{F}^{\tau,W}$ -measurable Brownian motion defined by

$$\sqrt{b}\tilde{W}_s = W(\tau_s) - \sum_{0 \le r \le s} \left( W(\tau_r) - W(\tau_{r^-}) \right)$$

satisfying  $b\mathbb{E}[\tilde{W}_s \tilde{W}_s^{\mathsf{T}}] = bsI$ .

#### B. Dynamic Programming

This subsection introduces the general control problem studied in this paper. First, the basic notions of controlled timechanged diffusions and admissible systems are defined. Then, the optimal control problem in controller time is presented, and the associated dynamic programming verification theorem is stated. The concepts in this subsection are based on the general theory of controlled Markov processes [26]. A nonlinear portfolio problem is solved at the end of the subsection.

Controlled Time-Changed Diffusions: Consider a controlled diffusion

$$dY_{t} = F(\zeta_{t}, Y_{t^{-}}, U(\zeta_{t})) dt + G(\zeta_{t}, Y_{t^{-}}, U(\zeta_{t})) dW_{t}$$
(11)

with state Y and input U. Recall that  $\zeta_t$  is defined in (8) as the inverse process of a subordinator,  $\tau_s$ . Let  $X_s$  denote the time-changed process,  $X_s = Y(\tau_s)$ . The processes,  $X_s$  is thus a time-changed controlled diffusion.

Admissible Systems: For  $s \ge 0$ , let  $\mathcal{F}_s^{\zeta,X}$  be the  $\sigma$ -algebra generated by  $(s, X_s)$ , and let  $\mathcal{F}^{\zeta,X}$  be the associated filtration.

Let  $\mathcal{X} \subset \mathbb{R}^n$  and  $\mathcal{U} \subset \mathbb{R}^p$  be a set of states and a set of inputs, respectively. A state and input trajectory  $(X_s, U_s)$  is called an admissible system if

- $X_s \in \mathcal{X}$  for all  $s \ge 0$   $U_s$  is a càglàd,  $\mathcal{F}^{\zeta, X}$ -adapted process such that  $U_s \in \mathcal{U}$  for all s > 0.

Note that the requirement that  $U_s$  is càglàd and  $\mathcal{F}^{\zeta,X}$ adapted implies that  $U(\zeta_t)$  may depend on the "noisy clock" process,  $\zeta_t$ , as well as  $X_r$ , with  $r < \zeta_t$ . If  $\zeta_t \neq t$ , then  $U(\zeta_t)$ cannot directly measure t.

Problem 1: The time-changed optimal control problem over controller time horizon [0, S] is to find a policy  $U_s$  that solves

$$\min_{U} \mathbb{E}\left[\int_{0}^{S} c(s, X_{s}, U_{s}) ds + \Psi(X_{S})\right]$$

where the minimum is taken over all admissible systems  $(X_s, U_s).$ 

Given a policy, U, and  $(s, x) \in [0, S] \times \mathbb{R}^n$ , the cost-to-go function J(s, x; U), is defined by

$$J(s,x;U) = \mathbb{E}\left[\int_{s}^{S} c(s,X_{r},U_{r})dr + \Psi(X_{S}) \middle| X_{s} = x\right].$$

Note, then, that the optimal control problem can be equivalently cast as minimizing J(0, x; U) over all admissible systems.

Backward Evolution Operator: As in standard continuoustime optimal control, the backward evolution operator

$$\mathcal{A}^{u}h(s,x) = \lim_{\sigma \downarrow 0} \frac{1}{\sigma} \left( \mathbb{E} \left[ h(s+\sigma, X_{s+\sigma}) | X_s = x, U_r = u \right] - h(s,x) \right) \quad (12)$$

is used to formulate the dynamic programming equations.

To calculate an explicit form for  $\mathcal{A}^u$ , an auxiliary stochastic process is introduced. For  $(s, x, u) \in [0, S) \times \mathcal{X} \times \mathcal{U}$ , define  $Y_{st}^{xu}$  by

$$Y_{st}^{xu} = x + \int_{0}^{t} F\left(s, Y_{sr}^{xu}, u\right) dr + \int_{0}^{t} G\left(s, Y_{sr}^{xu}, u\right) d\hat{W}_{r} \quad (13)$$

where  $\hat{W}_r$  is a unit Brownian motion independent of  $W_t$  and  $\tau_s$ .

Now the domain of  $\mathcal{A}^u$  is defined. Let  $\mathcal{D}$  be the set of  $h \in \mathcal{C}_p^{1,2}$  such that there exist K and q satisfying

$$\int_{0}^{\infty} \left| \mathbb{E}_{\hat{W}} \left[ h\left(s, Y_{st}^{xu}\right) \right] - h(s, x) \right| \lambda(dt) < K \left( 1 + \|x\|^{q} + \|u\|^{q} \right)$$
(14)

for all  $(s, x, u) \in [0, S) \times \mathcal{X} \times \mathcal{U}$ .

It will be shown in Section IV-B that for  $h \in D$ , the backward evolution operator for  $X_s$  is given by

$$\mathcal{A}^{u}h(s,x) = \frac{\partial h(s,x)}{\partial s} + b\frac{\partial h(s,x)}{\partial x}F(s,x,u) + \frac{1}{2}b\mathrm{Tr}\left(G(s,x,u)^{\mathsf{T}}\frac{\partial^{2}h(s,x)}{\partial x^{2}}G(s,x,u)\right) + \int_{0}^{\infty} \left(\mathbb{E}_{\hat{W}}\left[h\left(s,Y_{st}^{xu}\right)\right] - h(s,x)\right)\lambda(dt).$$
(15)

*Remark 1:* When the dynamics are time-homogeneous, i.e., F(s, y, u) = F(y, u) and G(s, y, u) = G(y, u), and the policy is Markov,  $U_s = U(X_{s^-})$ , the expression for  $\mathcal{A}^u$  in (15) is a special case of Phillips' Theorem [20], [27]. In this case, the formula can be derived using semigroup theory [27]. The derivation in this paper is instead based on Itô calculus.

*Finite Horizon Verification:* The following result is a dynamic programming verification theorem for Problem 1. The theorem is proved in Section IV-B by reducing it to a special case of finite-horizon dynamic programming for controlled Markov processes [26].

Theorem 2: Assume that there is a function  $V \in \mathcal{D}$  that satisfies

$$\inf_{u} \left[ c(s, x, u) + \mathcal{A}^{u} V(s, x) \right] = 0 \tag{16}$$

$$V(S,x) = \Psi(x) \tag{17}$$

where (16) holds for all  $(s, x, u) \in [0, S) \times \mathcal{X} \times \mathcal{U}$  and (17) holds for all  $x \in \mathcal{X}$ .

Then  $V(s, x) \leq J(s, x; U)$  for every admissible system,  $(X_{\cdot}, U_{\cdot})$  with  $X_s = x_{\cdot}$ .

Furthermore, if a policy  $U_r^*$  and associated state process  $X_r^*$ , with  $X_s^* = x$ , satisfy

$$U_r^* \in \operatorname*{arg\,min}_u \left[ c\left(r, X_r^*, u\right) + \mathcal{A}^u V\left(r, X_r^*\right) \right]$$
(18)

for almost all  $(r, \omega) \in [s, S] \times \Omega$ , then  $V(s, x) = J(s, x; U^*)$ .

*Remark 2:* The theorem gives a *sufficient* condition for a policy to be optimal. By the definition of admissible systems if (16), (17) hold and  $(X_r^*, U_r^*)$  satisfy (18), then  $U_r^*$  is optimal over all causal policies, including those with memory.

*Example 5:* Consider the problem of maximizing  $\mathbb{E}[X_S^{\eta}]$ , with  $\eta \in (0, 1)$  subject to the time-changed dynamics

$$dY_t = U(\zeta_t)Y_t (\mu_1 dt + \sigma_1 dW_1(t))$$
$$+ (1 - U(\zeta_t))Y_t (\mu_2 dt + \sigma_2 dW_2(t))$$
$$X_s = Y(\tau_s)$$

where  $W_1(t)$  and  $W_2(t)$  are independent Brownian motions. The problem can be interpreted as allocating wealth between stocks modeled by time-changed geometric Brownian motions:  $Z_i(s) = R_i(\tau_s)$ , where  $dR_i(t) = R_i(t)(\mu_i dt + \sigma_i dW_i(t))$ . Let  $u^*$  maximize the following quadratic, f(u):

Let a maximize the following quadratic, f(a).

$$\frac{1}{2}\eta(\eta-1)\left((u\sigma_1)^2 + ((1-u)\sigma_2)^2\right) + \eta\left(u\mu_1 + (1-u)\mu_2\right)$$

and let  $\rho^* = f(u^*)$ . If  $\rho^* \in \text{dom}(\beta)$ , it can be verified by elementary stochastic calculus that V(s, x) given by

$$V(s,x) = e^{\beta(\rho^*)(S-s)}x^{\eta}$$

satisfies the dynamic programming equations, (16) and (17), with  $\mathcal{X} \times \mathcal{U} = [0, \infty) \times \mathbb{R}$  and max replacing min. The corresponding optimal input is  $U_s^* = u^*$ .

# C. Linear Quadratic Regulators

In this section, Theorem 2 is applied to linear systems with quadratic cost. The result (with no Brownian forcing) was originally presented in [3], using a proof technique specialized for linear systems.

Problem 2: Consider linear dynamics

$$dY_t = (AY_t + BU(\zeta_t)) dt + M dW_t$$
(19)

subject to the time change  $X_s = Y(\tau_s)$ . Here  $\mathcal{X} = \mathbb{R}^n$  and  $\mathcal{U} = \mathbb{R}^p$ .

The time-changed linear quadratic regulator problem over controller time horizon [0, S] is to find a policy  $U_s$  that solves

$$\min_{U} \mathbb{E}\left[\int_{0}^{S} \left(X_{s}^{\mathsf{T}}QX_{s} + U_{s}^{\mathsf{T}}RU_{s}\right)ds + X_{S}^{\mathsf{T}}\Phi X_{S}\right]$$

over all càglàd,  $\mathcal{F}^{\zeta,X}$ -adapted policies. Here Q and  $\Phi$  are positive semidefinite, while R is positive definite.

The following lemma introduces the mappings used to construct the optimal solution for the time-changed linear quadratic regulator problem. The lemma is proved in Appendix C by showing that each mapping may be computed from  $\beta(\tilde{A})$  for an appropriately defined matrix  $\tilde{A}$ .

*Lemma 2:* Let P be an  $n \times n$  matrix. If  $\{0\} \cup \operatorname{spec}(2A) \subset \operatorname{dom}(\beta)$ , then the following linear mappings are well defined:

$$\begin{split} F(P) &= b(A^{\mathsf{T}}P + PA) + \int_{0}^{\infty} \left( e^{A^{\mathsf{T}}t} P e^{At} - P \right) \lambda(dt) \\ G(P) &= bP + \int_{0}^{\infty} e^{A^{\mathsf{T}}t} P \int_{0}^{t} e^{Ar} dr \lambda(dt) \\ H(P) &= \int_{0}^{\infty} \int_{0}^{t} e^{A^{\mathsf{T}}r} dr P \int_{0}^{t} e^{A\rho} d\rho \lambda(dt) \\ g(P) &= \operatorname{Tr} \left( P \left( bMM^{\mathsf{T}} + \int_{0}^{\infty} \int_{0}^{t} e^{Ar} MM^{\mathsf{T}} e^{A^{\mathsf{T}}r} dr \lambda(dt) \right) \right). \end{split}$$

Furthermore, F, G, and H satisfy

$$\mathbb{E}\left[e^{A^{\mathsf{T}}\tau_s}Pe^{A\tau_s}\right] = P + sF(P) + O(s^2)$$
$$\mathbb{E}\left[e^{A^{\mathsf{T}}\tau_s}P\int_0^{\tau_s}e^{Ar}dr\right] = sG(P) + O(s^2)$$
$$\mathbb{E}\left[\int_0^{\tau_s}e^{A^{\mathsf{T}}r}drP\int_0^{\tau_s}e^{A\rho}d\rho\right] = sH(P) + O(s^2).$$

*Remark 3:* The descriptions of F, G, and H in terms of expectations are not required for the proof below. They are given to demonstrate that the formulas in terms of  $(b, \lambda)$  coincide with the formulas from [3].

Example 6: With no temporal noise, the mappings become

$$F(P) = A^{\mathsf{T}}P + YP, \quad G(P) = P,$$
  
$$H(P) = 0, \quad g(P) = \operatorname{Tr}(PMM^{\mathsf{T}}). \tag{20}$$

Furthermore, since  $\beta(z) = z$  is analytic everywhere, these formulas are true for any state matrix, A.

*Example 7:* Consider an arbitrary strictly increasing subordinator with Laplace exponent  $\psi$ . Let  $A = \mu$  where  $\mu$  is a real, non-zero scalar with  $2\mu \in \text{dom}(\beta)$ . Let M be a scalar. Combining (2) with the formula  $\int_0^t e^{\mu\sigma} d\sigma = \mu^{-1}(e^{\mu t} - 1)$ shows that

$$\begin{split} F(P) &= \beta(2\mu)P \\ G(P) &= \mu^{-1} \left( \beta(2\mu) - \beta(\mu) \right) P \\ H(P) &= \mu^{-2} \left( \beta(2\mu) - 2\beta(\mu) \right) P \\ g(P) &= \frac{1}{2} \mu^{-1} \beta(2\mu) M^2 P. \end{split}$$

Theorem 3: Say that  $\{0\} \cup \operatorname{spec}(2A) \subset \operatorname{dom}(\beta)$ . Define the function  $V(s, x) = x^{\mathsf{T}} P_s x + h_s$  by the backward differential equations

$$-\frac{d}{ds}P_s = Q + F(P_s)$$
$$-G(P_s)B\left(R + B^{\mathsf{T}}H(P_s)B\right)^{-1}B^{\mathsf{T}}G(P_s)^{\mathsf{T}}$$
$$-\frac{d}{ds}h_s = g(P_s)$$

with final conditions  $P_S = \Phi$  and  $h_S = 0$ . The function V(s, x) satisfies dynamic programming equations, (16) and (17), and the optimal policy is given by

$$U_s = K_s X_{s^-}$$
$$K_s = -\left(R + B^{\mathsf{T}} H(P_s) B\right)^{-1} B^{\mathsf{T}} G(P_s)^{\mathsf{T}}.$$

The policy is computed by minimizing (18). Thus, Theorem 2 implies that linear state feedback is optimal over all causal policies.

A straightforward variation on the proof of Theorem 3 shows that for any linear policy,  $U_s = L_s X_{s^-}$ , the cost-to-go is given by

$$J(s,x;U) = x^{\mathsf{T}}Z_s x + p_s$$



Fig. 2. (a) Plots of  $X_1(s)$  under the optimal policy and the LQR policy for 10 realizations of  $\tau_s$ . The initial condition is  $x = [0, 1]^{\mathsf{T}}$ . (b) The same plots under time variable t. The black line shows the LQR trajectory with no temporal noise. In the case of no temporal noise, the classical LQR uses high gains near t = 0 to produces high-speed trajectories such that  $Y_1$  approaches 0 at final time. In this case, timing errors lead to wide variation in the final position. The optimal policy reduces the speed of the trajectory near s = 0 order to minimize the effects of temporal noise. (c) The optimal cost V(s,x) and J(s, x; U) for the LQR policy (i.e. U is optimal without time noise) are plotted for  $x = [0, 1]^{\top}$ . As expected,  $V(s, x) \leq J(s, x; U)$ . Furthermore, as the timehorizon increases, the LQR policy depends strongly on timing information, and so temporal noise leads to higher cost as s goes to 0. (d) A histogram of the final positions,  $X_1(S)$ , for 1000 realizations of  $\tau_s$ . The optimal controller leads to  $X_1(S)$  being tightly distributed around 0, while the LQR controller gives a wide spread of  $X_1(S)$  values. The errors in the final position lead to increased cost for the LQR controller.

where  $Z_s$  and  $k_s$  satisfy the backward differential equations

$$-\frac{d}{ds}Z_s = Q + F(Z_s) + L_s^{\mathsf{T}}B^{\mathsf{T}}G(Z_s)^{\mathsf{T}} + G(Z_s)BL_s$$
$$+ L_s^{\mathsf{T}}\left(R + B^{\mathsf{T}}H(Z_s)B\right)L_s$$
$$-\frac{d}{ds}p_s = g(Z_s).$$

In the following example, these formulas are used in order to compare the performance of the policy from Theorem 3 with the policy  $U_s = L_s X_{s^-}$ , where  $L_s$  is the standard LQR gain, not compensating for temporal noise.

*Example* 8: Consider the system defined by the state matrices

$$A = \begin{bmatrix} 0.75 & 1\\ 0 & 0.75 \end{bmatrix}, \qquad B = \begin{bmatrix} 0\\ 1 \end{bmatrix}, \qquad M = 0$$

with cost matrices given by

$$R = 0.5, \qquad Q = 0, \qquad \Phi = \begin{bmatrix} 1 & 0 \\ 0 & 0 \end{bmatrix}.$$

Let  $\tau_s$  be the inverse Gaussian subordinator with  $\gamma = \delta = 2$ . The condition,  $\operatorname{spec}(2A) \subset \operatorname{dom}(\beta)$ , is satisfied since  $2 \cdot 0.75 = 1.5 < \gamma^2/2 = 2$ . Fig. 2 compares the optimal policy with the standard LQR policy.

### IV. PROOFS OF MAIN RESULTS

A. Proof of Theorem 1

From the definition of  $X_s$ 

$$X_{s} = X_{0} + \int_{0}^{\tau_{s}} F_{t} dt + \int_{0}^{\tau_{s}} G_{t} dW_{t}.$$
 (21)

Note that  $X_s$  is  $\mathcal{F}_s^{\tau,W}$ -measurable. Thus,  $X_s$  will be an  $\mathcal{F}^{\tau,W}$  semimartingale, provided that 1)  $\int_0^{\tau_s} F_t dt$  has finite variation and 2)  $\int_0^{\tau_s} G_t dW_t$  is an  $\mathcal{F}^{\tau,W}$  martingale.

Finite variation follows since:

$$\operatorname{Var}\left(\int_{0}^{\tau_{s}} F_{t} dt\right) \leq \int_{0}^{\tau_{s}} \|F_{t}\| dt < \infty \text{ almost surely.}$$

To prove the martingale property, note that for  $0 \leq r \leq s$ , we have

$$\mathbb{E}\left[\int_{0}^{\tau_{s}} G_{t} dW_{t} \middle| \mathcal{F}_{r}^{\tau,W}\right] = \int_{0}^{\tau_{r}} G_{t} dW_{t} + \mathbb{E}\left[\int_{\tau_{r}}^{\tau_{s}} G_{t} dW_{t} \middle| \mathcal{F}_{r}^{\tau,W}\right]$$
$$= \int_{0}^{\tau_{r}} G_{t} dW_{t}.$$

Furthermore

$$\mathbb{E}\left[\left\|\int_{0}^{\tau_{s}}G_{t}dW_{t}\right\|\right]^{2} \leq \mathbb{E}\left[\left\|\int_{0}^{\tau_{s}}G_{t}dW_{t}\right\|^{2}\right] < \infty$$

by Jensen's inequality. Thus 2) holds.

Now (10) must be proved. For more compact notation, define the processes  $H_t$  and  $Z_t$  as

$$H_t = \begin{bmatrix} F_t & G_t \end{bmatrix} \qquad Z_t = \begin{bmatrix} t \\ W_t \end{bmatrix}$$

so that  $X_s$  may be written as

$$X_s = \int_0^{\tau_s} H_t dZ_t.$$
(22)

Note that  $Z(\tau_s) = [\tau_s, W(\tau_s)^{\mathsf{T}}]^{\mathsf{T}}$ . Since  $\tau_s$  is a subordinator,  $W(\tau_s)$  is a Lévy process on  $\mathbb{R}^d$ , with Lévy symbol

$$\eta_{W_{\tau}}(z) = -\frac{1}{2}bz^{\mathsf{T}}z + \int\limits_{\mathbb{R}^d} \left(e^{iz^{\mathsf{T}}x} - 1\right)\mu_{W,\tau}(dt)$$

for some Lévy measure  $\mu_{W,\tau}$ . (See [20, Theorem 1.3.25 and Theorem 1.3.33].) Thus, the continuous part of  $W(\tau_s)$  is a Brownian motion with  $\mathbb{E}[W(\tau_s)W(\tau_s)^{\mathsf{T}}] = bI$ .

Define the  $\tilde{Z}_s$  by removing the jumps from  $Z(\tau_s)$ 

$$\tilde{Z}_s = Z(\tau_s) - \sum_{0 \le r \le s} \left( Z(\tau_r) - Z(\tau_r^-) \right).$$

It follows that  $\tilde{Z}_s = [bs, \sqrt{b}\tilde{W}_s^{\mathsf{T}}]^{\mathsf{T}}$ , where  $\tilde{W}_s$  is the Brownian motion from the theorem statement. Thus, (10) can be equivalently written as

$$X_{s} = \int_{0}^{s} H(\tau_{r^{-}}) d\tilde{Z}_{s} + \sum_{0 \le r \le s} \int_{\tau_{r^{-}}}^{\tau_{r}} H_{s} dZ_{t}$$
(23)

So, the proof will follow by deriving (23) from (22). If b = 0, then  $\tilde{Z}_s = 0$  and  $\tau_s = \sum_{0 \le r \le s} \tau_r - \tau_{r^-}$ . Thus

$$X_s = \sum_{0 \le r \le s} \int_{\tau_{r^-}}^{\tau_r} H_s dZ_t$$

so, in this case, (10) holds.

Now assume b > 0. The cases when  $\tau_s$  has finite rate  $(\lambda((0,\infty)) < \infty)$  and infinite rate  $(\lambda((0,\infty)) = \infty)$  will be treated separately.

*Finite Rate:* Let  $r_0 = 0$  and let  $r_1, r_2, \ldots$  be the jump times of  $\tau_s$ . With probability 1, there exist a finite (random) integer L such that L jumps occur over [0, s]. Note that (22) may be expanded as

$$X_{s} = \int_{\tau(r_{L})}^{\tau_{s}} H_{t} dZ_{t} + \sum_{k=0}^{L-1} \left[ \int_{\tau(r_{k})}^{\tau(r_{k+1})} H_{t} dZ_{t} + \int_{\tau(r_{k+1})}^{\tau(r_{k+1})} H_{t} dZ_{t} \right]$$
(24)

Let  $s_0^n \leq s_1^n \leq \cdots \leq s_{K_n}^n$  be a sequence of partitions such that

$$\begin{split} \lim_{n \to \infty} s_{K_n}^n &= \infty & a.s. \\ \lim_{n \to \infty} \sup \left\{ \left| s_{k+1}^n - s_k^n \right| : k = 0, \dots, K_n - 1 \right\} = 0 & a.s. \\ \left\{ r_i : r_i \leq s_{K_n}^n \right\} \subset \left\{ s_0^n, \dots, s_{K_n}^n \right\}. \end{split}$$

The last condition ensures that the jump times are contained in the partition.

Note that for  $s \in [r_k, r_{k+1})$ ,  $\tau_s = bs + \tau^d(r_k)$ , where  $\tau_s^d$  is the discontinuous part of  $\tau_s$ . Since b > 0 follows that the sequence  $\tau(s_0^n), \tau(s_1^n), \ldots$ , satisfies the following properties, almost surely:

$$\lim_{n \to \infty} \tau \left( s_{K_n}^n \right) = \infty$$
$$\lim_{n \to \infty} \sup \left\{ \left| \tau \left( s_{i+1}^n \right) - \tau \left( s_i^n \right) \right| : i \ge 0 \right\} = 0$$

Using a standard argument from stochastic integration (see [28, Theorem II.21]), the integral from  $\tau(r_k)$  to  $\tau(r_{k+1})$  may be evaluated as

$$\begin{aligned} & \tau(r_{k+1}) \\ & \int_{\tau(r_{k})} H_{s} dZ_{t} \\ &= \lim_{n \to \infty} \sum_{r_{k} \leq s_{i}^{n} < r_{k+1}} H\left(\tau(s_{i}^{n})\right) \left(Z\left(\tau(s_{i+1})\right) - Z\left(\tau(s_{i})\right)\right) \\ &= \lim_{n \to \infty} \sum_{r_{k} \leq s_{i}^{n} < r_{k+1}} H\left(\tau\left(s_{i}^{n}\right)\right) \left(\tilde{Z}(s_{i+1}) - \tilde{Z}(s_{i})\right) \\ &= \int_{r_{k}} H\left(\tau_{s^{-}}\right) d\tilde{Z}_{s}. \end{aligned}$$

$$(25)$$

The second equality uses the fact that no jumps occur over  $(r_k, r_{k+1})$ . Equation (23) follows by combining (24) and (25).

Infinite Rate: Let  $\epsilon_n > 0$  be a sequence decreasing to 0, at a rate to be specified later. Define  $\tau_s^n$  to be the process by removing all jumps of size at most  $\epsilon_n$  from  $\tau_s$ 

$$\tau_s^n = bs + \int_{\epsilon_n}^{\infty} tN(s, dt).$$
(26)

Let  $r_0^n = 0$ , and let  $r_1^n, r_2^n, \ldots$  be the jump times of  $\tau_s^n$ . Let  $L_s^n = \sup\{k : r_k^n \le s\}$ . With probability 1,  $L_s^n < \infty$ . If  $\epsilon_n$  are chosen as in Lemma 3 from Appendix A, then  $X_s$  may be computed as a limit

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$$X_{s} = \lim_{n \to \infty} \left[ \int_{\tau\left(r_{L_{s}^{n}}^{n}\right)}^{\tau_{s}} H_{t} dZ_{t} + \sum_{k=0}^{L_{s}^{n}-1} H\left(\tau\left(r_{k}^{n}\right)\right) \left(Z\left(\tau\left(r_{k+1}^{n-1}\right)\right) - Z\left(\tau\left(r_{k}^{n}\right)\right)\right) + \sum_{k=0}^{L_{s}^{n}-1} \int_{\tau\left(r_{k+1}^{n-1}\right)}^{\tau\left(r_{k+1}^{n}\right)} H_{t} dZ_{t} \right].$$
(27)

Note that  $Z(\tau(r_{k+1}^{n-})) - Z(\tau(r_{k}^{n}))$  may be expressed as  $Z\left(\tau\left(r_{k+1}^{n-}\right)\right) - Z\left(\tau\left(r_{k}^{n}\right)\right)$  $= \tilde{Z}\left(r_{k+1}^{n}\right) - \tilde{Z}\left(r_{k}^{n}\right) + \sum_{\substack{r_{k}^{n} < r \leq r_{k+1}^{n} \\ \Delta \tau_{r} \leq \epsilon_{n}}} \left(Z(\tau_{r}) - Z(\tau_{r^{-}})\right).$ 

Note that the terms in the summation all vanish as  $\epsilon_n \rightarrow 0$ . Furthermore,  $r_{L_s^n}^n \uparrow s$ , almost surely. Thus, (27) can be expressed as

$$X_{s} = \lim_{n \to \infty} \left[ \sum_{k=0}^{L_{s}^{n}-1} H\left(\tau\left(r_{k}^{n}\right)\right) \left(\tilde{Z}(r_{k+1}) - \tilde{Z}(r_{k})\right) + \sum_{k=0}^{L_{s}^{n}-1} \int_{\tau\left(r_{k+1}^{n}\right)}^{\tau\left(r_{k+1}^{n}\right)} H_{t} dZ_{t} \right]$$

and (23) now follows using [28, Theorem II.21].

## B. Proof of Theorem 2

Theorem 2 is a special case of finite-horizon dynamic programming for controlled Markov processes ([26, Theorem III.8.1]), provided that the following two conditions hold for all  $h \in \mathcal{D}$ :

- (i) The backward evolution operator, defined in (12) is given by the formula in (15).
- (ii) If *h* satisfies

$$\mathbb{E}\left[\left|h(S, X_S)\right| \left|X_s = x\right] < \infty, \text{ and } \\ \mathbb{E}\left[\int\limits_{s}^{S} \left|\mathcal{A}^{U_r} h(r, X_r)\right| dr \left|X_s = x\right] < \infty\right]$$

then the Dynkin formula holds

$$\mathbb{E}\left[h(S, X_S)|X_s = x\right] - h(s, x)$$
$$= \mathbb{E}\left[\int_{s}^{S} \mathcal{A}^{U_r} h(r, X_r) dr |X_s = x\right]. \quad (28)$$

First, using Theorem 1, a more explicit formula for  $X_s$  is derived, and then using Itô's formula for semimartingales, a formula for  $h(s, X_s)$  is given. Using the formula for  $h(s, X_s)$ , (15) and (28) are then proved.

Note that  $Y(\tau_{r^-}) = X_{r^-}$  and for all  $t \in [\tau_{r^-}, \tau_r]$ ,  $\zeta_t = r$ . Therefore

$$\begin{split} F\left(\zeta(\tau_{r^{-}}),Y(\tau_{r^{-}}),U\left(\zeta(\tau_{r^{-}})\right)\right) &= F(r,X_{r^{-}},U_{r}), \quad \text{and} \\ F\left(\zeta_{t},Y_{t},U(\zeta_{t})\right) &= F(r,Y_{t},U_{r}) \quad \text{for all} \quad t \in [\tau_{r^{-}},\tau_{r}]. \end{split}$$

The expressions for G are similar. Thus, Theorem 1 implies that  $X_s$  is given by

$$X_{s} = X_{0} + b \int_{0}^{s} F(r, X_{r^{-}}, U_{r}) dr + \sqrt{b} \int_{0}^{s} G(r, X_{r^{-}}, U_{r}) d\tilde{W}_{r}$$
$$+ \sum_{0 \le r \le s} \left( \int_{\tau_{r^{-}}}^{\tau_{r}} F(r, Y_{t}, U_{r}) dt + \int_{\tau_{r^{-}}}^{\tau_{r}} G(r, Y_{t}, U_{r}) dW_{t} \right).$$
(29)

Now a formula for  $h(s, X_s)$  will be derived. Note that for any càglàd,  $\mathcal{F}^{\tau,W}$ -adapted process,  $Z_s$ , the stochastic integral with respect to  $X_s$  is given by

$$\int_{0}^{s} Z_{r} dX_{r} = \int_{0}^{s} Z_{r} bF(r, X_{r^{-}}, U_{r}) dr + \int_{0}^{s} Z_{r} \sqrt{b} G(r, X_{r^{-}}, U_{r}) d\tilde{W}_{r} + \sum_{0 \le r \le s} Z_{r} (X_{r} - X_{r^{-}}).$$

Furthermore, the continuous part of the quadratic variation is given by

$$[X,X]_s^c = \int_0^{\infty} \frac{1}{2} b G(r,X_{r^-},U_r) G(r,X_{r^-},U_r)^{\mathsf{T}} dr.$$

Thus Itô's formula for semimartingales (see [28]) implies that  $h(s, X_s)$  is given by

$$h(s, X_s) = h(0, X_0)$$

$$+ \int_0^s \left( \frac{\partial h(r, X_{r^-})}{\partial r} + b \frac{\partial h(r, X_{r^-})}{\partial x} F(r, X_{r^-}, U_r) \right) dr$$

$$+ \int_0^s \frac{1}{2} b \operatorname{Tr} \left( G(r, X_{r^-}, U_r)^{\mathsf{T}} \frac{\partial^2 h(r, X_{r^-})}{\partial x^2} G(r, X_{r^-}, U_r) \right) dr$$

$$+ \int_0^s \sqrt{b} \frac{\partial h(r, X_{r^-})}{\partial x} G(r, X_{r^-}, U_r) d\tilde{W}_r$$

$$+ \sum_{0 \le r \le s} \left( h(r, X_r) - h(r, X_{r^-}) \right). \tag{30}$$

Now (15) will be derived from (30). Assume that  $X_s = x$ and  $U_r = u$  for  $r \in [s, s + h]$ . Then (30) implies that

$$\mathbb{E}\left[h(s+h, X_{s+h})\right] = h(s, x)$$

$$+ \mathbb{E}\left[\int_{s}^{s+h} \left(\frac{\partial h(r, X_{r^{-}})}{\partial r} + b\frac{\partial h(r, X_{r^{-}})}{\partial x}F(r, X_{r^{-}}, u)\right)dr\right]$$

$$+ \mathbb{E}\left[\frac{1}{2}b\mathrm{Tr}\left(G(r, X_{r^{-}}, u)^{\mathsf{T}}\frac{\partial^{2}h(r, X_{r^{-}})}{\partial x^{2}}G(r, X_{r^{-}}, u)\right)dr\right]$$

$$+ \mathbb{E}\left[\sum_{s < r \le s+h} \left(h(r, X_{r}) - h(r, X_{r^{-}})\right)\right].$$
(31)

If r > s, the Brownian motions  $W_t$  for  $t \in [\tau_{r^-}, \tau_r]$  and  $\hat{W}_t$  for  $t \in [0, \Delta \tau_r]$  are identically distributed and independent of  $\mathcal{F}_s^{\tau,W}$ . Therefore, using (13), and given that  $X_s = x$  and  $U_r = u$ , the expectations of the jump terms may be written as

$$\mathbb{E}\left[h(r, X_r) - h(r, X_{r^-})\right] = \mathbb{E}\left[\Theta(r, \Delta \tau_r, X_{r^-}, u)\right]$$

where  $\Theta(r, \Delta \tau_r, X_{r^-}, u) = \mathbb{E}_{\hat{W}}[h(r, Y_{r, \Delta \tau_r}^{X_r - u})] - h(r, X_{r^-}).$ Thus, the term at the bottom of (31) may be evaluated as a Poisson integral

$$\mathbb{E}\left[\sum_{s < r \le s+h} \left(h(r, X_r) - h(r, X_{r^-})\right)\right]$$
$$= \mathbb{E}\left[\int_{s}^{s+h} \int_{0}^{\infty} \Theta(r, t, X_{r^-}, u) N(dr, dt)\right]$$
$$= \mathbb{E}_W\left[\int_{s}^{s+h} \int_{0}^{\infty} \Theta(r, t, X_{r^-}, u) \lambda(dt) dr\right] \qquad (32)$$

where the second is equation follows from Fubini's theorem and (14).

Now, (15) results from combining (12) with (31) and (32). Turning to (28), since  $X_s$  is  $\mathcal{F}_s^{\tau,W}$  measurable, it suffices to prove that

$$\mathbb{E}\left[h(S, X_S) - h(s, X_s) - \int\limits_{s}^{S} \mathcal{A}^{U_s} h(r, X_r) dr | \mathcal{F}_s^{\tau, W}\right] = 0.$$

Since  $X_r(\omega) = X_{r^-}(\omega)$  for almost all  $(r, \omega) \in [s, S] \times \Omega$ , and omitting  $\mathcal{F}_s^{\tau, W}$  from the expectations for brevity, it follows that:

$$\mathbb{E}\left[\int\limits_{s}^{S} \mathcal{A}^{U_{s}} h(r, X_{r}) dr\right] = \mathbb{E}\left[\int\limits_{s}^{S} \mathcal{A}^{U_{s}} h(r, X_{r^{-}}) dr\right].$$

Combining (15) and (30) implies that

$$\mathbb{E}\left[h(S, X_S) - h(s, X_s) - \int_s^S \mathcal{A}^{U_r} h(r, X_{r^-}) dr\right]$$
$$= \mathbb{E}\left[\sum_{s < r \le S} \left(h(r, X_r) - h(r, X_{r^-})\right)\right]$$
$$- \mathbb{E}\left[\int_s^S \int_0^\infty \left(\mathbb{E}_{\hat{W}}\left[h\left(r, Y_t^{X_{r^-}U_r}\right)\right] - h(r, X_{r^-})\right)\lambda(dt) dr\right]$$
(33)

As in the proof of (15), the two terms at the bottom of (33)are equal in expectation. Thus (28) holds and the proof is complete.

# C. Proof of Theorem 3

Assume that  $V(s, x) = x^{\mathsf{T}} P_s x + h_s$ . Applying the backward evolution operator corresponding to (19) to V(s, x) results in

$$\mathcal{A}^{u}V(s,x) = x^{\mathsf{T}}\dot{P}_{s}x + \dot{h}_{s} + bx^{\mathsf{T}}(A^{\mathsf{T}}P_{s} + P_{s}A)x + 2bx^{\mathsf{T}}P_{s}Bu + b\mathrm{Tr}(P_{s}MM^{\mathsf{T}}) + \int_{0}^{\infty} \Bigl(\mathbb{E}_{\hat{W}}\Bigl[Y_{t}^{xu}{}^{\mathsf{T}}P_{s}Y_{t}^{xu}\Bigr] - x^{\mathsf{T}}P_{s}x\Bigr)\lambda(dt).$$
(34)

Note that  $\mathbb{E}_{\hat{W}}[Y_t^{xu}Y_t^{xu^{\mathsf{T}}}]$  may be evaluated as

$$\mathbb{E}_{\hat{W}}\left[Y_t^{xu} Y_t^{xu\mathsf{T}}\right] = y_t^{xu} y_t^{xu\mathsf{T}} + \Sigma_t$$

where  $y_t^{xu}$  is the mean of  $Y_t^{xu}$  and  $\Sigma_t$  is the covariance. A standard argument in linear stochastic differential equations shows that the mean and covariance are given by

$$y_t^{xu} = e^{At}x + \int_0^t e^{Ar} dr Bu, \quad \Sigma_t = \int_0^t e^{Ar} M M^{\mathsf{T}} e^{A^{\mathsf{T}} r} dr.$$

Thus, the integral in (34) may be written as

$$\int_{0}^{\infty} \left( \mathbb{E}_{\hat{W}} \left[ Y_{t}^{xu^{\mathsf{T}}} P_{s} Y_{t}^{xu} \right] - x^{\mathsf{T}} P_{s} x \right) \lambda(dt)$$

$$= \int_{0}^{\infty} \left[ x \\ u \right]^{\mathsf{T}} \left[ \begin{array}{c} \hat{F}(t, P_{s}) & \hat{G}(t, P_{s})B \\ B^{\mathsf{T}} \hat{G}(t, P_{s})^{\mathsf{T}} & B^{\mathsf{T}} \hat{H}(t, P_{s})B \end{array} \right] \left[ x \\ u \end{array} \right] \lambda(dt)$$

$$+ \int_{0}^{\infty} \int_{0}^{t} \operatorname{Tr} \left( P_{s} e^{Ar} M M^{\mathsf{T}} e^{A^{\mathsf{T}} r} \right) dr \lambda(dt)$$
(35)

$$\begin{split} \hat{F}(t,P) &= e^{A^{\mathsf{T}}t} P e^{At} - P\\ \hat{G}(t,P) &= e^{A^{\mathsf{T}}t} P \int_{0}^{t} e^{A\rho} d\rho\\ \hat{H}(t,P) &= \int_{0}^{t} e^{A^{\mathsf{T}}r} dr P \int_{0}^{t} e^{A\rho} d\rho. \end{split}$$

Combining (34) and (35), and using the linear operators from Lemma 2 gives

$$\mathcal{A}^{u}V(s,x) = \begin{bmatrix} x \\ u \end{bmatrix}^{\mathsf{T}} \begin{bmatrix} \dot{P}_{s} + F(P_{s}) & G(P_{s})B \\ B^{\mathsf{T}}G(P_{s}) & B^{\mathsf{T}}H(P_{s})B \end{bmatrix} \begin{bmatrix} x \\ u \end{bmatrix} + \dot{h}_{s} + g(P_{s}).$$

Therefore, adding the cost gives

$$\begin{aligned} x^{\mathsf{T}}Qx + u^{\mathsf{T}}Ru + \mathcal{A}^{u}V(s,x) &= \dot{h}_{s} + g(P_{s}) \\ &+ \begin{bmatrix} x \\ u \end{bmatrix}^{\mathsf{T}} \begin{bmatrix} Q + \dot{P}_{s} + F(P_{s}) & G(P_{s})B \\ B^{\mathsf{T}}G(P_{s}) & R + B^{\mathsf{T}}H(P_{s})B \end{bmatrix} \begin{bmatrix} x \\ u \end{bmatrix} \end{aligned}$$

The result now follows from quadratic minimization.

#### V. DISCUSSION

The work in this paper lays a theoretical foundation for future research on biological motor control [29], [30], finance [31], and multi-agent control [32], [33] in context that the controller is uncertain about the time of the plant. Possible extensions include time estimation from sensory data, optimal control control with different time horizons, and control with multiple noisy clocks. It will also be useful to consider alternative clock models, and practical methods for controller computation for nonlinear systems.

*Time Estimation:* In our setting, the system is a controlled Markov process with respect to controller time, and the optimal solution only depends on the current values of the state and the controller time. If the state cannot be perfectly measured, the horizon is taken with respect to plant time, or the clock  $\tau_s$  does not have independent increments, then control performance might be improved by using the measurement history to estimate time. For example, in option pricing, inferences about the "business time" can be used to estimate the volatility of stock prices [34]. To perceive time, humans appear to integrate sensory cues about the passage of time in a Bayesian manner [35]; humans also appear to incorporate sensory information about the timing of events to improve state estimation [36]. Formal analysis of optimal control with time estimation would be interesting, but appears to be challenging.

It would also be interesting to test if probing control strategies could improve time estimation. For example the distance that a mass moves in response to a probing force gives information about how much time has passed. It may be possible to use the improved time estimate to achieve better control performance at the cost of extra actuation.

*Plant Time Control:* This paper studies a controller horizon [0, S], which is an interval of time with respect to the clock measured by the controller. For portfolio optimization, in which the

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controller measures calendar time, such a horizon is sensible. In other problems, such as control with real-time constraints, it may be more natural to consider a horizon [0, T] over plant time. In this case, the challenges of time estimation discussed above must be addressed.

Multiple Clocks: In this paper, we assumed that the plant dynamics evolve according to one clock, while the controller can measure a different clock. If the plant consists of numerous subsystems, then each could potentially evolve according to a different clock. This scenario arises in portfolio problems, in which the goal is to allocate wealth between a bank process which accrues interest at known, fixed rate and a stock process that evolves in a variable rate market [22]. Here, the bank process may be interpreted as evolving with respect to a perfect clock, while the stock process may be viewed as evolving with respect to a noisy clock. In engineering applications, such as mobile sensor networks, multiple autonomous agents with their own clocks solve cooperative control problems. Currently, problems arising from drift between clocks are mitigated by using expensive clocks and time synchronization protocols. The work in this paper will be extended to reduce the need for precision timing and synchronization.

*Alternative Clock Models:* This paper modeled clocks with Lévy processes largely for mathematical convenience. It would be interesting to study similar control problems based on jitter [12], [13], phase noise [37], stochastic integrals [4], or biological timing models [38].

*Efficient Computation:* Currently, the dynamic programming method in this paper is limited to special problems that can be solved exactly. It may be possible to solve a broader class of problems by integrating approximate dynamic programming techniques [39]–[41].

#### VI. CONCLUSION

This paper gives basic results on control with uncertainty in time. The technical backbone of the paper is Theorem 1 which expresses the original plant dynamics in terms of the controller's clock index. Using the new representation, the system becomes a controlled Markov process, and thus existing dynamic programming theory can be applied. Given the dynamic programming equations, time changed versions of linear quadratic control and a nonlinear portfolio problem problem are solved explicitly.

# APPENDIX A

# A TECHNICAL LEMMA FOR THEOREM 1

Lemma 3: Let  $\tau_s$  be an infinite rate subordinator. Let  $r_0^n = 0$ and let  $r_1^n \leq r_2^n \leq \cdots$  be the jump times of  $\tau_s^n$ , from (26). For any sequence  $S_n \to \infty$  there is a sequence  $\epsilon_n \downarrow 0$  such that the following limits hold, almost surely:

$$\lim_{n \to \infty} \sup \left\{ r_i^n : r_i^n \le S_n \right\} = \infty$$
(36)

$$\lim_{n \to \infty} \sup\left\{\tau_{r_i^n} : r_i^n \le S_n\right\} = \infty$$
(37)

$$\lim_{n \to \infty} \sup\left\{r_{i+1}^n - r_i^n : r_i^n \le S_n\right\} = 0$$
(38)

$$\lim_{n \to \infty} \sup \left\{ \tau_{r_{i+1}^n} - \tau_{r_i^n} : r_i^n \le S_n \right\} = 0.$$
(39)

*Proof:* First it will be shown that for any sequence  $S_n \rightarrow \infty$ , a sequence  $\epsilon_n \downarrow 0$  can be chosen such that (38) and (39) hold. Then it will be shown that (36) and (37) hold.

Consider (38). Using the Borel-Cantelli lemma, it is sufficient to prove that for some constant R > 0, and  $\epsilon_n$  sufficiently small

$$\mathbb{P}\left(\sup_{r_i^n \le S_n} |r_{i+1}^n - r_i^n| \ge \frac{1}{2^n}\right) < \frac{R}{2^n}.$$
 (40)

For ease of notation, the superscripts on  $r_i^n$  and the subscripts on  $\epsilon_n$  and  $S_n$  will be dropped.

With probability 1,  $\tau_s^n$  has only a finite number of jumps over [0, S], so let  $K = \max\{i : r_i \leq S\}$ .

Consider (40). Define the function  $g(\epsilon)$  by

$$g(\epsilon) = \int_{\epsilon}^{\infty} \lambda(dt).$$

Note that the differences are  $r_{i+1} - r_i$  are exponential random variables with rate parameter  $g(\epsilon)$ . Thus, the event that  $r_{i+1} - r_i \ge 1/2^n$  is a Bernoulli random variable with probability  $p(\epsilon)$  given by

$$p(\epsilon) = \mathbb{P}\left(r_{i+1} - r_i \ge \frac{1}{2^n}\right) = g(\epsilon) \int_{\frac{1}{2^n}}^{\infty} e^{-g(\epsilon)x} dx$$
$$= e^{-g(\epsilon)/2^n}$$

Let J be the geometric random variable defined by

$$J = \min\left\{i : r_{i+1} - r_i \ge \frac{1}{2^n}\right\}$$

Then the probability of J is given by

$$\mathbb{P}(J=k) = (1-p(\epsilon))^k p(\epsilon).$$

Using the definitions of K and J, the probability in (40) may be written as

$$\mathbb{P}\left(\sup_{r_i \le S} |r_{i+1} - r_i| \ge \frac{1}{2^n}\right) = \mathbb{P}(J \le K).$$

Furthermore, given any constant M > 0

$$\mathbb{P}(J \le K) \le \mathbb{P}(J \le M) + \mathbb{P}(M \le K).$$
(41)

Thus, (40) may be bounded by bounding the terms on the right of (41) separately.

Now  $\mathbb{P}(M \leq K)$  will be bounded. Note that K is a Poisson random variable with parameter  $Sg(\epsilon)$ . Markov's inequality thus shows that

$$\mathbb{P}(M \le K) \le \frac{1}{M} \mathbb{E}[K] = \frac{Sg(\epsilon)}{M}$$
(42)

The term  $\mathbb{P}(J \leq M)$  can be computed exactly as

$$\mathbb{P}(J \le M) = p(\epsilon) \sum_{k=0}^{M} (1 - p(\epsilon))^{k} = 1 - (1 - p(\epsilon))^{M+1}.$$

Thus, (40) will hold if M can be chosen such that

$$Sg(\epsilon)/M < 1/2^n$$
 and  $1 - (1 - p(\epsilon))^{M+1} < 1/2^n$ . (43)

Rearranging terms, (43) is equivalent to

$$2^n Sg(\epsilon) < M < \frac{\log\left(1 - \frac{1}{2^n}\right)}{\log\left(1 - p(\epsilon)\right)} - 1.$$

Therefore, a suitable constant M exists if

$$(2^n Sg(\epsilon) + 1)\log\left(1 - p(\epsilon)\right) > \log\left(1 - \frac{1}{2^n}\right).$$
(44)

It will be shown that (44) holds if  $\epsilon$  is sufficiently small. Since  $\tau_s$  has infinite rate,  $g(\epsilon) \to \infty$  and  $p(\epsilon) \to 0$  as  $\epsilon \to 0$ . Setting  $\gamma = g(\epsilon)/2^n$ , the left side of (44) converges to 0 by L'Hôspital's rule

$$\lim_{h \to \infty} \gamma \log(1 - e^{-\gamma}) = \lim_{\gamma \to \infty} \frac{\frac{e^{-\gamma}}{1 - e^{-\gamma}}}{-\gamma^{-2}} = -\lim_{\gamma \to \infty} \frac{\gamma^2}{e^{\gamma} - 1} = 0$$

Thus, when  $\epsilon$  is sufficiently small, (40) must hold.

Now consider (39). Note that  $\tau_{r_{i+1}^{n-}} - \tau_{r_i^n}$  can be expressed as

$$\tau_{r_{i+1}^{n-}} - \tau_{r_{i}^{n}} = b\left(r_{i+1}^{n} - r_{i}^{n}\right) + \sum_{\substack{r_{i}^{n} < u \le r_{i+1}^{n} \\ \Delta \tau_{u} \le \epsilon_{n}}} \Delta \tau_{u}$$

Thus  $\sup_{r_i^n \leq S_n} |\tau_{r_{i+1}^n} - \tau_{r_i^n}|$  is bounded above by

$$b \sup_{\substack{r_i^n \le S_n}} \left| r_{i+1}^n - r_i^n \right| + \sup_{\substack{r_i \le S_n \\ \Delta \tau_u \le \epsilon_n}} \sum_{\substack{r_i^n < u \le r_{i+1}^n \\ \Delta \tau_u \le \epsilon_n}} \Delta \tau_u$$

It has already been shown that the first term on the right converges to 0 almost surely. Thus to prove (39), it suffices to prove that

$$\lim_{n \to \infty} \sup_{\substack{r_i^n \le S_n \\ r_i^n < u \le r_i^n \\ \Delta \tau_u \le \epsilon_n}} \Delta \tau_u = 0$$
(45)

almost surely, when  $\epsilon_n \downarrow 0$  sufficiently quickly. By the Borel-Cantelli lemma, (45) will follow if  $\epsilon_n$  is chosen such that

$$\mathbb{P}\left(\sup_{\substack{r_i^n \le S_n \\ r_i^n < u \le r_n^n \\ \Delta \tau_u \le \epsilon_n}} \Delta \tau_u \ge \frac{1}{2^n}\right) < \frac{R}{2^n}$$
(46)

for some R > 0.

As before, suppress the superscripts on  $r_i^n$  and the subscripts on  $\epsilon_n$  and  $S_n$ . Recall that  $r_{i+1} - r_i$  are exponential random variables with rate parameter  $g(\epsilon)$ . Furthermore, the jump times of  $\tau_s^n$  are independent of the small-jumps process

$$\tau_s - \tau_s^n = \sum_{\substack{0 \le r \le s \\ \Delta \tau_r \le \epsilon}} \Delta \tau_r.$$

Define  $h(\epsilon)$  by

$$h(\epsilon) = \int_{0}^{\epsilon} t\lambda(dt)$$

Let  $q(\epsilon)$  be the probability that  $\sum_{\substack{r_i < u \leq r_{i+1} \\ \Delta \tau_u \leq \epsilon}} \Delta \tau_u \geq (1/2^n)$ . Define  $\hat{q}(\epsilon)$  as the upper bound on  $q(\epsilon)$  given by Markov's inequality

$$q(\epsilon) = \mathbb{P}\left(\sum_{\substack{r_i < u \le r_{i+1} \\ \Delta \tau_u \le \epsilon}} \Delta \tau_u \ge \frac{1}{2^n}\right) \le 2^n \mathbb{E}\left[\sum_{\substack{r_i < u \le r_{i+1} \\ \Delta \tau_u \le \epsilon}} \Delta \tau_u\right]$$
$$= 2^n \mathbb{E}\left[\int_{r_i}^{r_{i+1}} \int_{0}^{\epsilon} tN(dr, dt)\right] = \frac{2^n h(\epsilon)}{g(\epsilon)} = \hat{q}(\epsilon).$$
(47)

As in the proof of (40), the bound in (46) will be recast as a more tractable inequality.

Let L be the geometric random variable defined by

$$L = \min\left\{i: \sum_{\substack{r_i < u \le r_{i+1} \\ \Delta \tau_u \le \epsilon}} \Delta \tau_u \ge \frac{1}{2^n}\right\}.$$

So L has probability given by  $\mathbb{P}(L = k) = (1 - q(\epsilon))^k q(\epsilon)$ . As in the proof of (40), for any constant M > 0

$$\mathbb{P}\left(\sum_{\substack{r_i < u \le r_{i+1} \\ \Delta \tau_u \le \epsilon}} \Delta \tau_u \ge \frac{1}{2^n}\right)$$
$$= \mathbb{P}(L \le K) \le \mathbb{P}(L \le M) + \mathbb{P}(M \le K).$$

The first term on the right can be bounded as

$$\mathbb{P}(L \le M) = 1 - (1 - q(\epsilon))^{M+1} \le 1 - (1 - \hat{q}(\epsilon))^{M+1}.$$

Furthermore, as in the proof of (40), if

$$(2^n Sg(\epsilon) + 1)\log\left(1 - \hat{q}(\epsilon)\right) < \log\left(1 - \frac{1}{2^n}\right)$$
(48)

the constant M can be chosen such that

$$\mathbb{P}(L \le M) + \mathbb{P}(M \le K) \le \frac{2}{2^n}.$$

Thus, if (48) holds, then so does (46). Note that  $\hat{q}(\epsilon) \to 0$  as  $\epsilon \to 0$ , and thus  $\log(1 - \hat{q}(\epsilon)) \to 0$  as well. Thus, for (48) to hold for sufficiently small  $\epsilon$ , it suffices to show that  $g(\epsilon) \log(1 - \hat{q}(\epsilon)) \to 0$  as  $\epsilon \to 0$ .

The Taylor series expansion of  $\log(1 - \hat{q}(\epsilon))$  gives

$$\begin{split} |g(\epsilon)\log\left(1-\hat{q}(\epsilon)\right)| &= g(\epsilon)\sum_{k=1}^{\infty}\frac{\hat{q}(\epsilon)^{k}}{k} = g(\epsilon)\sum_{k=1}^{\infty}\frac{2^{nk}h(\epsilon)^{k}}{g(\epsilon)^{k}k} \\ &= 2^{n}h(\epsilon)\sum_{k=0}^{\infty}\frac{\hat{q}(\epsilon)^{k}}{k+1} \\ &\leq 2^{n}h(\epsilon)\sum_{k=0}^{\infty}\hat{q}(\epsilon)^{k} = \frac{2^{n}h(\epsilon)}{1-\hat{q}(\epsilon)}. \end{split}$$

Now  $\lim_{\epsilon \to 0} h(\epsilon) = 0$  implies that  $\lim_{\epsilon \to 0} g(\epsilon) \log(1 - \hat{q}(\epsilon)) = 0$ . Therefore (48) holds for small  $\epsilon$  and so (39) is proved.

Now (36) and (37) will be proved. As long as  $S_n \to \infty$ , the limit in (36) is immediate from (38) since

$$S_n - \sup \{r_i^n : r_i^n \le S_n\} \le \sup \{r_{i+1}^n - r_i^n : r_i^n \le S_n\}.$$

Now (37) will be proved. If b > 0, then (37) follows for any sequence with  $S_n \to \infty$ . Thus, assume that b = 0. There exists  $\hat{\epsilon} > 0$  such that  $\lambda((\hat{\epsilon}, \infty)) > 0$ . For all n such that  $\epsilon_n \leq \hat{\epsilon}$ , the following holds:

$$\sup\left\{\tau_{r_i^n}: r_i^n \le S_n\right\} \ge \tau_{S_n}^n \ge \hat{\epsilon} N\left(S_n, (\hat{\epsilon}, \infty)\right).$$

Here  $N(s, (\hat{\epsilon}, \infty))$  is a Poisson process with rate  $\lambda((\hat{\epsilon}, \infty)) > 0$ . Thus, the lower bound goes to  $\infty$  almost surely as  $S_n \to \infty$ . Therefore (37) holds.

# APPENDIX B Proof of Lemma 1

First note that  $\beta$  is analytic at z if  $\int_0^\infty (e^{zt} - 1)\lambda(dt)$  is. Furthermore, by the Lévy-Itô decomposition

$$\mathbb{E}[e^{z\tau_s}] = e^{zbs} \mathbb{E}\left[\exp\left(z\int_0^\infty tN(s,dt)\right)\right].$$

Thus, it suffices to prove the lemma for the case that b = 0.

It will be shown that  $\beta$  is analytic. Let  $y : [0,1] \to \mathbb{C}$  be a continuous, piecewise continuously differentiable curve with  $y(s) \in \text{dom}(\beta)$ , with derivative y'(s) continuous on all but a finite set. Since y is arbitrary, Morera's theorem implies that  $\beta$  is analytic, provided that

$$\oint_C \int_0^\infty \left( e^{yt} - 1 \right) \lambda(dt) dy = 0.$$
(49)

For any  $t \ge 0$ ,  $e^{yt} - 1$  is analytic, and so  $\oint_C (e^{yt} - 1)dy = 0$ . Thus, (49) will hold, provided that the order of integration can be switched. This will be justified by Fubini's theorem.

Since dy(s)/dr = y'(s) almost everywhere, Fubini's theorem can be applied if the following holds:

$$\int_{0}^{1} \int_{0}^{\infty} \left| e^{y(s)t} - 1 \right| |y'(s)| \,\lambda(dt) ds < \infty.$$
 (50)

Since y'(s) is continuous, except on a finite set, it follows that there is a constant M such that |y'(s)| < M almost everywhere. Now  $e^{y(s)t} - 1$  will be bounded. For  $t \in [0, 1]$ , the following bound holds:

$$\left| e^{y(r)t} - 1 \right| \le \sum_{k=1}^{\infty} t^k \left| y(r) \right|^k k! \le t \left( e^{|y(r)|} - 1 \right).$$
 (51)

Let  $r = \sup \{ \operatorname{Rey}(s) : s \in [0, 1] \}$ . By continuity,  $r < r_{\max}$ . The function  $e^{y(s)t} - 1$  can also be bounded using the triangle inequality:

$$\left| e^{y(r)t} - 1 \right| \le \left| e^{y(s)t} \right| + 1 \le e^{rt} + 1.$$
 (52)

Thus, the integral in (50) can be bounded as

$$\int_{0}^{1} \int_{0}^{\infty} \left| e^{y(s)t} - 1 \right| |y'(s)| \,\lambda(dt) ds$$
$$\leq \int_{0}^{1} \left( \int_{0}^{1} t\lambda(dt) \left( e^{|y(r)|} - 1 \right) + \int_{1}^{\infty} (e^{rt} + 1)\lambda(dt) \right) M ds.$$

Note that the integral on the right is finite since  $r < r_{max}$ . Thus, (49) holds and  $\beta$  is analytic.

Since  $\beta$  is analytic, the integral in (5) converges for any matrix A with spec(A)  $\subset$  dom( $\beta$ ). See [21].

Now (4) will be proved. The proof is similar to the proof of Theorem 2.3.8 in [20].

First, the function t will be approximated by step functions over  $(0,\infty)$ . The construction is similar to the approach in the proof of Theorem 1.17 in [42]. Consider a sequence a sequence of integers  $j_n \to \infty$  so that  $\gamma_n = 2^{-j_n} \downarrow 0$  at a rate to be specified later. Let  $k_n(t)$  be the unique integer such that  $k\gamma_n \leq t < (k+1)\gamma_n$ . Define the function  $\varphi_n(t)$  by

$$\varphi_n(t) = \begin{cases} k_n(t)\gamma_n & t \in (0,n) \\ n & t \ge n. \end{cases}$$

Then  $\varphi_n(t)$  is a simple function such that  $\varphi_n(t) = 0$  for  $t \in$  $(0, \gamma_n), t - \gamma_n < \varphi_n(t) \le t$  for  $t \in [\gamma_n, n]$ , and  $\varphi_n(t) \le t$  for t > 0. The formula, (4), is a consequence of the following chain of equalities:

$$\mathbb{E}\left[\exp\left(z\int_{0}^{\infty}tN(s,dt)\right)\right]$$
$$=\lim_{n\to\infty}\mathbb{E}\left[\exp\left(z\int_{0}^{\infty}\varphi_{n}(t)N(s,dt)\right)\right]$$
(53)

$$= \lim_{n \to \infty} \exp\left(s \int_{0}^{\infty} \left(e^{z\varphi_n(t)} - 1\right) \lambda(dt)\right)$$
(54)

$$= \exp\left(s\int_{0}^{\infty} \left(e^{zt} - 1\right)\lambda(dt)\right).$$
(55)

The first equation is the most challenging, and will be handled last. To prove (54), note that  $z\varphi_n(t)$  is a simple function. Thus, there are constants  $c_i \in \mathbb{C}$  and disjoint  $\lambda$ -measurable sets,  $A_i$ , such that

$$z\varphi_n(t) = \sum_{i=1}^q c_i \mathbf{1}_{A_i}(t).$$

Since  $\varphi_n(t) = 0$  over  $(0, \gamma_n)$ , it follows that 0 is not in the closure of any  $A_i$ . Thus, the integral on the right of (53) may be written as

$$\int_{0}^{\infty} z\varphi_n(t)N(s,dt) = \sum_{i=1}^{q} c_i N(s,A_i)$$

where  $N(s, A_i)$  are independent Poisson processes with rate  $\lambda(A_i)$ . Thus, the expectation on the right of (53) may be calculated as

$$\mathbb{E}\left[\exp\left(z\int_{0}^{\infty}\varphi_{n}(t)N(s,dt)\right)\right]$$
$$=\prod_{i=1}^{q}\mathbb{E}\left[\exp\left(c_{i}N(s,A_{i})\right)\right]$$
$$=\prod_{i=1}^{q}\exp\left(-s\lambda(A_{i})\right)\sum_{k=0}^{\infty}\frac{(s\lambda(A_{i}))^{k}}{k!}e^{c_{i}k}$$
$$=\prod_{i=1}^{q}\exp\left(s\lambda(A_{i})\left(e^{c_{i}}-1\right)\right).$$

Now (54) follows because:

$$s\sum_{i=1}^{q} (e^{c_i} - 1) \lambda(A_i) = s\int_{0}^{\infty} \left(e^{z\varphi_n(t)} - 1\right).$$

To prove (55), note that the construction of  $\varphi_n$  implies that  $e^{z\varphi_n(t)} - 1$  is 0 for  $t < \gamma_n$  and constant for  $t \ge n$ . Thus, it is absolutely integrable, so by Lebesgue's dominated convergence theorem

$$\lim_{n \to \infty} \int_{0}^{\infty} \left( e^{z\varphi_n(t)} - 1 \right) \lambda(dt) = \int_{0}^{\infty} (e^{zt} - 1)\lambda(dt).$$

Therefore, (55) holds.

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Now (53) will be proved. First, it will be shown that

$$\lim_{n \to \infty} \int_{0}^{\infty} \varphi_n(t) N(s, dt) = \int_{0}^{t} t N(s, dt), \qquad a.s.$$
(56)

Then, dominated convergence will be applied.

Assume that  $\gamma_{n-1}$  is fixed. The difference of the right and left of (56) may be bounded as

$$0 \leq \int_{0}^{\infty} (t - \varphi_n(t)) N(s, dt)$$
$$\leq \int_{0}^{\gamma_{n-1}} t N(s, dt) + \gamma_n \int_{\gamma_{n-1}}^{n} N(s, dt) + \int_{n}^{\infty} (t - n) N(s, dt).$$
(57)

To bound the first term on the right of (57), note that

$$\int\limits_{0}^{\gamma_{1}}tN(s,dt)=\sum_{i=1}^{\infty}\int\limits_{\gamma_{i+1}}^{\gamma_{i}}tN(s,dt)<\infty, \quad \text{almost surely}.$$

Thus, the first term on the right of (57) may be expressed as the tail sum

$$\int_{0}^{\gamma_{n-1}} tN(s,dt) = \sum_{i=n-1}^{\infty} \int_{\gamma_{i+1}}^{\gamma_i} tN(s,dt)$$

which converges to 0 almost surely, provided that  $\gamma_n \downarrow 0$  sufficiently quickly. (See [20].)

Now consider the second term on the right of (57). For fixed  $\gamma_{n-1}$ , the next term  $\gamma_n$  may be chosen sufficiently small to give the following probability bound:

$$\mathbb{P}\left(\gamma_n \int_{\gamma_{n-1}}^n N(s, dt) \ge 2^{-n}\right)$$
$$= e^{-s\lambda([\gamma_{n-1}, n))} \sum_{k \ge \frac{1}{\gamma_n 2^n}} \frac{\left(s\lambda\left([\gamma_{n-1}, n\right)\right)\right)^k}{k!} < \frac{1}{2^n}.$$

The Borel-Cantelli lemma implies that the second term converges to 0 almost surely.

The last term on the right of (57) is 0 if  $\tau_s < n$ , which holds for sufficiently large n almost surely. Thus (56) holds.

Now it will be shown that Lebesgue's dominated convergence applies to (53). Let r = Re z. Note that the function on the right has magnitude given by

$$\exp\left(r\int_{0}^{\infty}\varphi_{n}(t)N(s,dt)\right).$$
(58)

Thus, it suffices to show that the term on the right has finite expectation. If  $r \leq 0$ , then the term is bounded above by 1 and so finiteness is immediate. So, consider the case that r > 0.

The expectation will be bounded using monotone convergence. Note that  $\varphi_n(t) \leq \varphi_{n+1}(t)$ . Indeed, recall that  $\gamma_n = 2^{-j_n}$  for an integer  $j_n$ . Thus  $\gamma_n = 2^{j_{n+1}-j_n}\gamma_{n+1}$ . By construction  $k_n(t)$  satisfies  $k_n(t)\gamma_n = k_n(t)2^{j_{n+1}-j_n}\gamma_{n+1} \leq t$ . It follows that  $k_n(t)2^{j_{n+1}-j_n} \leq k_{n+1}(t)$ . It follows that  $\varphi_n(t) \leq \varphi_{n+1}(t)$  for all  $t \geq 0$ .

The monotone convergence theorem now implies that (53) holds for z = r. Since (54) and (55) have already been proved, it follows that the magnitude from (58) has expectation bounded by  $e^{s\beta(r)}$ , which is finite.

Finally, (6) will be proved. Since  $\beta$  is analytic on dom( $\beta$ ), it follows that both sides of (4) must be analytic as well. Say that spec(A)  $\subset$  dom( $\beta$ ) and consider a contour C, contained in dom( $\beta$ ), which encloses spec(A). Using the holomorphic functional calculus (see [21]), (6) can be derived as follows:

$$\mathbb{E}\left[e^{A\tau_s}\right] = \frac{1}{2\pi i} \oint_C \mathbb{E}[e^{y\tau_s}](yI - A)^{-1}dy$$
$$= \frac{1}{2\pi i} \oint_C e^{s\beta(y)}(yI - A)^{-1}dy$$
$$= e^{s\beta(A)}.$$

# APPENDIX C Proof of Lemma 2

Define the matrices Z and A by

$$Z = \begin{bmatrix} I \\ 0 \end{bmatrix} P[I \quad 0] \quad \text{and} \quad \tilde{A} = \begin{bmatrix} A & I \\ 0 & 0 \end{bmatrix}$$

respectively. Note that  $e^{\tilde{A}t}$  is given by

$$e^{\tilde{A}t} = \begin{bmatrix} e^{At} & \int_0^t e^{Ar} dr \\ 0 & I \end{bmatrix}.$$

Thus, the matrix-valued mappings may be written as

$$\begin{bmatrix} F(P) & G(P) \\ G(P)^{\mathsf{T}} & H(P) \end{bmatrix} = b(\tilde{A}^{\mathsf{T}}Z + Z\tilde{A}) + \int_{0}^{\infty} \left( e^{\tilde{A}^{\mathsf{T}}t} Z e^{\tilde{A}t} - Z \right) \lambda(dt)$$

Since  $e^{\tilde{A}^{\mathsf{T}}t} \otimes e^{\tilde{A}^{\mathsf{T}}t} = e^{\tilde{A}^{\mathsf{T}} \oplus \tilde{A}^{\mathsf{T}}t}$ , the equation may be vectorized as

$$\operatorname{vec}\left(\begin{bmatrix}F(P) & G(P)\\G(P)^{\mathsf{T}} & H(P)\end{bmatrix}\right)$$
$$=\left(b\tilde{A}^{\mathsf{T}} \oplus \tilde{A}^{\mathsf{T}} + \int_{0}^{\infty} \left(e^{\tilde{A}^{\mathsf{T}} \oplus \tilde{A}^{\mathsf{T}}t} - I\right)\lambda(dt)\right)\operatorname{vec}(Z)$$
$$=\beta(\tilde{A}^{\mathsf{T}} \oplus \tilde{A}^{\mathsf{T}})\operatorname{vec}(Z).$$

Thus according to Lemma 1, F, G, and H are well defined, as long as  $\operatorname{spec}(\tilde{A}^{\mathsf{T}} \oplus \tilde{A}^{\mathsf{T}}) \subset \operatorname{dom}(\beta)$ . By construction, the spectrum is given by

$$\operatorname{spec}(\tilde{A}^{\mathsf{T}} \oplus \tilde{A}^{\mathsf{T}}) = \operatorname{spec}(\tilde{A}^{\mathsf{T}}) + \operatorname{spec}(\tilde{A}^{\mathsf{T}})$$
$$= \{0\} \cup \operatorname{spec}(A) \cup \left(\operatorname{spec}(A) + \operatorname{spec}(A)\right).$$

Let  $r = \max\{\operatorname{Re} \mu : \mu \in \operatorname{spec}(A)\}$ . If  $r \leq 0$ , then the maximum real part of any eigenvalue of  $\tilde{A}^{\mathsf{T}} \oplus \tilde{A}^{\mathsf{T}}$  is 0. If r > 0, then the corresponding maximum real part must be 2r. Since  $\{0\} \cup \operatorname{spec}(2A) \subset \operatorname{dom}(\beta)$ , it follows that  $\operatorname{spec}(\tilde{A}^{\mathsf{T}} \oplus \tilde{A}^{\mathsf{T}}) \subset \operatorname{dom}(\beta)$ , and so the mappings are defined.

Furthermore, the relevant expectations may be vectorized and evaluated using (6)

$$\operatorname{vec}\left(\mathbb{E}\left[e^{\tilde{A}^{\mathsf{T}}\tau_{s}}Ze^{\tilde{A}^{\mathsf{T}}\tau_{s}}\right]\right) = \mathbb{E}\left[e^{\tilde{A}^{\mathsf{T}}\oplus\tilde{A}^{\mathsf{T}}\tau_{s}}\right]\operatorname{vec}(Z)$$
$$= \operatorname{vec}(Z) + s\beta(\tilde{A}^{\mathsf{T}}\oplus\tilde{A}^{\mathsf{T}})\operatorname{vec}(Z)$$
$$+ O(s^{2}).$$

The proof for g is similar, noting that

$$\operatorname{vec}\left(\int_{0}^{t} e^{At} M M^{\mathsf{T}} e^{A^{\mathsf{T}} t} dt\right) = \begin{bmatrix} I & 0 \end{bmatrix} (e^{\hat{A}t} - I) \begin{bmatrix} 0 \\ I \end{bmatrix} \operatorname{vec}(M M^{\mathsf{T}})$$
  
where  $\hat{A} = \begin{bmatrix} A \oplus A & I \\ 0 & 0 \end{bmatrix}$ .

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