MODEL ROBUSTNESS OFFINITESTATE NONLINEAR
FILTERING OVER THE INFINITE TIME HORIZON

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We investigate the robustness of nonlinear filtering for continuous time finite state Markov chains, observed in white noise, with respect to misspecification of the model parameters. It is shown that the distance between the optimal filter and that with incorrect model parameters converges to zero uniformly over the infinite time interval as the misspecified model converges to the true model, provided the signal obeys a mixing condition. The filtering error is controlled through the exponential decay of the derivative of the nonlinear filter with respect to its initial condition. We allow simultaneously for misspecification of the initial condition, of the transition rates of the signal, and of the observation function. The first two cases are treated by relatively elementary means, while the latter case requires the use of Skorokhod integrals and tools of anticipative stochastic calculus.

1. Introduction. The theory of nonlinear filtering concerns the estimation of a signal corrupted by white noise, and has diverse applications in target tracking, signal processing, automatic control, finance, and so on. The basic setting of the theory involves a Markov signal process, for example, the solution of a (nonlinear) stochastic differential equation or a finite-state Markov process, observed in independent corrupting noise. The calculation of the resulting filters is a classical topic in stochastic analysis [14]. Of course, the filtering equations will depend explicitly on the model chosen for the signal process and observations; in almost all realistic applications, however, the model that underlies the filter is only an approximation of the true system that generates the observations. In order for the theory to be practically useful, it is important to establish that the filtered estimates are not too sensitive to the choice of underlying model.

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Continuity with respect to the model parameters of nonlinear filtering estimates on a fixed finite time interval is well established, for example, [3, 4, 10]; generally speaking, it is known that the error incurred in a finite time interval due to the choice of incorrect model parameters can be made arbitrarily small if the model parameters are chosen sufficiently close to those of the true model. As the corresponding error bounds grow rapidly with the length of the time interval, however, such estimates are of little use if we are interested in robustness of the filter over a long period of time. One would like to show that the approximation errors do not accumulate, so that the error remains bounded uniformly over an infinite time interval.

The model robustness of nonlinear filters on the infinite time horizon was investigated in discrete time in [6, 12, 13, 19]. The key idea that allows one to control the accumulation of approximation errors is the asymptotic stability property of many nonlinear filters, which is the focus of much recent work (see [1, 2] and the references therein) and can be summarized as follows. The optimal nonlinear filter is a recursive equation that is initialized with the true distribution of the signal process at the initial time. If the filter is initialized with a different distribution, then the resulting filtered estimates are no longer optimal (in the least-squares sense). The filter is called asymptotically stable if the solution of the wrongly initialized filter converges to the solution of the correctly initialized filter at large times; that is, the filter “forgets” its initial condition after a period of observation.

Using an approximate filter rather than the optimal filter is equivalent to using the optimal filter where we make an approximation error after every time step. Now suppose the optimal filter forgets its initial condition at an exponential rate; then also the approximation error at each time step is forgotten at an exponential rate, and the errors cannot accumulate in time. If the approximation error at each time step is bounded (finite time robustness), then the total approximation error will be bounded uniformly in time. Model robustness on the infinite time horizon is thus a consequence of finite time robustness together with the exponential forgetting property of the filter. This is precisely the method used in [6, 12, 13, 19], and its implementation is fairly straightforward once bounds on the exponential forgetting rate of the filter have been obtained. However, the method used there does not extend to nonlinear filtering in continuous time; even the continuous time model with point process observations studied in [6], though more involved, reduces essentially to discrete (but random) observation times. The continuous time case requires different tools, which we develop in this paper in the setting of nonlinear filtering of a finite-state Markov signal process. (We also mention [7], where a different but related problem is solved.)

We consider the following filtering setup. The signal process \( X = (X_t)_{t \geq 0} \) is a continuous time, homogeneous Markov chain with values in the finite alphabet \( S = \{a_1, \ldots, a_d\} \), transition intensities matrix \( \Lambda = (\lambda_{ij}) \) and initial
distribution $\nu^i = P(X_0 = a_i)$. The observation process $Y = (Y_t)_{t \geq 0}$ is given by
\begin{equation}
Y_t = \int_0^t h(X_s) \, ds + B_t,
\end{equation}
where $h : S \rightarrow \mathbb{R}$ is the observation function [we will also write $h^i = h(a_i)$] and $B$ is a Wiener process that is independent of $X$. The filtering problem for this model concerns the calculation of the conditional probabilities $\pi_t^i = P(X_t = a_i | \mathcal{F}_t^Y)$ from the observations $\{Y_s : s \leq t\}$, where $\mathcal{F}_t^Y = \sigma\{Y_s : s \leq t\}$. It is well known that $\pi_t$ satisfies the Wonham equation \cite{14, 22}
\begin{equation}
d\pi_t = \Lambda^* \pi_t \, dt + (H - h^* \pi_t) \pi_t (dY_t - h^* \pi_t \, dt), \quad \pi_0 = \nu,
\end{equation}
where $x^*$ denotes the transpose of $x$ and $H = \text{diag} \, h$. Note that the Wonham equation is initialized with the true distribution of $X_0$; we will denote by $\pi_t(\mu)$ the solution of the Wonham equation at time $t$ with an arbitrary initial distribution $\pi_0 = \mu$, and by $\pi_{s,t}(\mu)$ the solution of the Wonham equation at time $t \geq s$ with the initial condition $\pi_s = \mu$. In \cite{1, 2} the exponential forgetting property of the Wonham filter was established as follows: the $\ell_1$-distance $|\pi_t(\mu) - \pi_t(\nu)|$ decays exponentially a.s., provided the initial distributions are equivalent $\mu \sim \nu$ and that the mixing condition $\lambda_{ij} > 0 \ \forall i \neq j$ is satisfied. Now consider the Wonham filter with incorrect model parameters:
\begin{equation}
d\tilde{\pi}_t = \tilde{\Lambda}^* \tilde{\pi}_t \, dt + (\tilde{H} - \tilde{h}^* \tilde{\pi}_t) \tilde{\pi}_t (dY_t - \tilde{h}^* \tilde{\pi}_t \, dt), \quad \tilde{\pi}_0 = \nu,
\end{equation}
where $\tilde{\Lambda}$ and $\tilde{h}$ denote a transition intensities matrix and observation function that do not match the underlying signal-observation model $(X,Y)$, $\tilde{H} = \text{diag} \, \tilde{h}$, and we denote by $\tilde{\pi}_t(\mu)$ the solution of this equation with initial condition $\tilde{\pi}_0 = \mu$ and by $\tilde{\pi}_{s,t}(\mu)$ the solution with $\tilde{\pi}_s = \mu$. The following is the main result of this paper.

**Theorem 1.1.** Suppose $\nu^i, \mu^i > 0 \ \forall i$ and $\lambda_{ij}, \tilde{\lambda}_{ij} > 0 \ \forall i \neq j$. Then
\begin{equation}
\sup_{t \geq 0} E\|\tilde{\pi}_t(\mu) - \pi_t(\nu)\|^2 \leq C_1 |\mu - \nu| + C_2 |\tilde{h} - h| + C_3 |\tilde{\Lambda}^* - \Lambda^*|,
\end{equation}
where $|\tilde{\Lambda}^* - \Lambda^*| = \sup\{|(\tilde{\Lambda}^* - \Lambda^*) \tau^i : \tau^i > 0 \ \forall i, |\tau| = 1\}$ and the quantities $C_1, C_2, C_3$ are bounded on any compact subset of parameters $\{(\nu, \Lambda, h, \mu, \tilde{\Lambda}, \tilde{h}) : \nu^i, \mu^i > 0 \ \forall i, |\nu| = |\mu| = 1, \lambda_{ij}, \tilde{\lambda}_{ij} > 0 \ \forall i \neq j, \sum_j \lambda_{ij} = \sum_j \tilde{\lambda}_{ij} = 0 \ \forall i\}$. Additionally we have the asymptotic estimate
\begin{equation}
\limsup_{t \rightarrow \infty} E\|\tilde{\pi}_t(\mu) - \pi_t(\nu)\|^2 \leq C_2 |\tilde{h} - h| + C_3 |\tilde{\Lambda}^* - \Lambda^*|.
\end{equation}
In particular, this implies that if $\nu^i > 0 \ \forall i$, $\lambda_{ij} > 0 \ \forall i \neq j$, then
\begin{equation}
\lim_{(\tilde{h}, \tilde{\Lambda}, \nu) \rightarrow (h, \Lambda, \nu)} \sup_{t \geq 0} E\|\tilde{\pi}_t(\mu) - \pi_t(\nu)\| = \lim_{(\tilde{h}, \tilde{\Lambda}) \rightarrow (h, \Lambda)} \limsup_{t \rightarrow \infty} E\|\tilde{\pi}_t(\mu) - \pi_t(\nu)\| = 0.
\end{equation}
Let us sketch the basic idea of the proof. Rather than considering the Wonham filter, let us demonstrate the idea using the following simple caricature of a filtering equation. Consider a smooth “observation” $y_t$ and a “filter” whose state $x_t$ is propagated by the ordinary differential equation $dx_t/dt = f(x_t, y_t)$. Similarly, we consider the “approximate filter” $d\tilde{x}_t/dt = \tilde{f}(\tilde{x}_t, y_t)$ and assume that everything is sufficiently smooth, so that for fixed $y$ both equations generate a two-parameter flow $x_t = \varphi_{s,t}^y(x_s), \tilde{x}_t = \tilde{\varphi}_{s,t}^y(\tilde{x}_s)$. The following calculation is straightforward:

$$\varphi_{0,t}^y(x) - \varphi_{0,t}^y(x) = \int_0^t \frac{d}{ds} [\varphi_{s,t}^y(\varphi_{0,s}^y(x))] ds$$

$$= \int_0^t D\varphi_{s,t}^y(\varphi_{0,s}^y(x)) \cdot (f(\varphi_{0,s}^y(x), y_s) - \tilde{f}(\varphi_{0,s}^y(x), y_s)) ds,$$

where $D\varphi_{s,t}^y(\varphi_{0,s}^y(x)) \cdot v$ denotes the directional derivative of $\varphi_{s,t}^y(x)$ in the direction $v$. Hence we obtain the following estimate on the approximation error:

$$|\varphi_{0,t}^y(x) - \varphi_{0,t}^y(x)| \leq \int_0^t |D\varphi_{s,t}^y(\varphi_{0,s}^y(x))| \cdot |f(\varphi_{0,s}^y(x), y_s) - \tilde{f}(\varphi_{0,s}^y(x), y_s)| ds.$$

Now suppose that $|f(\cdot, \cdot) - \tilde{f}(\cdot, \cdot)| \leq K$, where $K \to 0$ as $\tilde{f} \to f$; this is an expression of finite-time robustness, as it ensures that $|\varphi_{0,t}^y(x) - \tilde{\varphi}_{0,t}^y(x)| \leq Kt \to 0$ (for fixed $t$) as $\tilde{f} \to f$. Suppose furthermore that we can establish a bound of the form $|D\varphi_{s,t}^y(\cdot)| \leq Ce^{-\lambda(t-s)}$, that is an infinitesimal perturbation to the initial condition is forgotten at an exponential rate. Then the estimate above is uniformly bounded and converges to zero uniformly in time as $\tilde{f} \to f$. Conceptually this is similar to the logic used in discrete time, but we have to replace the exponential forgetting of the initial condition by the requirement that the derivative of the filter with respect to its initial condition decays exponentially.

Returning to the Wonham filter, this procedure can be implemented in a fairly straightforward way if $\tilde{h} = h$. In this case, most of the work involves finding a suitable estimate on the exponential decay of the derivative of the filter with respect to its initial condition; despite the large number of results on filter stability, such estimates are not available in the literature to date. We obtain estimates by adapting methods from [2], together with uniform estimates of the concentration of the optimal filter near the boundary of the simplex.

The general case with $\tilde{h} \neq h$ is significantly more involved. The problem is already visible in the simple demonstration above. Note that the integrand on the right-hand side of the error estimate is not adapted; it depends on the observations on the entire interval $[0, t]$. As the Wonham filter is defined in terms of an Itô-type stochastic integral, this will certainly get us into
trouble. When \( \hat{h} = h \) the stochastic integral cancels in the error bound and the problems are kept to a minimum; in the general case, however, we are in no such luck. Nonetheless this problem is not prohibitive, but it requires us to use the stochastic calculus for anticipating integrands developed by Nualart and Pardoux [15, 16] using Skorokhod integrals rather than Itô integrals and using Malliavin calculus tools.

An entirely different application of the Malliavin calculus to problems of filter stability can be found in [8].

The remainder of this paper is organized as follows. In Section 2 we prove some regularity properties of the solution of the Wonham equation. We also demonstrate the error estimate discussed above in the simpler case \( \hat{h} = h \), and comment on the more general applicability of such a bound. In Section 3 we obtain exponential bounds on the derivative of the Wonham filter with respect to its initial condition. Section 4 treats the general case \( \hat{h} \neq h \) using anticipative stochastic calculus; some of the technical estimates appear in Appendix B. Finally, Appendix A contains a brief review of the results from the Malliavin calculus and anticipative stochastic calculus that are needed in the proofs.

**Notation.** The signal-observation pair \((X, Y)\) is defined on the standard probability space \((\Omega, \mathcal{F}, P)\). The expectation with respect to \(P\) is denoted by \(E\) or sometimes \(E_P\). For \(x \in \mathbb{R}^d\), we denote by \(|x|\) the \(\ell_1\)-norm, by \(\|x\|\) the \(\ell_2\)-norm, and by \(\|x\|_p\) the \(\ell_p\)-norm. We write \(x > y\) (resp. \(<, \geq, \leq\)) if \(x_i > y_i\) (resp. \(<, \geq, \leq\)) \(\forall i\).

The following spaces will be used throughout. Probability distributions on \(\mathbb{S}\) are elements of the simplex \(\Delta^{d-1} = \{x \in \mathbb{R}^d : x \geq 0, |x| = 1\}\). Usually, we will be interested in the interior of the simplex \(S^{d-1} = \{x \in \mathbb{R}^d : x > 0, |x| = 1\}\). The space of vectors tangent to \(S^{d-1}\) is denoted by \(TS^{d-1} = \{x \in \mathbb{R}^d : \sum_i x_i = 0\}\). Finally, we will denote the positive orthant by \(\mathbb{R}^d_{++} = \{x \in \mathbb{R}^d : x > 0\}\).

**2. Preliminaries.** Equation (1.2) is a nonlinear equation for the conditional distribution \(\pi_t\). It is well known however (e.g. [9]) that \(\pi_t\) can also be calculated in a linear fashion: \(\pi_t = \rho_t/|\rho_t|\), where the unnormalized density \(\rho_t\) is propagated by the Zakai equation

\[
d\rho_t = \Lambda^*\rho_t\, dt + H\rho_t\, dY_t, \quad \rho_0 = \nu.
\]

We will repeatedly exploit this representation in what follows. As before \(\rho_t(\mu)\) and \(\rho_{s,t}(\mu)\) \((t \geq s)\) denote the solution of the Zakai equation at time \(t\) with the initial condition \(\rho_0 = \mu\) and \(\rho_s = \mu\), respectively, and \(\pi_{s,t}(\mu) = \rho_{s,t}(\mu)/|\rho_{s,t}(\mu)|\).
We also recall the following interpretation of the norm $|\rho_t|$ of the unnormalized conditional distribution. If we define a new measure $Q \sim P$ through
\begin{equation}
\frac{dP}{dQ} = |\rho_t(\nu)| = |\rho_t|,
\end{equation}
then under $Q$ the observation process $Y_t$ is an $\mathcal{F}_t^Y$-Wiener process. This observation will be used in Section 4 to apply the Malliavin calculus.

The main goal of this section is to establish some regularity properties of the solutions of the Wonham and Zakai equations. In particular, as we will want to calculate the derivative of the filter with respect to its initial condition, we have to establish that $\pi_{s,t}(\mu)$ is in fact differentiable. We will avoid problems at the boundary of the simplex by disposing of it altogether: we begin by proving that if $\mu \in \mathcal{S}^{d-1}$, then a.s. $\pi_{s,t}(\mu) \in \mathcal{S}^{d-1}$ for all times $t > s$.

**Lemma 2.1.** $\mathbb{P}(\rho_{s,t}(\mu) \in \mathbb{R}^{d}_{++} \text{ for all } \mu \in \mathbb{R}^{d}_{++}, 0 \leq s \leq t < \infty) = 1$.

**Proof.** The following variant on the pathwise filtering method reduces the Zakai equation to a random differential equation. First, we write $\Lambda^* = S + T$ where $S$ is the diagonal matrix with $S_{ii} = \lambda_{ii}$. Note that the matrix $T$ has only nonnegative entries. We now perform the transformation $f_{s,t}(\mu) = L_{s,t} \rho_{s,t}(\mu)$ where
\begin{equation}
L_{s,t} = \exp(\left(\frac{1}{2}H^2 - S\right)(t - s) - H(Y_t - Y_s)).
\end{equation}
Then $f_{s,t}(\mu)$ satisfies
\begin{equation}
\frac{df_{s,t}}{dt} = L_{s,t}T L_{s,t}^{-1} f_{s,t}, \quad f_{s,s} = \mu.
\end{equation}
Let $\Omega_c \subset \Omega$, $\mathbb{P}(\Omega_c) = 1$ be a set such that $t \mapsto B_t(\omega)$ is continuous for every $\omega \in \Omega_c$. Then $t \mapsto L_{s,t}$, $t \mapsto L_{s,t}^{-1}$ are continuous in $t$ and have strictly positive diagonal elements for every $\omega \in \Omega_c$. By standard arguments, there exists for every $\omega \in \Omega_c$, $\mu \in \mathbb{R}^d$ and $s \geq 0$ a unique solution $f_{s,t}(\mu)$ to equation (2.3) where $t \mapsto f_{s,t}(\mu)$ is a $C^1$-curve. Moreover, note that $L_{s,t}T L_{s,t}^{-1}$ has nonnegative matrix elements for every $\omega \in \Omega_c$, $s \leq t < \infty$. Hence if $\mu \in \mathbb{R}^d_{++}$ then clearly $f_{s,t}(\mu)$ must be nondecreasing, that is, $f_{s,t} \geq f_{s,r}$ for every $t \geq r \geq s$ and $\omega \in \Omega_c$. But then $\mathbb{R}^d_{++}$ must be forward invariant under equation (2.3) for every $\omega \in \Omega_c$, and as $L_{s,t}$ has strictly positive diagonal elements the result follows. □

**Corollary 2.2.** $\mathbb{P}(\pi_{s,t}(\mu) \in \mathcal{S}^{d-1} \text{ for all } \mu \in \mathcal{S}^{d-1}, 0 \leq s \leq t < \infty) = 1$. 

Let us now investigate the map \( \rho_{s,t}(\mu) \). As this map is linear in \( \mu \), we can write \( \rho_{s,t}(\mu) = U_{s,t}\mu \) a.s. where the \( d \times d \) matrix \( U_{s,t} \) is the solution of

\[
(2.4) \quad dU_{s,t} = \Lambda^* U_{s,t} dt + HU_{s,t} dY_t, \quad U_{s,s} = I.
\]

The following lemma establishes that \( U_{s,t} \) defines a linear stochastic flow in \( \mathbb{R}^d \).

**Lemma 2.3.** For a.e. \( \omega \in \Omega \) (i) \( \rho_{s,t}(\mu) = U_{s,t}\mu \) for all \( s \leq t \); (ii) \( U_{s,t} \) is continuous in \( (s,t) \); (iii) \( U_{s,t} \) is invertible for all \( s \leq t \), where \( U_{s,t}^{-1} \) is given by

\[
(2.5) \quad dU_{s,t}^{-1} = -U_{s,t}^{-1} \Lambda^* dt + U_{s,t}^{-1} H^2 dt - U_{s,t}^{-1} H dY_t, \quad U_{s,s}^{-1} = I;
\]

(iv) \( U_{r,t} U_{s,r} = U_{s,t} \) (and hence \( U_{s,t} U_{r,s}^{-1} = U_{r,t} \)) for all \( s \leq r \leq t \).

**Proof.** Continuity of \( U_{s,t} \) (and \( U_{s,t}^{-1} \)) is a standard property of solution of Lipschitz stochastic differential equations. Invertibility of \( U_{0,t} \) for all \( 0 \leq t < \infty \) is established in [20], page 326, and it is evident that \( U_{s,t} = U_{0,t} U_{0,s}^{-1} \) satisfies equation (2.4). The remaining statements follow, where we can use continuity to remove the time dependence of the exceptional set as in the proof of [20], page 326. \( \square \)

We now turn to the properties of the map \( \pi_{s,t}(\mu) \).

**Lemma 2.4.** The Wonham filter generates a smooth stochastic semiflow in \( \mathcal{S}^{d-1} \), that is, the solutions \( \pi_{s,t}(\mu) \) satisfy the following conditions:

1. For a.e. \( \omega \in \Omega \), \( \pi_{s,t}(\mu) = \pi_{r,t}(\pi_{s,r}(\mu)) \) for all \( s \leq r \leq t \) and \( \mu \).
2. For a.e. \( \omega \in \Omega \), \( \pi_{s,t}(\mu) \) is continuous in \( (s,t,\mu) \).
3. For a.e. \( \omega \in \Omega \), the injective map \( \pi_{s,t}(\cdot) : \mathcal{S}^{d-1} \rightarrow \mathcal{S}^{d-1} \) is \( C^\infty \) for all \( s \leq t \).

**Proof.** For \( x \in \mathbb{R}^d_{++} \) define \( \Sigma(x) = x/|x| \), so that \( \pi_{s,t}(\mu) = \Sigma(\rho_{s,t}(\mu)) \) \( (\mu \in \mathcal{S}^{d-1}) \). Note that \( \Sigma \) is smooth on \( \mathbb{R}^d_{++} \). Hence continuity in \( (s,t,\mu) \) and smoothness with respect to \( \mu \) follow directly from the corresponding properties of \( \rho_{s,t}(\mu) \). The semiflow property \( \pi_{s,t}(\mu) = \pi_{r,t}(\pi_{s,r}(\mu)) \) follows directly from Lemma 2.3. It remains to prove injectivity.

Suppose that \( \pi_{s,t}(\mu) = \pi_{s,t}(\nu) \) for some \( \mu, \nu \in \mathcal{S}^{d-1} \). Then \( U_{s,t}\mu/|U_{s,t}\mu| = U_{s,t}\nu/|U_{s,t}\nu| \) and as \( U_{s,t} \) is invertible we have \( \mu = (|U_{s,t}\mu|/|U_{s,t}\nu|) \nu \). But as \( \mu \) and \( \nu \) must lie in \( \mathcal{S}^{d-1} \), it follows that \( \mu = \nu \). Hence \( \pi_{t}(\cdot) \) is injective. \( \square \)

**Remark 2.5.** The results in this section hold identically if we replace \( \Lambda \) by \( \hat{\Lambda} \), \( h \) by \( \hat{h} \). We will use the obvious notation \( \tilde{\pi}_{s,t}(\mu), \tilde{\rho}_{s,t}(\mu), \tilde{U}_{s,t}, \) and so on.
We finish this section by obtaining an expression for the approximation error in the case \( \tilde{h} = h \); in fact, we will demonstrate the bound for this simple case in a more general setting than is considered in the following. Rather than considering the approximate Wonham filter with modified \( \Lambda \), consider the equation
\[
\frac{d}{dt} \pi_t = f(\pi_t) dt + (H - h^* \pi_t) \pi_t (dY_t - h^* \pi_t dt),
\]
where \( f: \mathcal{S}^{d-1} \to TS^{d-1} \) is chosen in such a way that this equation has a strong solution and \( \inf \{ t > 0 : \pi_t \notin \mathcal{S}^{d-1} \} = \infty \) a.s. In the sequel we consider the case \( f(\pi) = \tilde{\Lambda} \pi \), which clearly satisfies the requirements. We formulate the more general result here, as it might be of interest in other contexts (see Remark 2.8).

**Proposition 2.6.** Let \( \pi_t \) be as above. Then the difference between \( \pi_t \) and the Wonham filter started at \( \mu \) is a.s. given by
\[
\pi_t - \pi_t(\mu) = \int_0^t D\pi_s, \pi_s(f(\pi_s) - \Lambda^* \pi_s) ds,
\]
where \( D\pi_s, \pi_s(\mu) \cdot v \) is the derivative of \( \pi_s, \pi_s(\mu) \) in the direction \( v \in TS^{d-1} \).

**Proof.** Define the (scalar) process \( \Gamma_t \) by
\[
\Gamma_t = \exp \left( \int_0^t h^* \pi_s dY_s - \frac{1}{2} \int_0^t (h^* \pi_s)^2 ds \right).
\]
Using Itô’s rule, we evaluate
\[
\frac{d}{ds}(\Gamma_s U_{s,t}^{-1} \pi_s) = \Gamma_s U_{s,t}^{-1} (f(\pi_s) - \Lambda^* \pi_s).
\]
Multiplying both sides by \( U_{0,t} \), we obtain
\[
\frac{d}{ds}(\Gamma_s U_{s,t} \pi_s) = \Gamma_s U_{s,t} (f(\pi_s) - \Lambda^* \pi_s).
\]
Now introduce as before the map \( \Sigma: \mathbb{R}^d_{++} \to \mathcal{S}^{d-1} \), \( \Sigma(x) = x/|x| \), which is smooth on \( \mathbb{R}^d_{++} \). Define the matrix \( D\Sigma(x) \) with elements
\[
[D\Sigma(x)]^{ij} = \frac{\partial \Sigma^i(x)}{\partial x^j} = \frac{1}{|x|} [\delta_{ij} - \Sigma^i(x)].
\]
Note that \( \Sigma(\alpha x) = \Sigma(x) \) for any \( \alpha > 0 \). Hence
\[
\frac{d}{ds} \Sigma(U_{s,t} \pi_s) = \frac{d}{ds} \Sigma(\Gamma_s U_{s,t} \pi_s) = D\Sigma(\Gamma_s U_{s,t} \pi_s) \frac{d}{ds}(\Gamma_s U_{s,t} \pi_s).
\]
But then we have, using $D\Sigma(\alpha x) = \alpha^{-1}D\Sigma(x)$ ($\alpha > 0$),
\[
\frac{d}{ds} \Sigma(U_{s,t} \tilde{\pi}_s) = D\Sigma(\Gamma_s U_{s,t} \tilde{\pi}_s) \Gamma_s U_{s,t} (f(\tilde{\pi}_s) - \Lambda^* \tilde{\pi}_s) \\
= D\Sigma(U_{s,t} \tilde{\pi}_s) U_{s,t} (f(\tilde{\pi}_s) - \Lambda^* \tilde{\pi}_s).
\]

On the other hand, we obtain from the representation $\pi_{s,t}(\mu) = \Sigma(U_{s,t} \mu)$
\[
D\pi_{s,t}(\mu) \cdot v = D\Sigma(U_{s,t} \mu) U_{s,t} v, \quad \mu \in S^{d-1}, \quad v \in TS^{d-1}.
\]

Note that $f(\tilde{\pi}_s) - \Lambda^* \tilde{\pi}_s \in TS^{d-1}$ as we required that $f : S^{d-1} \to TS^{d-1}$, so that $D\Sigma(U_{s,t} \tilde{\pi}_s) U_{s,t} (f(\tilde{\pi}_s) - \Lambda^* \tilde{\pi}_s) = D\pi_{s,t}(\tilde{\pi}_s) \cdot (f(\tilde{\pi}_s) - \Lambda^* \tilde{\pi}_s)$. Finally, note that
\[
\int_0^t \frac{d}{ds} \Sigma(U_{s,t} \tilde{\pi}_s) \, ds = \Sigma(\tilde{\pi}_t) - \Sigma(U_{0,t} \tilde{\pi}_0) = \tilde{\pi}_t - \pi_t(\mu),
\]
and the proof is complete. □

**Corollary 2.7.** The following estimate holds:
\[
|\tilde{\pi}_t - \pi_t(\mu)| \leq \int_0^t |D\pi_{s,t}(\tilde{\pi}_s)||f(\tilde{\pi}_s) - \Lambda^* \tilde{\pi}_s| \, ds,
\]
where $|D\pi_{s,t}(\mu)| = \sup\{|D\pi_{s,t}(\mu) \cdot v| : v \in TS^{d-1}, \ |v| = 1\}$. Moreover
\[
|\tilde{\pi}_t - \pi_t(\nu)| \leq |\pi_t(\mu) - \pi_t(\nu)| + \int_0^t |D\pi_{s,t}(\tilde{\pi}_s)||f(\tilde{\pi}_s) - \Lambda^* \tilde{\pi}_s| \, ds.
\]

**Remark 2.8.** Corollary 2.7 suggests that the method used here could be applicable to a wider class of filter approximations than those obtained by misspecification of the underlying model. In particular, in the infinite-dimensional setting it is known [5] that by projecting the filter onto a properly chosen finite-dimensional manifold, one can obtain finite-dimensional approximate filters that take a form very similar to equation (2.6). In order to obtain useful error bounds for such approximations one would need to have a fairly tight estimate on the derivative of the filter with respect to its initial condition. Unfortunately, worst-case estimates of the type developed in Section 3 are not sufficiently tight to give quantitative results on the approximation error, even in the finite-state case. In the remainder of the article we will restrict ourselves to studying the robustness problem.

In the following, it will be convenient to turn around the role of the exact and approximate filters in Corollary 2.7, that is, we will use the estimate
\[
|\pi_t(\nu) - \tilde{\pi}_t(\mu)| \leq |\tilde{\pi}_t(\nu) - \tilde{\pi}_t(\mu)| + \int_0^t |D\pi_{s,t}(\tilde{\pi}_s)|(\Lambda^* - \tilde{\Lambda}^*) \tilde{\pi}_s| \, ds,
\]
which holds provided $\tilde{h} = h$. The proof is identical to the one given above.
3. **Exponential estimates for the derivative of the filter.** In order for the bound equation (2.8) to be useful, we must have an exponential estimate for $|D \tilde{\pi}_{s,t}(\cdot)|$. The goal of this section is to obtain such an estimate. We proceed in two steps. First, we use native filtering arguments as in [2] to obtain an a.s. exponential estimate for $|D\pi_{0,t}(\nu)|$. As the laws of the observation processes generated by signals with different initial distributions and jump rates are equivalent, we can extend this a.s. bound to $|D \tilde{\pi}_{s,t}(\mu)|$. We find, however, that the proportionality constant in the exponential estimate depends on $\mu$ and diverges as $\mu$ approaches the boundary of the simplex. This makes a pathwise bound on $|D \tilde{\pi}_{s,t}(\pi_s)|$ difficult to obtain, as $\pi_s$ can get arbitrarily close to the boundary of the simplex on the infinite time interval. Instead, we proceed to find a uniform bound on $E|D \tilde{\pi}_{s,t}(\pi_s)|$.

We begin by recalling a few useful results from [2].

**Lemma 3.1.** Assume $\mu, \nu$ are in the interior of the simplex. Then

$$\pi^i_t(\mu) = \sum_j (\mu^j/\nu^j)P(X_0 = a_j, X_t = a_i | F^Y_t) / \sum_j (\mu^j/\nu^j)P(X_0 = a_j | F^Y_t). \quad (3.1)$$

**Proof.** Define a new measure $P^\mu \sim P$ through

$$\frac{dP^\mu}{dP} = \frac{d\mu}{d\nu}(X_0).$$

It is not difficult to verify that under $P^\mu$, $X_t$ is still a finite-state Markov process with intensities matrix $\Lambda$ but with initial distribution $P^\mu(X_0 = a_i) = \mu^i$. Hence evidently $\pi^i_t(\mu) = P^\mu(X_t = a_i | F^Y_t)$. Using the usual change of measure formula for conditional expectations, we can write

$$\pi^i_t(\mu) = E_{P^\mu}(I_{X_t = a_i} | F^Y_t) = \frac{E(I_{X_t = a_i} (d\mu/d\nu)(X_0) | F^Y_t)}{E((d\mu/d\nu)(X_0) | F^Y_t)}.$$

The result now follows immediately. □

For the proof of the following lemma we refer to [2], Lemma 5.7, page 662.

**Lemma 3.2.** Define $\rho^i_j = P(X_0 = a_j | F^Y_t, X_t = a_i)$. Assume that $\lambda_{ij} > 0 \forall i \neq j$. Then for any $t \geq 0$ we have the a.s. bound

$$\max_{j,k,\ell} |\rho^j_k - \rho^j_\ell| \leq \exp\left(-2t \min_{p,q \neq p} \sqrt{\lambda_{pq}^\lambda_{qp}}\right).$$

We are now ready to obtain some useful estimates.
PROPOSITION 3.3. Let \( \lambda_{ij} > 0 \) \( \forall i \neq j \) and \( \nu \in \mathcal{S}^{d-1} \), \( v \in T\mathcal{S}^{d-1} \). Then a.s.

\[
|D\pi_t(\nu) \cdot v| \leq \sum_k |\nu_k| \frac{\exp\left(-2t \min_{p,q \neq p} \sqrt{\lambda_{pq}\lambda_{qp}}\right)}{\mu^k}.
\]

PROOF. We can calculate directly the directional derivative of (3.1):

\[
(D\pi_t(\mu) \cdot v)^i = \frac{\sum_j (v^j / \nu^j)(P(X_0 = a_j, X_t = a_j|\mathcal{F}_t^Y) - \pi^i_t(\mu)P(X_0 = a_j|\mathcal{F}_t^Y))}{\sum_j (\nu^j / \nu^j)P(X_0 = a_j|\mathcal{F}_t^Y)}.
\]

Setting \( \mu = \nu \), we obtain after some simple manipulations

\[
(D\pi_t(\nu) \cdot v)^i = \pi^i_t(\nu) \sum_{j,k}(v^j / \nu^j) \pi^k_t(\nu) (p^j_{ik} - \tilde{p}^j_{ik}).
\]

The result follows from Lemma 3.2. \( \square \)

To obtain this bound we had to use the true initial distribution \( \nu \), jump rates \( \lambda_{ij} \) and observation function \( h \). However, the almost sure nature of the result allows us to drop these requirements.

COROLLARY 3.4. Let \( \tilde{\lambda}_{ij} > 0 \) \( \forall i \neq j \) and \( \mu \in \mathcal{S}^{d-1} \), \( v \in T\mathcal{S}^{d-1} \). Then a.s.

\[
|D\tilde{\pi}_{s,t}(\mu) \cdot v| \leq \sum_k |\nu_k| \frac{\exp\left(-2(t-s) \min_{p,q \neq p} \sqrt{\tilde{\lambda}_{pq}\tilde{\lambda}_{qp}}\right)}{\mu^k}.
\]

Moreover, the result still holds if \( \mu, v \) are \( \mathcal{F}_s \)-measurable random variables with values a.s. in \( \mathcal{S}^{d-1} \) and \( T\mathcal{S}^{d-1} \), respectively.

PROOF. Note that we can write \( \tilde{\pi}_{0,t}(\mu) = \tilde{P}^\mu(X_t = a_i|\mathcal{F}_t^Y) \), where \( \tilde{P}^\mu \) is the measure under which \( X_t \) has transition intensities matrix \( \tilde{\Lambda} \) and initial distribution \( \mu \), and \( dY_t = \tilde{h}(X_t) dt + dB_t \) where \( dB_t \) is a Wiener process independent of \( X_t \). But \( P^\mu \) and \( P \) are equivalent measures (by the Girsanov theorem and [21], Section IV.22), so that the result for \( s = 0 \) follows trivially from Proposition 3.3. The result for \( s > 0 \) follows directly as the Wonham equation is time homogeneous.

To show that the result still holds when \( \mu, v \) are random, note that \( \tilde{\pi}_{s,t} \) only depends on the observation increments in the interval \([s,t]\), that is, \( D\tilde{\pi}_{s,t}(\mu) \cdot v \) is \( \mathcal{F}_{[s,t]} \)-measurable where \( \mathcal{F}_{[s,t]} = \sigma\{Y_r - Y_s: s \leq r \leq t\} \). Under the equivalent measure \( Q \) introduced in Section 2, \( Y \) is a Wiener process and hence \( \mathcal{F}_{[s,t]} \) and \( \mathcal{F}_s^Y \) are independent. It follows from the bound with constant \( \mu, v \) that

\[
\mathbb{E}_Q(I_{|D\tilde{\pi}_{s,t}(\mu) \cdot v| \leq \epsilon(s)}|\sigma\{\mu, v\}) = 1, \quad Q\text{-a.s.},
\]

where \((*)\) is the right-hand side of (3.2). Hence \( \mathbb{E}_Q(I_{|D\tilde{\pi}_{s,t}(\mu) \cdot v| \leq \epsilon(*)}) = 1 \), and the statement follows from \( P \sim Q \). \( \square \)
Proposition 3.5. Let \( \lambda_{ij} > 0 \ \forall i \neq j \) and \( \mu_1, \mu_2 \in S^{d-1} \). Then a.s.

\[
|\tilde{\pi}_{s,t}(\mu_2) - \tilde{\pi}_{s,t}(\mu_1)| \leq C|\mu_2 - \mu_1| \exp \left( -2(t-s) \min_{p,q \neq p} \sqrt{\lambda_{pq}\lambda_{qp}} \right),
\]
where \( C = \max\{1/\mu_1^k, 1/\mu_2^k : k = 1, \ldots, d\} \).

Proof. Define \( \gamma(u) = \tilde{\pi}_{s,t}(\mu_1 + u(\mu_2 - \mu_1)) \), \( u \in [0,1] \). Then

\[
\tilde{\pi}_{s,t}(\mu_2) - \tilde{\pi}_{s,t}(\mu_1) = \int_0^1 \frac{d\gamma}{du} du = \int_0^1 D\tilde{\pi}_{s,t}(\mu_1 + u(\mu_2 - \mu_1)) \cdot (\mu_2 - \mu_1) du.
\]

We can thus estimate

\[
|\tilde{\pi}_{s,t}(\mu_2) - \tilde{\pi}_{s,t}(\mu_1)| \leq \sup_{u \in [0,1]} |D\tilde{\pi}_{s,t}(\mu_1 + u(\mu_2 - \mu_1)) \cdot (\mu_2 - \mu_1)|.
\]

The result now follows from Corollary 3.4. \( \square \)

Corollary 3.4 and Proposition 3.5 are exactly what we need to establish boundedness of equation (2.8). Note, however, that the right-hand side of (3.2) is proportional to \( 1/\mu^i \), and we must estimate \( |D\tilde{\pi}_{s,t}(\pi_s)| \). Though we established in Section 2 that \( \pi_s \) cannot hit the boundary of the simplex in finite time, it can get arbitrarily close to the boundary during the infinite time interval, thus rendering the right-hand side of equation (3.2) arbitrarily large. If we can establish that \( \sup_{s \geq 0} E(1/\min_k \pi_s^k) < \infty \), however, then we can control \( E|D\tilde{\pi}_{s,t}(\pi_s)| \) to obtain a useful bound.

We begin with an auxiliary integrability property of \( \pi_t \):

Lemma 3.6. Let \( \nu \in S^{d-1} \) and \( T < \infty \). Then

\[
E \int_0^T (\pi_s^i)^{-k} ds < \infty \quad \forall i = 1, \ldots, d, k \geq 1.
\]

Proof. Applying Itô’s rule to the Wonham equation gives

\[
d \log \pi_t^i = \left( \lambda_{ii} - \frac{1}{2}(h^i - h^* \pi_t)^2 \right) dt + \sum_{j \neq i} \lambda_{ij} \frac{\pi_t^j}{\pi_t^i} dt + (h^i - h^* \pi_t) dW_t,
\]

where the innovation \( dW_t = dY_t - h^* \pi_t dt \) is an \( \mathcal{F}_t \)-Wiener process. The application of Itô’s rule is justified by a standard localization argument, as \( \pi_t \) is in \( S^{d-1} \) for all \( t \geq 0 \) a.s. and \( \log x \) is smooth in \( (0,1) \). As \( \lambda_{ij} \geq 0 \) for \( j \neq i \), we estimate

\[
-k \log \pi_t^i \leq -k \log \nu^j - k \lambda_{ii} t + \frac{k}{2} \lambda_{ij} (h^i - h^j)^2 t - k \int_0^t (h^i - h^* \pi_s) dW_s.
\]
But as \( h^i - h^* \pi_t \) is bounded, Novikov’s condition is satisfied and hence
\[
E \exp \left( -k \int_0^t (h^i - h^* \pi_s) \, dW_s - \frac{k^2}{2} \int_0^t (h^i - h^* \pi_s)^2 \, ds \right) = 1.
\]

Estimating the time integral, we obtain
\[
E(\pi_t^i)^{-k} \leq (\nu^i)^{-k} \exp \left( -k \lambda_{ii} t + \frac{1}{2} k (k + 1) \max_j (h^j - h^i)^2 t \right).
\]

The lemma now follows by the Fubini–Tonelli theorem, as \((\pi_s^i)^{-k} \geq 0\) a.s. \(\square\)

We are now in a position to bound \( \sup_{t \geq 0} E(1/\min_i \pi_t^i) \).

**Proposition 3.7.** Let \( \nu \in S^{d-1} \) and suppose that \( \lambda_{ij} > 0 \ \forall i \neq j \). Then
\[
\sup_{t \geq 0} E \left( \frac{1}{\min_i \pi_t^i} \right) < \infty.
\]

**Proof.** By Itô’s rule and using the standard localization argument, we obtain
\[
(\pi_t^i)^{-1} = (\nu^i)^{-1} - \int_0^t \lambda_{ii} (\pi_s^i)^{-1} \, ds - \int_0^t (\pi_s^i)^{-2} \sum_{j \neq i} \lambda_{ji} \pi_s^j \, ds
\]
\[
- \int_0^t (\pi_s^i)^{-1} (h^i - h^* \pi_s) \, dW_s + \int_0^t (\pi_s^i)^{-1} (h^i - h^* \pi_s)^2 \, ds,
\]
where \( W_t \) is the innovations Wiener process. Using Lemma 3.6 we find
\[
E \int_0^t (\pi_s^i)^{-2} (h^i - h^* \pi_s)^2 \, ds \leq \max_j (h^i - h^j)^2 E \int_0^t (\pi_s^i)^{-2} \, ds < \infty,
\]
so the expectation of the stochastic integral term vanishes. Using the Fubini–Tonelli theorem, we can thus write
\[
E((\pi_t^i)^{-1}) = (\nu^i)^{-1} - \int_0^t \lambda_{ii} E((\pi_s^i)^{-1}) \, ds
\]
\[
- \int_0^t E \left( (\pi_s^i)^{-2} \sum_{j \neq i} \lambda_{ji} \pi_s^j \right) \, ds + \int_0^t E((\pi_s^i)^{-1} (h^i - h^* \pi_s)^2) \, ds.
\]
Taking the derivative and estimating each of the terms, we obtain
\[
\frac{dM_t^i}{dt} \leq - \min_{j \neq i} \lambda_{ji} (M_t^i)^2 + \left( |\lambda_{ii}| + \min_{j \neq i} \lambda_{ji} + \max_j (h^i - h^j)^2 \right) M_t^i,
\]
where we have written $M_i^t = \mathbb{E}((\pi_t^i)^{-1})$ and we have used $(M_i^t)^2 \leq \mathbb{E}(\pi_t^i)^{-2}$ by Jensen’s inequality. Using the estimate

$$-K_1^i(M_i^t)^2 + K_2^i M_i^t \leq -K_2^i \tilde{M}_i^t + \frac{(K_2^i)^2}{K_1^i}$$

for $K_1^i > 0$,

we now obtain

$$\frac{dM_i^t}{dt} \leq K_2^i \left( \frac{K_2^i}{K_1^i} - M_i^t \right),$$

where $K_1^i = \min_{j \neq i} \lambda_{ij} > 0$. Consequently we obtain

$$M_i^t \leq e^{-K_2^i t}(\nu^i)^{-1} + \frac{(K_2^i)^2}{K_1^i} e^{-K_2^i t} \int_0^t e^{K_2^i s} ds = e^{-K_2^i t}(\nu^i)^{-1} + \frac{K_2^i}{K_1^i} \left( 1 - e^{-K_2^i t} \right).$$

We can now estimate

$$\sup_{t \geq 0} \mathbb{E} \left( \frac{1}{\min_i \pi_t^i} \right) \leq \sum_{i=1}^d \sup_{t \geq 0} \mathbb{E} \left( \frac{1}{\pi_t^i} \right) \leq \sum_{i=1}^d \left( \frac{1}{\nu^i} \frac{K_2^i}{K_1^i} \right) < \infty,$$

which is what we set out to prove. □

We can now prove Theorem 1.1 for the special case $\tilde{h} = h$. Using equation (2.8), Corollary 3.4, Proposition 3.5 and Proposition 3.7, we obtain

$$\mathbb{E}|\pi_t - \tilde{\pi}_t(\mu)|$$

$$\leq |\mu - \nu| \max_k \left\{ \frac{1}{\mu^k} \vee \frac{1}{\nu^k} \right\} \exp \left( -2 t \min_{p,q \neq p} \sqrt{\lambda_{pq} \lambda_{qp}} \right)$$

$$+ |\Lambda^* - \tilde{\Lambda}^*| \sup_{s \geq 0} \mathbb{E} \left( \frac{1}{\min_i \pi_s^k} \right) \int_0^t \exp \left( -2(t-s) \min_{p,q \neq p} \sqrt{\lambda_{pq} \lambda_{qp}} \right) ds,$$

where $|\Lambda^* - \tilde{\Lambda}^*| = \sup \{|(\Lambda^* - \tilde{\Lambda}^*)\mu| : \mu \in \mathcal{S}^{d-1}\}$. Thus

$$\mathbb{E}|\pi_t - \tilde{\pi}_t(\mu)| \leq |\mu - \nu| \max_k \left\{ \frac{1}{\mu^k} \vee \frac{1}{\nu^k} \right\} e^{-\beta t} + |\Lambda^* - \tilde{\Lambda}^*| \frac{\sup_{s \geq 0} \mathbb{E}(1/\min_k \pi_s^k)}{\beta},$$

where we have written $\beta = 2 \min_{p,q \neq p} (\lambda_{pq} \lambda_{qp})^{1/2}$. The result follows directly using $||\pi_t - \tilde{\pi}_t(\mu)||^2 \leq |\pi_t - \tilde{\pi}_t(\mu)|$ [as $|\pi_t - \tilde{\pi}_t(\mu)| \leq 1$].

4. Model robustness of the Wonham filter. We are now ready to proceed to the general case where the initial density, the transition intensities matrix and the observation function can all be misspecified. The simplicity of the special case $\tilde{h} = h$ that we have treated up to this point is due to the fact that in the calculation of equation (2.7), the stochastic integral term drops out and we can proceed with the calculation using only ordinary calculus.
In the general case we cannot get rid of the stochastic integral, and hence we run into anticipativity problems in the next step of the calculation.

We solve this problem by using anticipative stochastic integrals in the sense of Skorokhod, rather than the usual Itô integral (which is a special case of the Skorokhod integral defined for adapted processes only). Though the Skorokhod integral is more general than the Itô integral in the sense that it allows some anticipating integrands, it is less general in that we have to integrate against a Wiener process (rather than against an arbitrary semimartingale), and that the integrands should be functionals of the driving Wiener process. In our setup, the most convenient way to deal with this is to operate exclusively under the measure \( Q \) of Section 2, under which the observation process \( Y \) is a Wiener process. At the end of the day we can calculate the relevant expectation with respect to the measure \( P \) by using the explicit expression for the Radon–Nikodym derivative \( dP/dQ \). The fact that the integrands must be functionals of the underlying Wiener process is not an issue, as both the approximate and exact filters are functionals of the observations only.

Our setup is further detailed in Appendix A, together with a review of the relevant results from the Malliavin calculus and anticipative stochastic calculus. Below we will use the notation and results from this appendix without further comment. We will also refer to Appendix B for some results on smoothness of the various integrands we encounter; these results are not central to the calculations, but are required for the application of the theory in Appendix A.

We begin by obtaining an anticipative version of Proposition 2.6. Note that this result is precisely of the form one would expect. The first two lines follow the formula for the distance between two flows as one would guess, for example, from the discussion in the Introduction; the last line is an Itô correction term which contains second derivatives of the filter with respect to its initial condition.

**PROPOSITION 4.1.** The difference between \( \pi_t \) and \( \tilde{\pi}_t \) satisfies

\[
\pi_t - \tilde{\pi}_t = \int_0^t D\tilde{\pi}_{r,t}(\pi_r) \cdot \Delta_A \pi_r \, dr + \int_0^t D\tilde{\pi}_{r,t}(\pi_r) \cdot \Delta_H(\pi_r) \, dY_r
\]

\[
- \int_0^t \left[ \partial^2 \tilde{\pi}_{r,t}(\pi_r) \cdot (H - h^* \pi_r) \pi_r - \partial^2 \tilde{\pi}_{r,t}(\pi_r) \cdot (\hat{H} - \hat{h}^* \pi_r) \pi_r \right] \, dr
\]

\[
+ \frac{1}{2} \int_0^t \left[ \partial^2 \tilde{\pi}_{r,t}(\pi_r) \cdot (H - h^* \pi_r) \pi_r - \partial^2 \tilde{\pi}_{r,t}(\pi_r) \cdot (\hat{H} - \hat{h}^* \pi_r) \pi_r \right] \, dr,
\]

where the stochastic integral is a Skorokhod integral and we have written \( \Delta_A = \Delta^* - \Delta^t \), \( \Delta_H(\pi) = (H - h^* \pi) \pi - (\hat{H} - \hat{h}^* \pi) \pi \), and \( \partial^2 \tilde{\pi}_{r,t}(\mu) \cdot v \) is the directional derivative of \( D\tilde{\pi}_{r,t}(\mu) \cdot v \) with respect to \( \mu \in S^{d-1} \) in the direction \( v \in T S^{d-1} \).
Proof. Fix some $T > t$. We begin by evaluating, using Itô’s rule and equation (2.5),

$$
\tilde{U}_{0,s}^{-1} U_{0,s} \nu = \nu + \int_0^s \tilde{U}_{0,r}^{-1}(A^* - \tilde{A}^*) U_{0,r} \nu dr
$$

$$
- \int_0^s \tilde{U}_{0,r}^{-1}(H - \tilde{H}) U_{0,r} \nu dr + \int_0^s \tilde{U}_{0,r}^{-1}(H - \tilde{H}) U_{0,r} \nu dY_r.
$$

Now multiply from the left by $\tilde{U}_{0,t}$; we wish to use Lemma A.5 to bring $\tilde{U}_{0,t}$ into the Skorokhod integral term, that is, we claim that

$$
\tilde{U}_{s,t} U_{0,s} \nu = \tilde{U}_{0,t} \nu + \int_0^s \tilde{U}_{r,t}(A^* - \tilde{A}^*) U_{0,r} \nu dr - \int_0^s \tilde{U}_{r,t} \tilde{H}(H - \tilde{H}) U_{0,r} \nu dr
$$

$$
+ \int_0^s \tilde{U}_{r,t}(H - \tilde{H}) U_{0,r} \nu dY_r + \int_0^s (D_r \tilde{U}_{0,t}) \tilde{U}_{0,r}^{-1}(H - \tilde{H}) U_{0,r} \nu dr.
$$

To justify this expression we need to verify the integrability conditions of Lemma A.5. Note that all matrix elements of $\tilde{U}_{s,t}$ are in $\mathbb{D}_\infty$ for $0 \leq s \leq t < T$, and that

$$
D_r \tilde{U}_{s,t} = \begin{cases} 0, & \text{a.e. } r \notin [s,t], \\ \tilde{U}_{r,t} \tilde{H} U_{s,r}, & \text{a.e. } r \in [s,t]. \end{cases}
$$

This follows directly from Proposition A.4 and Lemma 2.3 (note that the same result holds for $U_{s,t}$ if we replace $\tilde{H}$ by $H$ and $\tilde{U}$ by $U$). Once we plug this result into the expression above, the corresponding integrability conditions can be verified explicitly, see Lemma B.1, and hence we have verified that

$$
\tilde{U}_{s,t} U_{0,s} \nu = \tilde{U}_{0,t} \nu + \int_0^s \tilde{U}_{r,t}(A^* - \tilde{A}^*) U_{0,r} \nu dr + \int_0^s \tilde{U}_{r,t}(H - \tilde{H}) U_{0,r} \nu dY_r.
$$

Next we would like to apply the anticipating Itô rule, Proposition A.6, with the function $\Sigma : \mathbb{R}_+^d \rightarrow \mathbb{S}_+^{d-1}$, $\Sigma(x) = x/|x|$. To this end we have to verify a set of technical conditions, see Lemma B.2. We obtain

$$
\Sigma(\tilde{U}_{s,t} U_{0,s} \nu)
$$

$$
= \Sigma(\tilde{U}_{0,t} \nu) + \int_0^s D \Sigma(\tilde{U}_{r,t} U_{0,r} \nu) \tilde{U}_{r,t}(A^* - \tilde{A}^*) U_{0,r} \nu dr
$$

$$
+ \frac{1}{2} \sum_{k,l} \int_0^s \frac{\partial^2 \Sigma}{\partial x^k \partial x^l}(\nabla_r \tilde{U}_{r,t} U_{0,r} \nu)(\nabla_r \tilde{U}_{r,t} U_{0,r} \nu)^{k}(\tilde{U}_{r,t}(H - \tilde{H}) U_{0,r} \nu)^l dr
$$

$$
+ \int_0^s D \Sigma(\tilde{U}_{r,t} U_{0,r} \nu) \tilde{U}_{r,t}(H - \tilde{H}) U_{0,r} \nu dY_r.
$$

We need to evaluate $\nabla_r \tilde{U}_{r,t} U_{0,r} \nu$. Using Proposition A.2, we calculate

$$
\lim_{\varepsilon \searrow 0} D \tilde{U}_{r,+\varepsilon,t} U_{0,r+\varepsilon} \nu = \lim_{\varepsilon \searrow 0} \tilde{U}_{r,+\varepsilon,t} U_{r,r+\varepsilon} H U_{0,r} \nu = \tilde{U}_{r,t} H U_{0,r} \nu,
$$

$$
\lim_{\varepsilon \searrow 0} D \tilde{U}_{r,t} U_{0,r} \nu = \tilde{U}_{r,t} U_{0,r} H U_{0,r} \nu.
$$

$$
\lim_{\varepsilon \searrow 0} D \tilde{U}_{r,t} U_{0,r} \nu = \tilde{U}_{r,t} U_{0,r} H U_{0,r} \nu.
$$

$$
\lim_{\varepsilon \searrow 0} D \tilde{U}_{r,t} U_{0,r} \nu = \tilde{U}_{r,t} U_{0,r} H U_{0,r} \nu.
$$

$$
\lim_{\varepsilon \searrow 0} D \tilde{U}_{r,t} U_{0,r} \nu = \tilde{U}_{r,t} U_{0,r} H U_{0,r} \nu.
$$
and similarly
\[
\lim_{\varepsilon \to 0} D_r \bar{U}_{r-\varepsilon, t} U_{0, r-\varepsilon} \nu = \lim_{\varepsilon \to 0} \bar{U}_{r, t} \bar{H} \bar{U}_{r-\varepsilon, t} U_{0, r-\varepsilon} \nu = \bar{U}_{r, t} \bar{H} U_{0, r} \nu.
\]

After some rearranging, we obtain
\[
\Sigma(\bar{U}_{s, t} U_{0, s} \nu) = \Sigma(U_{0, t} \nu) + \int_0^s D \Sigma(U_{r, t} U_{0, r} \nu) \bar{U}_{r, t} (\Lambda^* - \tilde{\Lambda}^*) U_{0, r} \nu \, dr
\]
\[
+ \frac{1}{2} \sum_{k, \ell} \int_0^s \frac{\partial^2 \Sigma}{\partial x^k \partial x^\ell}(\bar{U}_{r, t} U_{0, r} \nu)(\bar{U}_{r, t} H U_{0, r} \nu)^k(\bar{U}_{r, t} H U_{0, r} \nu)^\ell \, dr
\]
\[
- \frac{1}{2} \sum_{k, \ell} \int_0^s \frac{\partial^2 \Sigma}{\partial x^k \partial x^\ell}(\bar{U}_{r, t} U_{0, r} \nu)(\bar{U}_{r, t} \bar{H} U_{0, r} \nu)^k(\bar{U}_{r, t} \bar{H} U_{0, r} \nu)^\ell \, dr
\]
\[
+ \int_0^s D \Sigma(\bar{U}_{r, t} U_{0, r} \nu) \bar{U}_{r, t} (H - \bar{H}) U_{0, r} \nu \, dY_r.
\]

From this point onward we will set \(s = t\). We will need (on \(\mathbb{R}^d_+\))
\[
D^2 \Sigma^{ik\ell}(x) = \frac{\partial^2 \Sigma(x)}{\partial x^k \partial x^\ell} = -\frac{1}{|x|} (D \Sigma^{ik}(x) + D \Sigma^{i\ell}(x)).
\]

Recall that \(D \Sigma(\alpha x) = \alpha^{-1} D \Sigma(x)\); it follows that also \(D^2 \Sigma(\alpha x) = \alpha^{-2} D^2 \Sigma(x)\) for \(\alpha > 0\). Using these expressions with \(\alpha = |U_{0, r} \nu|\), we get
\[
\pi_t - \bar{\pi}_t = \int_0^t D \Sigma(\bar{U}_{r, t} \pi_r) \bar{U}_{r, t} \Delta \pi_r \, dr + \int_0^t D \Sigma(\bar{U}_{r, t} \pi_r) \bar{U}_{r, t} (H - \bar{H}) \pi_r \, dY_r
\]
\[
+ \frac{1}{2} \sum_{k, \ell} \int_0^t \frac{\partial^2 \Sigma}{\partial x^k \partial x^\ell}(\bar{U}_{r, t} \pi_r)(\bar{U}_{r, t} H \pi_r)^k(\bar{U}_{r, t} H \pi_r)^\ell \, dr
\]
\[
- \frac{1}{2} \sum_{k, \ell} \int_0^t \frac{\partial^2 \Sigma}{\partial x^k \partial x^\ell}(\bar{U}_{r, t} \pi_r)(\bar{U}_{r, t} \bar{H} \pi_r)^k(\bar{U}_{r, t} \bar{H} \pi_r)^\ell \, dr.
\]

Next we want to express the integrands in terms of \(D \bar{\pi}_{r, t}(\pi_r) \cdot v\), and so on, rather than in terms of \(D \Sigma(x)\). Recall that \(D \bar{\pi}_{r, t}(\pi_r) \cdot v = D \Sigma(\bar{U}_{r, t} \pi_r) \bar{U}_{r, t} v\) when \(v \in TS^{d-1}\). Similar terms appear in the expression above, but, for example, \(H \pi_r \notin TS^{d-1}\). To rewrite the expression in the desired form, we use that \(D \Sigma(\bar{U}_{r, t} \pi_r) \bar{U}_{r, t} \pi_r = 0\). Hence
\[
D \Sigma(\bar{U}_{r, t} \pi_r) \bar{U}_{r, t} \bar{H} \pi_r = D \Sigma(\bar{U}_{r, t} \pi_r) \bar{U}_{r, t} (\bar{H} - \bar{h}^* \pi_r) \pi_r
\]
\[
= D \bar{\pi}_{r, t}(\pi_r) \cdot (\bar{H} - \bar{h}^* \pi_r) \pi_r
\]
and similarly for the other terms. Note also that
\[
\sum_k D^2 \Sigma^{ik\ell}(\bar{U}_{r, t} \pi_r)(\bar{U}_{r, t} \pi_r)^k = -D \Sigma^{i\ell}(\bar{U}_{r, t} \pi_r).
\]
Substituting this into the expression for $\pi_t - \tilde{\pi}_t$ and rearranging, we obtain

$$
\pi_t - \tilde{\pi}_t = \int_0^t D\pi_{r,t}(\pi_r) \cdot \Delta^r \pi_r \, dr + \int_0^t D\pi_{r,t}(\pi_r) \cdot \Delta H(\pi_r) \, dY_r
- \int_0^t D\tilde{\pi}_{r,t}(\pi_r) \cdot [h^* \pi_r (H - h^* \pi_r) \pi_r - \tilde{h}^* \pi_r (\tilde{H} - \tilde{h}^* \pi_r) \pi_r] \, dr
+ \frac{1}{2} \sum_{k,\ell} \int_0^t \frac{\partial^2 \Sigma}{\partial x^k \partial x^\ell} (\tilde{U}_{r,t} \pi_r)(\tilde{U}_{r,t} (H - h^* \pi_r) \pi_r)^k (\tilde{U}_{r,t} (H - h^* \pi_r) \pi_r)^\ell \, dr
- \frac{1}{2} \sum_{k,\ell} \int_0^t \frac{\partial^2 \Sigma}{\partial x^k \partial x^\ell} (\tilde{U}_{r,t} \pi_r)(\tilde{U}_{r,t} (\tilde{H} - \tilde{h}^* \pi_r) \pi_r)^k (\tilde{U}_{r,t} (\tilde{H} - \tilde{h}^* \pi_r) \pi_r)^\ell \, dr.
$$

It remains to note that we can write

$$
(D^2 \pi_{s,t}(\mu) \cdot v)^i = \sum_{k,\ell} D^2 \Sigma^{ik\ell} (\tilde{U}_{s,t} \mu)(\tilde{U}_{s,t} v)^k (\tilde{U}_{s,t} v)^\ell.
$$

The result follows immediately. □

**Remark 4.2.** We have allowed misspecification of most model parameters of the Wonham filter. One exception is the observation noise intensity: we have not considered observations of the form $dY_t = h(X_t) \, dt + \sigma \, dB_t$ with $\sigma \neq 1$; in other words, the quadratic variation of $Y_t$ is assumed to be known $[Y, Y]_t = t$. We do not consider this a significant drawback as the quadratic variation can be determined directly from the observation process $Y_t$. On the other hand, the model parameters $\nu, \Lambda, h$ are “hidden” and would have to be estimated, making these quantities much more prone to modeling errors.

If we allow misspecification of $\sigma$, we would have to be careful to specify in which way the filter is implemented: in this case, the normalized solution of the misspecified Zakai equation no longer coincides with the solution of the misspecified Wonham equation. Hence one obtains a different error estimate depending on whether the normalized solution of the misspecified Zakai equation, or the solution of the misspecified Wonham equation, is compared to the exact filter. Both cases can be treated using similar methods, but we do not pursue this here.

Let $e_t = \pi_t - \tilde{\pi}_t$. We wish to estimate the norm of $e_t$. Unfortunately, we can no longer use the triangle inequality as in Section 2 due to the presence of the stochastic integral; instead, we choose to calculate $\|e_t\|^2$, which is readily estimated.
Lemma 4.3. The filtering error can be estimated by
\[ \mathbb{E}_P \| e_t \|^2 \leq \int_0^t \mathbb{E}_P |D\tilde{\pi}_{r,t}(\pi_r) \cdot \Delta_A \pi_r| \, dr + K \int_0^t \mathbb{E}_P |D\tilde{\pi}_{r,t}(\pi_r) \cdot \Delta_H(\pi_r)| \, dr \\
+ \int_0^t \mathbb{E}_P |D\tilde{\pi}_{r,t}(\pi_r) \cdot (h^* \pi_r (H - h^* \pi_r) \pi_r - \tilde{h}^* \pi_r (\tilde{H} - \tilde{h}^* \pi_r) \pi_r)| \, dr \\
+ \frac{1}{2} \int_0^t \mathbb{E}_P |D^2\tilde{\pi}_{r,t}(\pi_r) \cdot (H - h^* \pi_r) \pi_r - D^2\tilde{\pi}_{r,t}(\pi_r) \cdot (\tilde{H} - \tilde{h}^* \pi_r) \pi_r| \, dr, \]

where \( K = 2 \max_k |h^k| + \max_k |\tilde{h}^k| \).

Proof. We wish to calculate \( \mathbb{E}_P \| e_t \|^2 = \mathbb{E}_P e_t^* e_t \). Using Proposition 4.1, we obtain
\[ \mathbb{E}_P \| e_t \|^2 = \int_0^t \mathbb{E}_P e_t^* D\tilde{\pi}_{r,t}(\pi_r) \cdot \Delta_A \pi_r \, dr \\
+ \mathbb{E}_P \left[ e_t^* \int_0^t D\tilde{\pi}_{r,t}(\pi_r) \cdot \Delta_H(\pi_r) \, dY_r \right] \\
- \int_0^t \mathbb{E}_P e_t^* D\tilde{\pi}_{r,t}(\pi_r) \cdot [h^* \pi_r (H - h^* \pi_r) \pi_r - \tilde{h}^* \pi_r (\tilde{H} - \tilde{h}^* \pi_r) \pi_r] \, dr \\
+ \frac{1}{2} \int_0^t \mathbb{E}_P e_t^* [D^2\tilde{\pi}_{r,t}(\pi_r) \cdot (H - h^* \pi_r) \pi_r \pi_r \\
- D^2\tilde{\pi}_{r,t}(\pi_r) \cdot (\tilde{H} - \tilde{h}^* \pi_r) \pi_r] \, dr. \]
The chief difficulty is the stochastic integral term. Using equation (2.2), we can write
\[ \mathbb{E}_P \left[ e_t^* \int_0^t D\tilde{\pi}_{r,t}(\pi_r) \cdot \Delta_H(\pi_r) \, dY_r \right] \\
= \mathbb{E}_Q \left[ |U_{0,t} \nu| e_t^* \int_0^t D\tilde{\pi}_{r,t}(\pi_r) \cdot \Delta_H(\pi_r) \, dY_r \right]. \]
We would like to apply equation (A.1) to evaluate this expression. First, we must establish that the integrand is in \( \text{Dom} \delta \); this does not follow directly from Proposition 4.1, as the anticipative Itô rule which was used to obtain that result can yield integrands which are only in \( L^{1,2}_{\text{loc}} \). We can verify directly, however, that the integrand in this case is indeed in \( \text{Dom} \delta \), see Lemma B.3. Next, we must establish that \( |U_{0,t} \nu| e_t^* \) is in \( D^{1,2} \) for every \( i \). Note that \( |U_{0,t} \nu| = \sum_i (U_{0,t} \nu)^i \), so \( |U_{0,t} \nu| \) is in \( D^\infty \). Moreover, we establish in Lemma B.4 that \( e_t \in D^{1,2} \) and that \( D_r e_t \) is a bounded random variable for every \( t \).
Hence it follows from Proposition A.1 that \( |U_{0,t}\nu| e_i^* \in D^{1,2} \). Consequently we can apply equation (A.1), and we obtain
\[
E_Q \left[ |U_{0,t}\nu| e_i^* \int_0^t D\tilde{\pi}_{r,t}(\pi_r) \cdot \Delta_H(\pi_r) dY_r \right]
\]
\[
= \int_0^t E_Q \left[ |U_{0,t}\nu| D_r e_i^* + D_r |U_{0,t}\nu| e_i^* D\tilde{\pi}_{r,t}(\pi_r) \cdot \Delta_H(\pi_r) \right] dr
\]
\[
= \int_0^t E_Q \left[ |U_{0,t}\nu| (D_r\pi_t - D_r\tilde{\pi}_t)^* D\tilde{\pi}_{r,t}(\pi_r) \cdot \Delta_H(\pi_r) \right] dr
\]
\[
+ \int_0^t E_Q \left[ \sum_i (U_{r,t}HU_{0,r}\nu)^i e_i^* D\tilde{\pi}_{r,t}(\pi_r) \cdot \Delta_H(\pi_r) \right] dr.
\]
Now note that \( |e_i| \leq 1 \), and that by Lemma B.4
\[
|(D_r\pi_t - D_r\tilde{\pi}_t)^i| \leq |(D_r\pi_t)^i| + |(D_r\tilde{\pi}_t)^i| \leq \max_k |h^k| + \max_k |\tilde{h}^k|.
\]
Furthermore we can estimate
\[
\left| \sum_i (U_{r,t}HU_{0,r}\nu)^i \right| \leq \frac{1}{|U_{0,t}\nu|} \sum U_{r,t}^j |h^j| U_{0,r}^{j,k} \leq \max_k |h^k|,
\]
where we have used a.s. nonnegativity of the matrix elements of \( U_{0,r} \) and \( U_{r,t} \) (this must be the case, as, for example, \( U_{r,t}\nu \) has nonnegative entries for any vector \( \nu \) with nonnegative entries). Hence we obtain
\[
E_Q \left[ |U_{0,t}\nu| e_i^* \int_0^t D\tilde{\pi}_{r,t}(\pi_r) \cdot \Delta_H(\pi_r) dY_r \right]
\]
\[
\leq \left( 2 \max_k |h^k| + \max_k |\tilde{h}^k| \right) \int_0^t E_Q |U_{0,t}\nu| |D\tilde{\pi}_{r,t}(\pi_r) \cdot \Delta_H(\pi_r)| dr.
\]
The result follows after straightforward manipulations. \( \square \)

Unlike in the case \( \tilde{h} = h \), we now have to deal also with second derivatives of the filter with respect to its initial condition. These can be estimated much in the same way as we dealt with the first derivatives.

**Lemma 4.4.** Let \( \tilde{\lambda}_{ij} > 0 \) \( \forall i \neq j \) and \( \mu \in S^{d-1} \), \( v, w \in TS^{d-1} \). Then a.s.
\[
|D^2\tilde{\pi}_{s,t}(\mu) \cdot v - D^2\tilde{\pi}_{s,t}(\mu) \cdot w|
\]
\[
\leq 2 \sum_k \frac{|v^k + w^k|}{\mu^k} \sum_j \frac{|v^j - w^j|}{\mu^j} \exp \left( -2(t-s) \min_{p,q \neq p} \sqrt{\lambda_{pq} \tilde{\lambda}_{qp}} \right).
\]
Moreover, the result still holds if \( \mu, v, w \) are \( \mathcal{F}_s^Y \)-measurable random variables with values a.s. in \( S^{d-1} \) and \( TS^{d-1} \), respectively.
PROOF. Proceeding as in the proof of Proposition 3.3, we can calculate directly the second derivative of (3.1):

\[(D^2 \pi_t(\mu) \cdot v)^i = -2(D\pi_t(\mu) \cdot v)^i \sum_j (\mu^j/\nu^j)P(X_0 = a_j | \mathcal{F}_t^Y) \sum_j (\mu^j/\nu^j)P(X_0 = a_j | \mathcal{F}_t^Y).\]

Setting \(\mu = \nu\) and using the triangle inequality, we obtain

\[|D^2 \pi_t(\nu) \cdot v - D^2 \pi_t(\nu) \cdot w| \leq 2 \sum_{i,j} \frac{|\nu^j(D\pi_t(\nu) \cdot v)^i - w^j(D\pi_t(\nu) \cdot w)^i|}{\nu^j}.\]

Another application of the triangle inequality and using Proposition 3.3 gives

\[|D^2 \pi_t(\nu) \cdot v - D^2 \pi_t(\nu) \cdot w| \leq \sum_k \frac{|v^k + w^k|}{\nu^k} |D\pi_t(\nu) \cdot (v - w)| + \sum_k \frac{|v^k - w^k|}{\nu^k} |D\pi_t(\nu) \cdot (v + w)| \leq 2 \sum_k \frac{|v^k + w^k|}{\nu^k} \sum_j \frac{|v^j - w^j|}{\nu^j} \exp\left(-2t \min_{p,q \neq p} \sqrt{\lambda_{pq} \lambda_{qp}}\right).\]

We can now repeat the arguments of Corollary 3.4 to establish that the result still holds if we replace \(\pi_{0,t}\) by \(\tilde{\pi}_{s,t}\) and \(v, w\) by \(\mathcal{F}_s^Y\)-measurable random variables \(\mu, v, w\). This completes the proof.

We are now ready to complete the proof of Theorem 1.1.

PROOF OF THEOREM 1.1. Set \(\beta = 2 \min_{p,q \neq p} (\lambda_{pq} \lambda_{qp})^{1/2}\). Let us collect all the necessary estimates. First, we have

\[\int_0^t E_P[D\tilde{\pi}_{r,t}(\pi_r) \cdot \Delta_D \pi_r] dr \leq \beta^{-1} \sup_{s \geq 0} E_P\left(1/ \min_k \pi^k_s\right) |\Lambda^* - \tilde{\Lambda}^*|,\]

as we showed in Section 3. Next, we obtain

\[\int_0^t E_P[D\tilde{\pi}_{r,t}(\pi_r) \cdot \Delta_H(\pi_r)] dr \leq \beta^{-1} \sup_{\pi \in S^{d-1}} \sum_k |h^k - \tilde{h}^k + \tilde{h}^* \pi - h^* \pi|\]

using Corollary 3.4. Using the triangle inequality, we can estimate this by

\[\int_0^t E_P[D\tilde{\pi}_{r,t}(\pi_r) \cdot \Delta_H(\pi_r)] dr \leq (d + 1) \beta^{-1} |h - \tilde{h}|.\]

Next, we estimate using Corollary 3.4

\[\int_0^t E_P[D\tilde{\pi}_{r,t}(\pi_r) \cdot (h^* \pi_r(H - h^* \pi_r) \pi_r - \tilde{h}^* \pi_r(\tilde{H} - \tilde{h}^* \pi_r) \pi_r)] dr\]
\[
\leq \beta^{-1} \sup_{\pi \in S^{d-1}} \sum_k |\pi^* (h^k - h^* \pi) - \hat{\pi}^* (\hat{h}^k - \hat{h}^* \pi)| \\
\leq \beta^{-1} \left((d + 1) \max_k |h^k| + d \max_{k,\ell} |\hat{h}^k - \hat{h}^\ell| \right) |h - \hat{h}|,
\]
where we have used the estimate
\[
\sum_k |\pi^* (h^k - h^* \pi) - \hat{\pi}^* (\hat{h}^k - \hat{h}^* \pi)| \\
\leq |\pi^*| \sum_k |h^k - \hat{h}^k + \hat{\pi}^* - h^* \pi| + |\pi^* - \hat{\pi}^*| \sum_k |\hat{h}^k - \hat{h}^* \pi| \\
\leq (d + 1) \max_k |h^k| |h - \hat{h}| + |h - \hat{h}| \sum_k |\hat{h}^k - \hat{h}^* \pi| \\
\leq \left((d + 1) \max_k |h^k| + d \max_{k,\ell} |\hat{h}^k - \hat{h}^\ell| \right) |h - \hat{h}|.
\]
Next we estimate using Lemma 4.4
\[
\frac{1}{2} \int_0^t \mathbb{E}_P |D^2 \pi_{r,t}(\pi_r) \cdot (H - h^* \pi_r) \pi_r - D^2 \hat{\pi}_{r,t}(\pi_r) \cdot (\hat{H} - \hat{h}^* \pi_r) \pi_r| \, dr \\
\leq \beta^{-1} \sup_{\pi \in S^{d-1}} \sum_k |h^k - h^* \pi + \hat{h}^k - \hat{h}^* \pi| \sum_j |\hat{h}^j - \hat{h}^j + \hat{h}^* \pi - h^* \pi| \\
\leq d(d + 1) \beta^{-1} \left(\max_{k,\ell} |h^k - h^\ell| + \max_{k,\ell} |\hat{h}^k - \hat{h}^\ell| \right) |h - \hat{h}|.
\]
We have now estimated all the terms in Lemma 4.3, and hence we have bounded \( \mathbb{E}_P \|e_t\|^2 = \mathbb{E}_P \|\pi_t(\nu) - \hat{\pi}_t(\nu)\|^2 \). It remains to allow for misspecified initial conditions. To this end, we estimate
\[
\|\pi_t(\nu) - \hat{\pi}_t(\mu)\|^2 \\
\leq \|e_t\|^2 + \|\hat{\pi}_t(\nu) - \hat{\pi}_t(\mu)\| \|\pi_t(\nu) - \hat{\pi}_t(\mu)\| + 2\|\pi_t(\nu) - \hat{\pi}_t(\nu)\|.
\]
Hence we obtain using the equivalence of finite-dimensional norms \( \|x\| \leq K_{21} |x| \)
\[
\|\pi_t(\nu) - \hat{\pi}_t(\mu)\|^2 \leq \|e_t\|^2 + 6K_{21} \|\pi_t(\nu) - \hat{\pi}_t(\mu)\|
\]
where we have used that the simplex is contained in the \((d - 1)\)-dimensional unit sphere, so \( \|\mu_1 - \mu_2\| \leq 2 \forall \mu_1, \mu_2 \in \Delta^{d-1} \). The statement of the theorem now follows directly from Lemma 4.3, Proposition 3.5 and the estimates above. \( \square \)
APPENDIX A: ANTICIPATIVE STOCHASTIC CALCULUS

The goal of this appendix is to recall briefly the main results of the Malliavin calculus, Skorokhod integrals and anticipative stochastic calculus that are needed in the proofs. In our application of the theory we wish to deal with functionals of the observation process \( (Y_t)_{t \in [0,T]} \), where \( T \) is some finite time (usually we will calculate integrals from 0 to \( t \), so we can choose any \( T > t \)). Recall from Section 2 that \( Y \) is an \( \mathcal{F}^Y_t \)-Wiener process under the measure \( Q \); it will thus be convenient to work always under \( Q \), as this puts us directly in the framework used, for example, in [15]. As the theory described below is defined \( Q \)-a.s. and as \( P \sim Q \), the corresponding properties under \( P \) are unambiguously obtained by using equation (2.2). We will presume this setup whenever the theory described here is applied.

A smooth random variable \( F \) is one of the form \( f(Y(h_1),\ldots,Y(h_n)) \), where \( Y(h) \) denotes the Wiener integral of the deterministic function \( h \in L^2([0,T]) \) with respect to \( Y \) and \( f \) is a smooth function which is of polynomial growth together with all its derivatives. For smooth \( F \) the Malliavin derivative \( D F \) is defined by

\[
D_t F = \sum_{i=1}^n \frac{\partial f}{\partial x^i}(Y(h_1),\ldots,Y(h_n))h_i(t).
\]

The Malliavin derivative \( D \) can be shown [15], page 26, to be closeable as an operator from \( L^p(\Omega, \mathcal{F}^Y_T, Q) \) to \( L^p(\Omega, \mathcal{F}^Y_T, Q; L^2([0,T])) \) for any \( p \geq 1 \), and we denote the domain of \( D \) in \( L^p(\Omega) \) by \( \mathbb{D}^{1,p} \) [for notational convenience we will drop the measure \( Q \) and \( \sigma \)-algebra \( \mathcal{F}^Y_T \) throughout this section, where it is understood that \( L^p(\Omega) \) denotes \( L^p(\Omega, \mathcal{F}^Y_T, Q) \), etc.]. More generally, we consider iterated derivatives \( D^k F \in L^p(\Omega; L^2([0,T]^k)) \) defined by \( D_{t_1,\ldots,t_k} F = D_{t_1} \cdots D_{t_k} F \), and the domain of \( D^k \) in \( L^p(\Omega) \) is denoted by \( \mathbb{D}^{k,p} \). The domains \( \mathbb{D}^{k,p} \) can also be localized ([15], pages 44–45), and we denote the corresponding localized domains by \( \mathbb{D}^{k,p}_{\text{loc}} \). Finally, we define the useful class \( \mathbb{D}^\infty = \bigcap_{p \geq 1} \bigcap_{k \geq 1} \mathbb{D}^{k,p} \).

We will use two versions of the chain rule for the Malliavin derivative.

**Proposition A.1.** Let \( \varphi: \mathbb{R}^m \to \mathbb{R} \) be \( C^1 \) and \( F = (F^1,\ldots,F^m) \) be a random vector with components in \( \mathbb{D}^{1,2} \). Then \( \varphi(F) \in \mathbb{D}^{1,2}_{\text{loc}} \) and

\[
D \varphi(F) = \sum_{i=1}^m \frac{\partial \varphi}{\partial x^i}(F) D F^i.
\]

If \( \varphi(F) \in L^2(\Omega) \) and \( D \varphi(F) \in L^2(\Omega \times [0,T]) \), then \( \varphi(F) \in \mathbb{D}^{1,2} \). These results still hold if \( F \) a.s. takes values in an open domain \( V \subset \mathbb{R}^m \) and \( \varphi \) is \( C^1(V) \).
The first (local) statement is [16], Proposition 2.9; the second statement can be proved in the same way as [17], Lemma A.1, and the proofs are easily adapted to the case where $F$ a.s. takes values in some domain. The next result is from [15], page 62:

**Proposition A.2.** Let $\varphi: \mathbb{R}^m \to \mathbb{R}$ be a smooth function which is of polynomial growth together with all its derivatives, and let $F = (F^1, \ldots, F^m)$ be a random vector with components in $D^\infty$. Then $\varphi(F) \in D^\infty$ and the usual chain rule holds. This implies that $D^\infty$ is an algebra, that is, $FG \in D^\infty$ for $F, G \in D^\infty$.

The following result follows from [15], page 32 (here $[s, t]^c = [0, T]\setminus[s, t]$).

**Lemma A.3.** If $F \in D^{1,2}$ is $\mathcal{F}^Y_s$-measurable, then $DF = 0$ a.e. in $\Omega \times [s, t]^c$.

It is useful to be able to calculate explicitly the Malliavin derivative of the solution of a stochastic differential equation. Consider $dx_t = f(x_t)\, dt + \sigma(x_t)\, dY_t$, $x_0 \in \mathbb{R}^m$, where $f(x)$ and $\sigma(x)$ are smooth functions of $x$ with bounded derivatives of all orders. It is well known that such equations generate a smooth stochastic flow of diffeomorphisms $x_t = \xi_t(x)$ [11]. We now have the following result.

**Proposition A.4.** All components of $x_t$ belong to $D^\infty$ for every $t \in [0, T]$. We have $D_x x_t = D\xi_t(x_0)D\xi_r(x_0)^{-1}\sigma(x_r) \ a.e. \ r < t$, where $(D\xi_t(x))^ij = \partial\xi^j_t(x)/\partial x^i$ is the Jacobian matrix of the flow, and $D_r x_t = 0 \ a.e. \ r > t$.

The first statement is given in [15], Theorem 2.2.2, page 105, the second on [15], equation (2.38), page 109, the third follows from adaptedness (Lemma A.3).

We now consider $D$ as a closed operator from $L^2(\Omega)$ to $L^2(\Omega \times [0, T])$ with domain $D^{1,2}$. Its Hilbert space adjoint $\delta = D^*$ is well defined in the usual sense as a closed operator from $L^2(\Omega \times [0, T])$ to $L^2(\Omega)$, and we denote its domain by $\text{Dom} \delta$. The operator $\delta$ is called the Skorokhod integral, and coincides with the Itô integral on the subspace $L^2_0(\Omega \times [0, T]) \subset \text{Dom} \delta$ of adapted square integrable processes ([15], Proposition 1.3.4, page 41). $\delta$ is thus an extension of the Itô integral to a class of possibly anticipative integrands. To emphasize this point we will write

$$\delta(uI_{[s, t]}) = \int_s^t u_r\, dY_r, \quad uI_{[s, t]} \in \text{Dom} \delta.$$
The Skorokhod integral has the following properties. First, its expectation vanishes $\mathbb{E}_Q \delta (u) = 0$ if $u \in \text{Dom } \delta$. Second, by its definition as the adjoint of $D$ we have

\begin{equation}
\mathbb{E}_Q(F \delta (u)) = \mathbb{E}_Q \left[ \int_0^T (D_t F) u_t \, dt \right]
\end{equation}

if $u \in \text{Dom } \delta$, $F \in \mathbb{D}^{1,2}$. We will also use the following result, the proof of which proceeds in exactly the same way as its one-dimensional counterpart ([15], page 40).

**Lemma A.5.** If $u$ is an $n$-vector of processes in $\text{Dom } \delta$ and $F$ is an $m \times n$-matrix of random variables in $\mathbb{D}^{1,2}$ such that $\mathbb{E}_Q \int_0^T \| F u_t \|_2^2 \, dt < \infty$, then

$$
\int_0^T F u_t \, dY_t = F \int_0^T u_t \, dY_t - \int_0^T (D_t F) u_t \, dt
$$

in the sense that $F u \in \text{Dom } \delta$ iff the right-hand side of this expression is in $L^2(\Omega)$.

As it is difficult to obtain general statements for integrands in $\text{Dom } \delta$, it is useful to single out restricted classes of integrands that are easier to deal with. To this end, define the spaces $\mathbb{L}^{k,p} = L^p([0,T]; \mathbb{D}^{k,p})$ for $k \geq 1$, $p \geq 2$. Note that $\mathbb{L}^{k,p} \subset \mathbb{L}^{1,2} \subset \text{Dom } \delta$ ([15], page 38). Moreover, the domains $\mathbb{L}^{k,p}$ can be localized to $\mathbb{L}_{\text{loc}}^{k,p}$ ([15], pages 43–45). We can now state an Itô change of variables formula for Skorokhod integrals, see [15, 16, 18]. The extension to processes that a.s. take values in some domain is straightforward through localization.

**Proposition A.6.** Consider an $m$-dimensional process of the form

$$
x_t = x_0 + \int_0^t v_s \, ds + \int_0^t u_s \, dY_s,
$$

where we assume that $x_t$ has a continuous version and $x_0 \in (\mathbb{D}_{\text{loc}}^{1,4})^m$, $v \in (\mathbb{L}_{\text{loc}}^{1,4})^m$, and $u \in (\mathbb{L}_{\text{loc}}^{2,4})^m$. Let $\varphi : \mathbb{R}^m \to \mathbb{R}$ be a $C^2$ function. Then

$$
\varphi(x_t) = \varphi(x_0) + \int_0^t D_\varphi(x_s) v_s \, ds + \int_0^t D_\varphi(x_s) u_s \, dY_s
$$

$$
+ \frac{1}{2} \int_0^t (D^2 \varphi(x_s) \nabla_x x_s, u_s) \, ds,
$$

where $\nabla_x x_s = \lim_{\varepsilon \to 0} \mathbb{D}_s (x_{s+\varepsilon} + x_{s-\varepsilon})$, $D_\varphi(x_s) u_s = \sum_i (\partial \varphi / \partial x^i)(x_s) u_s^i$, $(D^2 \varphi(x_s) \nabla_x x_s, u_s) = \sum_{ij} (\partial^2 \varphi / \partial x^i \partial x^j)(x_s) u_s^i \nabla_x x_s^j$. The result still holds if $x_s$ a.s. takes values in an open domain $V \subset \mathbb{R}^m \forall s \in [0,t]$ and $\varphi$ is $C^2(V)$.
APPENDIX B: SOME TECHNICAL RESULTS

Lemma B.1. The following equality holds:
\[
\tilde{U}_{0,t} \int_0^s \tilde{U}_{0,r}^{-1}(H - \tilde{H})U_{0,r}\nu dY_r = \int_0^s \tilde{U}_{r,t}(H - \tilde{H})U_{0,r}\nu dY_r + \int_0^s \tilde{U}_{r,t}\tilde{H}(H - \tilde{H})U_{0,r}\nu dr.
\]
The integral on the left-hand side is an Itô integral, on the right-hand side a Skorokhod integral.

Proof. We have already established in the proof of Proposition 4.1 that the matrix elements of \(\tilde{U}_{0,t}\) are in \(D_{\infty} \subset D_{1,2}\). Moreover,
\[
E_Q \| \tilde{U}_{r,t}(H - \tilde{H})U_{0,r}\nu \|^2 \\
\leq \| H - \tilde{H} \|^2 E_Q(\| \tilde{U}_{r,t} \|^2 \| U_{0,r} \|^2) \\
\leq \| H - \tilde{H} \|^2 \sqrt{E_Q \| \tilde{U}_{r,t} \|^4 E_Q \| U_{0,r} \|^4} \\
\leq C_4^4 \| H - \tilde{H} \|^2 \sqrt{E_Q \| \tilde{U}_{r,t} \|^4 E_Q \| U_{0,r} \|^4},
\]
where we have used the Cauchy–Schwarz inequality and \(\| \nu \| \leq 1\) for \(\nu \in S_{d-1}\). Here \(\| U \|_p = (\sum_{ij} U_{ij}^p)^{1/p}\) is the elementwise \(p\)-norm of \(U\), \(\| U \|\) is the usual matrix 2-norm, and \(C_p\) matches the norms \(\| U \| \leq C_p\| U \|_p\) (recall that all norms on a finite-dimensional space are equivalent). As \(U_{0,r}, \tilde{U}_{r,t}\) are solutions of linear stochastic differential equations, standard estimates give for any integer \(p \geq 2\)
\[
E_Q \left( \sup_{0 \leq r \leq t} \| \tilde{U}_{r,t} \|_p^p \right) \leq D_1(p) < \infty,
\]
and we obtain
\[
\int_0^s E_Q \| \tilde{U}_{r,t}(H - \tilde{H})U_{0,r}\nu \|^2 dr \leq s \sup_{0 \leq r \leq s} E_Q \| \tilde{U}_{r,t}(H - \tilde{H})U_{0,r}\nu \|^2 < \infty.
\]
Hence we can apply Lemma A.5 to obtain the result. By a similar calculation we can establish that the right-hand side of the expression in Lemma A.5 for our case is square integrable, so that the Skorokhod integral is well defined.

□

Lemma B.2. The anticipating Itô rule with \(\Sigma(x) = x/|x|\) can be applied to
\[
\tilde{U}_{s,t}U_{0,s}\nu = \tilde{U}_{0,t}\nu + \int_0^s \tilde{U}_{r,t}(\Lambda^* - \tilde{\Lambda}^*)U_{0,r}\nu dr + \int_0^s \tilde{U}_{r,t}(H - \tilde{H})U_{0,r}\nu dY_r.
\]
proof. Clearly the Skorokhod integral term has a.s. continuous sample paths, as both \( \tilde{U}_{s,t} U_{0,t} \) and the time integrals do; moreover, \( \tilde{U}_{0,t} \nu \in (\mathbb{D}^\infty)^d \). In order to be able to apply Proposition A.6, it remains to check the technical conditions \( v_r = \tilde{U}_{r,t}(\Lambda^* - \Lambda^*) U_{0,r} \nu \in (\mathbb{L}^1,^4)^d \), \( u_r = \tilde{U}_{r,t}(H - \tilde{H}) U_{0,r} \nu \in (\mathbb{L}^2,^4)^d \).

As \( \mathbb{D}^\infty \) is an algebra, \( u_t \) and \( v_t \) take values in \( \mathbb{D}^\infty \). Moreover, we can establish exactly as in the proof of Lemma B.1 that \( u \) and \( v \) are in \( L^3(\Omega \times [0, t]) \). To prove the proof we must establish that

\[
\sum_i \int_0^t E_Q \left[ \int_0^t (D_s u_r^i)^2 \, ds \right]^2 \, dr < \infty,
\]

\[
\sum_i \int_0^t E_Q \left[ \int_0^t (D_s v_r^i)^2 \, ds \right]^2 \, dr < \infty,
\]

thus ensuring that \( u, v \in (\mathbb{L}^1,^4)^d \), and

\[
\sum_i \int_0^t E_Q \left[ \int_0^t \left( D_s D_s u_r^i \right)^2 \, ds \right] \, dr < \infty
\]

which ensures that \( u \in (\mathbb{L}^2,^4)^d \). Using the Cauchy–Schwarz inequality we have

\[
\sum_i \int_0^t E_Q \left[ \int_0^t (D_s u_r^i)^2 \, ds \right]^2 \, dr \leq t \int_0^t \int_0^t E_Q \| D_s u_r \|^4 \, ds \, dr \leq t^3 \sup_{0 \leq r, s \leq t} E_Q \| D_s u_r \|^4,
\]

and similarly for \( v \). Moreover, we obtain

\[
\sum_i \int_0^t E_Q \left[ \int_0^t \left( D_s D_s u_r^i \right)^2 \, ds \right] \, dr \leq t^5 \sup_{0 \leq r, s, \sigma \leq t} E_Q \| D_s D_s u_r \|^4.
\]

But using the chain rule Proposition A.2 we can easily establish that

\[
D_s u_r = \begin{cases} 
\tilde{U}_{r,t}(H - \tilde{H}) U_{s,r} H U_{0,s} \nu, & \text{a.e. } 0 < s < r < t, \\
\tilde{U}_{s,t} H U_{r,s}(H - \tilde{H}) U_{0,r} \nu, & \text{a.e. } 0 < r < s < t,
\end{cases}
\]

and similarly

\[
D_\sigma D_s u_r = \begin{cases} 
\tilde{U}_{r,t}(H - \tilde{H}) U_{s,r} H U_{\sigma,s} H U_{0,\sigma} \nu, & \text{a.e. } 0 < \sigma < s < r < t, \\
\tilde{U}_{r,t}(H - \tilde{H}) U_{\sigma,s} H U_{s,r} H U_{0,s} \nu, & \text{a.e. } 0 < s < r < \sigma < t, \\
\tilde{U}_{s,t} H U_{r,s}(H - \tilde{H}) U_{\sigma,r} H U_{0,\sigma} \nu, & \text{a.e. } 0 < \sigma < r < s < t, \\
\tilde{U}_{s,t} H U_{\sigma,s} H U_{r,s}(H - \tilde{H}) U_{0,r} \nu, & \text{a.e. } 0 < r < s < \sigma < t, \\
\tilde{U}_{\sigma,t} H U_{r,s}(H - \tilde{H}) U_{0,r} \nu, & \text{a.e. } 0 < r < s < \sigma < t.
\end{cases}
\]

The desired estimates now follow as in the proof of Lemma B.1. \( \square \)
LEMMA B.3. The Skorokhod integrand obtained by applying the anticipative Itô formula as in Lemma B.2 is in $\text{Dom} \, \delta$.

PROOF. We use the notation $\rho_r = U_{0,r} \nu$. The Skorokhod integral in question is

$$\int_0^s \left| D \Sigma(\tilde{U}_{r,t}\rho_r)\tilde{U}_{r,t}(H - \tilde{H}) \rho_r \right| dY_r = \int_0^s f_r \, dY_r.$$ 

To establish $f \in \text{Dom} \, \delta$, it suffices to show that $f \in L^{1,2}$. We begin by showing

$$|D \Sigma(\tilde{U}_{r,t}\rho_r)\tilde{U}_{r,t}(H - \tilde{H}) \rho_r|$$

$$= \sum_i \left| \sum_{j,k} \delta_{ij} - \Sigma^i(\tilde{U}_{r,t}\rho_r)\tilde{U}_{r,t}^j h^k - \tilde{h}^k \rho_r \right|$$

$$\leq \frac{1}{|\tilde{U}_{r,t}\rho_r|} \sum_{i,j,k} \tilde{U}_{r,t}^j h^k - \tilde{h}^k | \rho_r$$

$$\leq \max_k |h^k - \tilde{h}^k| \sum_{i,j,k} \tilde{U}_{r,t}^j h^k$$

$$= d \max_k |h^k - \tilde{h}^k|,$$

where we have used the triangle inequality, $|\delta_{ij} - \Sigma^i(x)| \leq 1$ for any $x \in \mathbb{R}^2$, and the fact that $\tilde{U}_{r,t}$ and $\rho_r$ have nonnegative entries a.s. Hence $f_r$ is a bounded process. Similarly, we will show that $D_s f_r$ is a bounded process. Note that $f_r$ is a smooth function on $\mathbb{R}_+^d$ of positive random variables in $\mathbb{D}^\infty$; hence we can apply the chain rule Proposition A.1. This gives

$$\left( D_s f_r \right)^i = \begin{cases} 
\sum_{jk} D^2 \Sigma^{ijk}(\tilde{U}_{r,t}\rho_r)(\tilde{U}_{r,t}(H - \tilde{H}) \rho_r)^j(\tilde{U}_{r,t} U_{s,r} H \rho_s)^k \\
+ \sum_j D \Sigma^{ij}(\tilde{U}_{r,t}\rho_r)(\tilde{U}_{r,t}(H - \tilde{H}) U_{s,r} H \rho_s)^j, & \text{a.e. } s < r, \\
\sum_{jk} D^2 \Sigma^{ijk}(\tilde{U}_{r,t}\rho_r)(\tilde{U}_{r,t}(H - \tilde{H}) \rho_r)^j(\tilde{U}_{s,t} \tilde{H} U_{r,s} \rho_r)^k \\
+ \sum_j D \Sigma^{ij}(\tilde{U}_{r,t}\rho_r)(\tilde{U}_{s,t} \tilde{H} U_{r,s}(H - \tilde{H}) \rho_r)^j, & \text{a.e. } s > r.
\end{cases}$$

Proceeding exactly as before, we find that $D f \in L^\infty(\Omega \times [0,1]^2)$. But then by Proposition A.1 we can conclude that $D_s f_r \in \mathbb{D}^{1,2}$ for a.e. $(s,t) \in [0,1]^2$, and in particular $f \in \mathbb{D}^{1,2}$. Hence the proof is complete. \hfill \Box

LEMMA B.4. $D_r \pi_s = D \pi_{r,s}(\pi_r) : (H - h^s \pi_r) \pi_r$ a.e. $r < s$, $D_r \pi_s = 0$ a.e. $r > s$. Moreover $|\left( D_r \pi_s \right)^i| \leq \max_k |h^k|$ for every $i$. The equivalent results hold for $D_r \bar{\pi}_s$. In particular, this implies that $\pi_s$ and $\bar{\pi}_s$ are in $\mathbb{D}^{1,2}$. 


Proof. The case $r > s$ is immediate from adaptedness of $\pi_s$. For $r < s$, apply the chain rule to $\pi_s = \Sigma(U_{0,s} \nu) \in \mathbb{D}^{1,2}_{\text{loc}}$. Boundedness of the resulting expression follows, for example, as in the proof of Lemma B.3, and hence it follows that $\pi_s \in \mathbb{D}^{1,2}$. \hfill \Box

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