

Volatility and Correlation Analysis of Financial Market Data

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Overview

- ① Introduction
- ② High-Frequency Prices and Inference
- ③ Open, Close, High, Low Prices
 - Galerkin solution
 - Small-time Analytic Solution and Gap-Fill Approximation

Data type: High-Frequency Prices

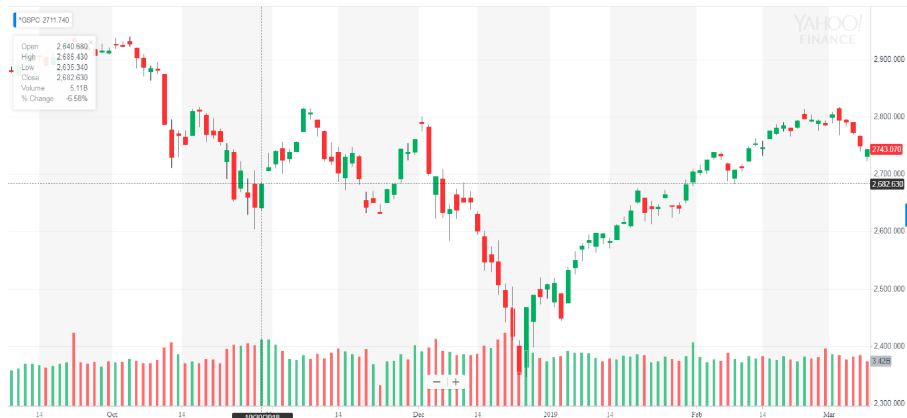
- Longitudinal view

Time-stamp	Volume	Price
2011-08-01 03:02:52.112434	1000	179.6400
2011-08-01 03:02:53.752456	1000	179.6200
2011-08-01 03:02:53.900010	1000	179.6300
2011-08-01 03:02:54.103493	1000	179.6050
2011-08-01 03:02:54.343493	1000	179.6700

- Cross-sectional view: the bid-ask spread

Volume	Price
101	179.6400
203	179.6200
305	179.6100
500	179.6000

Data type: Open, Close, High, Low prices



High-frequency returns and difficulties

- As sampling frequency approaches transaction-by-transaction frequency, **irregular spacing** between transactions and **discreteness** in transaction prices (such as the bid-ask spread) become dominant features of the data [10].
- These confounding effects in estimating volatility are generally termed **microstructure noise**, which is
 - independent of sampling interval
 - dominant over short time (< 5 minutes)
 - not cumulative (prices at later times are not affected)

Past work and model-based incoherence

- Attempts have been made to fit discrete time state-space models to high-frequency data and estimate price volatility [4, 2]

$$\begin{aligned}y_{t+\Delta t} &= F_t \theta_t + C_t W_{t,1} \\ \theta_{t+\Delta t} &= G_t \theta_t + V_t W_{t,1}.\end{aligned}$$

- As formulated, such types of models are not robust with respect to the specification of the sampling interval [7, 2, 12].
- Inference over shrinking sampling interval Δt leads to **incoherence**.
- Most recent work has focused on the **realized variance** [6] estimator:

$$RV_T := \sum_{i=1}^{N/\Delta t} r_{t_i}^2,$$

where $r_{t_i}^2$ is the squared *return* observed over $[t_{i-1}, t_i]$.

Realized variance estimator

In the absence of microstructure noise,

$$RV_t = \sum_{i=1}^{N/\Delta t} r_{t_i}^2 \rightarrow \int_0^T \sigma^2(X_t, t) dt$$

- Corrections for the presence of microstructure noise have been proposed in the form of:
 - sampling sparser grids and averaging [11],
 - combining estimators based on sub-sampled data at different frequencies [1],
 - kernel-based estimation [8, 3].
- Disadvantages: averages out information; no inherent model-based projection forward in time.

A filtering-based approach

modeling microstructure noise

Given the **true** price S_j at sample index j , we model the **measured** price P_j as contaminated by noise due to the bid-ask spread D :

$$P_j = S_j + \nu_j, \quad \nu_j \sim U[-D/2, D/2],$$

A first-order Taylor expansion of the **observed log price** $\log(P_j)$ produces

$$\log(P_j) \approx \log(S_j) + \frac{1}{S_j} \nu_j.$$

We further model the noise term with a Gaussian distribution

$$\zeta_j := \frac{1}{S_j} \nu_j \sim N\left(0, \frac{D}{4Q^2}\right)$$

where Q is an order-magnitude approximation of the true price S_j .

Stochastic evolution of price and volatility

the idealized joint log-price $\log(\hat{S}_t)$ and log-volatility $\log(\hat{\sigma}_{t,1}), \log(\hat{\sigma}_{t,2})$ diffusion process as the system of SDEs:

$$\begin{aligned}d \log(\hat{S}_t) &= \hat{\mu} dt + \sqrt{\hat{\sigma}_{t,1} \hat{\sigma}_{t,2}} \sqrt{dt} \hat{e}_t + dJ_t, \\d \log(\hat{\sigma}_{t,1}) &= -\hat{\theta}_1(\log(\hat{\sigma}_{t,1}) - \hat{\alpha}) dt + \hat{\tau}_1 \sqrt{dt} \hat{e}_{t,1}, \\d \log(\hat{\sigma}_{t,2}) &= -\hat{\theta}_2(\log(\hat{\sigma}_{t,2}) - \hat{\alpha}) dt + \hat{\tau}_2 \sqrt{dt} \hat{e}_{t,2}.\end{aligned}$$

- $\log(\hat{\sigma}_{t,1})$ evolves stochastically with long time scale $1/\hat{\theta}_1$,
- $\log(\hat{\sigma}_{t,2})$ evolves stochastically with short time scale $1/\hat{\theta}_2$,
- dJ_t is a compound Poisson process modeling jumps,
- $\hat{e}_t, \hat{e}_{t,1}, \hat{e}_{t,2}$ are Normally distributed with

$$\mathbb{E}[\hat{e}_t \hat{e}_{t,2}] = \rho, \quad \mathbb{E}[\hat{e}_t \hat{e}_{t,2}] = 0, \quad \mathbb{E}[\hat{e}_{t,1} \hat{e}_{t,2}] = 0.$$

Discrete interpretation of the formulation

Defining the log of the observed price P_j as $Y_j := \log(P_j)$

$$Y_j = \log(S_j) + \zeta_j,$$

$$\log(S_j) = \log(S_{j-1}) + \mu(\Delta) + \sqrt{\sigma_{j,1}\sigma_{j,2}} \epsilon_j + J_j(\Delta),$$

$$\log(\sigma_{j+1,1}) = \alpha(\Delta) + \theta_1(\Delta) \{\log(\sigma_{j,1}) - \alpha(\Delta)\} + \tau_1(\Delta) \epsilon_{j,1},$$

$$\log(\sigma_{j+1,2}) = \alpha(\Delta) + \theta_2(\Delta) \{\log(\sigma_{j,2}) - \alpha(\Delta)\} + \tau_2(\Delta) \epsilon_{j,2}.$$

$$\sigma_{j+1,i} = \hat{\sigma}_{(j+1)\Delta,i} \sqrt{\Delta}, \quad S_j = \hat{S}_{j\Delta}, \quad J_j(\Delta) = J((j+1)\Delta) - J(j\Delta)$$

$$\alpha(\Delta) = \hat{\alpha} + \frac{1}{2} \log(\Delta), \quad \mu(\Delta) = \hat{\mu} \Delta,$$

$$\tau_i(\Delta) = \hat{\tau}_i \sqrt{\frac{1 - \exp\{-2\hat{\theta}_i \Delta\}}{2\hat{\theta}_i}}, \quad \theta_i(\Delta) = \exp\{-\hat{\theta}_i \Delta\}.$$

Prior formulation

To ensure that the model estimation is coherent across choices of sampling periods Δ , we

- ① define the first two moments of each continuous-time parameter,
- ② use the Delta Method to calculate the first two moments of each discrete-time parameter,
- ③ define the prior for each discrete-time parameter.

Computation

- We separate the standard innovation term from volatility factors $\sigma_{j,1}$ and $\sigma_{j,2}$ by considering the transformation

$$\begin{aligned} \log(S_j) &= \log(S_{j-1}) + \mu(\Delta) + \sqrt{\sigma_{j,1}\sigma_{j,2}} \epsilon_j + J_j(\Delta) \\ \Leftrightarrow \underbrace{\log[|\log(S_j/S_{j-1}) - \mu(\Delta) - J_j(\Delta)|]}_{y_j^*} &= \frac{1}{2} \underbrace{\log(\sigma_{j,1})}_{h_{j,1}} + \frac{1}{2} \underbrace{\log(\sigma_{j,2})}_{h_{j,2}} \\ &\quad + \underbrace{\log(\epsilon_j^2)/2}_{\epsilon_j^*}. \end{aligned}$$

Computation

- We separate the standard innovation term from volatility factors $\sigma_{j,1}$ and $\sigma_{j,2}$ by considering the transformation

$$\begin{aligned} \log(S_j) &= \log(S_{j-1}) + \mu(\Delta) + \sqrt{\sigma_{j,1}\sigma_{j,2}} \epsilon_j + J_j(\Delta) \\ \Leftrightarrow \underbrace{\log[|\log(S_j/S_{j-1}) - \mu(\Delta) - J_j(\Delta)|]}_{y_j^*} &= \frac{1}{2} \underbrace{\log(\sigma_{j,1})}_{h_{j,1}} + \frac{1}{2} \underbrace{\log(\sigma_{j,2})}_{h_{j,2}} \\ &\quad + \underbrace{\log(\epsilon_j^2)/2}_{\epsilon_j^*}. \end{aligned}$$

- We approximate ϵ_j^* as a mixture of Normals

$$\epsilon_j^* = \log(\epsilon_j^2)/2 \sim \sum_{l=1}^{10} p_l N\left(\frac{m_l}{2}, \frac{v_l^2}{4}\right).$$

Computation, continued

The joint distribution $(\epsilon_j^*, \epsilon_{j,2})$ conditional on the latent mixture element γ_j becomes

$$\begin{aligned}
 p(\epsilon_j^*, \epsilon_{j,2} | \gamma_j) &= p(\epsilon_{j,2} | \epsilon_j^*, \gamma_j) p(\epsilon_j^* | \gamma_j) \\
 &= p(\epsilon_{j,2} | \underbrace{d_j \exp(\epsilon_j^*)}_{\epsilon_j}, \gamma_j) p(\epsilon_j^* | \gamma_j) \\
 &= N\left(\epsilon_{j,2} \mid d_j \exp(\epsilon_j^*), (1 - \rho^2)\right) N\left(\epsilon_j^* \mid \frac{m_{\gamma_j}}{2}, \frac{v_{\gamma_j}^2}{4}\right),
 \end{aligned}$$

where d_j is the sign of ϵ_j .

Computation, continued

The joint distribution $(\epsilon_j^*, \epsilon_{j,2})$ conditional on the latent mixture element γ_j becomes

$$\begin{aligned}
 p(\epsilon_j^*, \epsilon_{j,2} | \gamma_j) &= p(\epsilon_{j,2} | \epsilon_j^*, \gamma_j) p(\epsilon_j^* | \gamma_j) \\
 &= p(\epsilon_{j,2} | \underbrace{d_j \exp(\epsilon_j^*)}_{\epsilon_j}, \gamma_j) p(\epsilon_j^* | \gamma_j) \\
 &= N\left(\epsilon_{j,2} \mid d_j \exp(\epsilon_j^*), (1 - \rho^2)\right) N\left(\epsilon_j^* \mid \frac{m_{\gamma_j}}{2}, \frac{v_{\gamma_j}^2}{4}\right),
 \end{aligned}$$

where d_j is the sign of ϵ_j . The joint distribution becomes a linear combination of independent Normal random variables when we replace the nonlinear $\exp(\epsilon_j^*)$ term with the first-order approximation

$$\exp(\epsilon_j^*) | \gamma_j \approx \exp(m_{\gamma_j}/2) (a_{\gamma_j} + b_{\gamma_j} (2\epsilon_j^* - m_{\gamma_j}))$$

Blocked Gibbs sampler with predominantly Metropolis-Hastings and Forward Filtering Backward Sampling steps

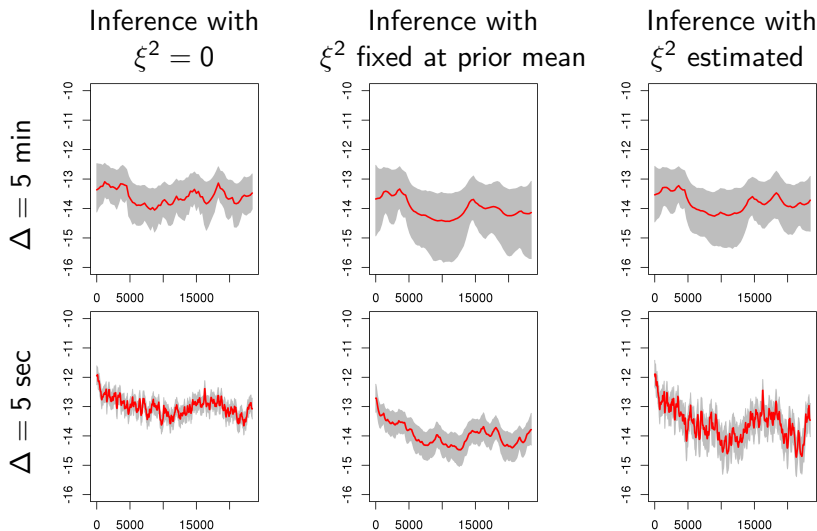
- ① sample observational parameters (MH)
- ② sample latent prices (FFBS)
- ③ sample volatility parameters (MH)
- ④ sample latent mixture indicators (Discrete)
- ⑤ sample volatility paths (FFBS)
- ⑥ sample jump parameters (MH)
- ⑦ sample jumps (Conjugate prior)

Estimating integrated volatility

- $\zeta \sim N(0, \xi^2)$
- We consider 300 simulated data sets over 1 trading day
- Coverage : percent of 95% probability/confidence intervals covering the true data-generating integrated volatility $\int \hat{\sigma}_t^2 dt$:

	Sampling period			
	60 sec	30 sec	15 sec	5 sec
Inference with $\xi^2 = 0$	72	28	3	0
Inference with ξ^2 fixed at prior mean	79	57	23	0
Inference with ξ^2 estimated	91	92	96	97
Inference with kernel-based estimator	51	48	59	76

Results for real data example



Log-volatility paths for the AAPL 03/06/2014 data.

Bivariate Open, Close, High, Low Prices

- Work has been done in the univariate case
- No equivalent results in the bivariate setting
- There exists only a single estimator in the literature which used bivariate OCHL to estimate asset correlation [9] (Rogers estimator)

Model formulation

We consider a two-dimensional correlated Brownian motion:

$$X(t) = x_0 + \mu_x t + \sigma_x W_x(t)$$

$$Y(t) = y_0 + \mu_y t + \sigma_y W_y(t)$$

where $W_x(t)$ and $W_y(t)$ are correlated standard Brownian motions with $\text{Cov}(W_x(t), W_y(t)) = \rho t$.

- We seek the 6-dimensional joint probability density function for the pair $(X(t), Y(t))$ and the random variables $M_X(t) = \max_{0 \leq s \leq t} X(s)$, $m_X(t) = \min_{0 \leq s \leq t} X(s)$, $M_Y(t) = \max_{0 \leq s \leq t} Y(s)$, $m_Y(t) = \min_{0 \leq s \leq t} Y(s)$.

Transition density

To calculate the joint density

$$p(X(t) = x, Y(t) = y, \\ m_X(t) = a_x, M_X(t) = b_x, m_Y(t) = a_y, M_Y(t) = b_y),$$

we first study cumulative-like distribution

$$p(X(t) = x, Y(t) = y, \\ m_X(t) \geq a_x, M_X(t) \leq b_x, m_Y(t) \geq a_y, M_Y(t) \leq b_y),$$

which is governed by a Fokker-Planck equation with absorbing boundaries:

Governing Fokker-Planck equation

Abbreviating

$$q(x, y, t) := p(X(t) = x, Y(t) = y, \\ m_X(t) \geq a_x, M_X(t) \leq b_x, m_Y(t) \geq a_y, M_Y(t) \leq b_y),$$

the FP equation is

$$\frac{\partial}{\partial t'} q(x, y, t') = -\mu_x \frac{\partial}{\partial x} q(x, y, t') - \mu_y \frac{\partial}{\partial y} q(x, y, t') + \\ \frac{1}{2} \sigma_x^2 \frac{\partial^2}{\partial x^2} q(x, y, t') + \rho \sigma_x \sigma_y \frac{\partial^2}{\partial x \partial y} q(x, y, t') + \frac{1}{2} \sigma_y^2 \frac{\partial^2}{\partial y^2} q(x, y, t'),$$

$$q(a_x, y, t') = q(b_x, y, t') = q(x, a_y, t') = q(x, b_y, t') = 0, \quad 0 < t' \leq t, \\ q(x, y, 0) = \delta(x_0 - x) \delta(y_0 - y)$$

Normalized problem

We can introduce a series of transformations to solve the equivalent initial-boundary problem

$$\frac{\partial}{\partial \tilde{t}} q(\tilde{x}, \tilde{y}, \tilde{t}) = \tilde{\mathcal{L}} q(\tilde{x}, \tilde{y}, \tilde{t}), \quad (\tilde{x}, \tilde{y}) \in \tilde{\Omega}$$

$$\tilde{\mathcal{L}} = \frac{1}{2} \frac{\partial^2}{\partial \tilde{x}^2} + \rho \sigma_{\tilde{y}} \frac{\partial^2}{\partial \tilde{x} \partial \tilde{y}} + \frac{1}{2} \sigma_{\tilde{y}}^2 \frac{\partial^2}{\partial \tilde{y}^2}, \quad \tilde{\Omega} := (0, 1) \times (0, 1)$$

where $0 < \sigma_{\tilde{y}}^2 \leq 1$.

- This new parameterization makes the computational domain invariant to data/parameter combinations.

Joint density is given by the fourth derivative with respect to boundaries

Denoting the joint density of the diffusion process and the attained extrema over the interval t as $f(x, y, a_x, b_x, a_y, b_y)$,

$$\frac{\partial^4}{\partial a_x \partial b_x \partial a_y \partial b_y} q(x, y, t) = f(x, y, a_x, b_x, a_y, b_y).$$

- The finite difference approximation of $\partial^4 / \partial a_x \partial b_x \partial a_y \partial b_y$ introduces a fundamental limitation due to finite precision round-off errors when the analytic solution $q(x, y, t)$ is not available

$$\mathcal{O}\left(\frac{\varepsilon_{mach}}{\varepsilon^4}\right) \rightarrow \infty \quad \text{as } \varepsilon \rightarrow 0,$$

where ε is the step size of the numerical differentiation with respect to the boundaries.

Approach 1: Finite Difference approximation to FP equation

Requires the solution to a system of ODEs:

$$\dot{c}(\tilde{t}) = Bc(\tilde{t}) \quad \Rightarrow c(\tilde{t}) = \exp(B\tilde{t}) c(0),$$

where the sparse system matrix B can be computed once and stored for a regular grid over $\tilde{\Omega}$ with finite step size h .

$$B = \frac{1}{2} \frac{1}{h^2} B_{\tilde{x},\tilde{x}} + \rho \sigma_{\tilde{y}} \frac{1}{4h^2} B_{\tilde{x},\tilde{y}} + \frac{1}{2} \sigma_{\tilde{y}}^2 \frac{1}{h^2} B_{\tilde{y},\tilde{y}}.$$

Approach 1: Finite Difference approximation to FP equation

Letting $k := 1/h$ and b denote the boundary parameters,

- FD approximation is fast but also limited by irregular truncation errors

$$\begin{aligned}
 q^{(k)}(\tilde{x}, \tilde{y}, \tilde{t}|b) - q(\tilde{x}, \tilde{y}, \tilde{t}|b) &= \underbrace{\left(\frac{1}{k}\right)^\alpha F_{reg}(b)}_{\text{smooth in } b, \alpha > 0} \\
 + \underbrace{\left(\frac{1}{k}\right)^\beta F_{irreg}(b)}_{\text{continuous but not differentiable w.r.t } b, \beta > 0} &+ \underbrace{\varepsilon_{mach} F_{round}(b)}_{\text{behaves as a R.V}}
 \end{aligned}$$

- With linear interpolation which arises when function arguments are not on grid points with respect to boundary differentiation

$$\mathcal{O}(F_{irreg}(b)) = k^\beta / \varepsilon \quad \Rightarrow \quad \left(\frac{1}{k}\right)^\beta F_{irreg}(b) \rightarrow \infty \text{ as } \varepsilon \rightarrow 0.$$

Approach 2: Trigonometric expansion of the differential operator

$$q(\tilde{x}, \tilde{y}, \tilde{t}) \approx \sum_{\nu} h_{\nu} \phi_{\nu}(\tilde{x}, \tilde{y}) e^{-\lambda_{\nu} \tilde{t}}.$$

- When $\rho = 0$,

$$q(\tilde{x}, \tilde{y}, \tilde{t}) \approx \sum_{l=1}^L \sum_{m=1}^M c_{l,m} \sin(2\pi l \tilde{x}) \sin(2\pi m \tilde{y})$$

- the system matrix is diagonal
- no new error is introduced in the time evolution
- truncated solution converges fast

Approach 2: Trigonometric expansion of the differential operator

- When $\rho \neq 0$, mixing term in the FP equation produces $\cos(2\pi l \tilde{x}) \cos(2\pi k \tilde{y})$ terms in the above expansion such that

$$\phi_\nu(\tilde{x}, \tilde{y}) = \sum_{l=1}^L \sum_{m=1}^M c_{l,m,\nu} \sin(2\pi l \tilde{x}) \sin(2\pi m \tilde{y}) := \Psi(\tilde{x}, \tilde{y})^T c_\nu,$$

- system matrix is dense and convergence is slow,
- new error is introduced in the time evolution

The Galerkin approach: weak solution to the PDE

We propose a solution $q^{(k)}(\tilde{x}, \tilde{y}, \tilde{t})$ of similar form

$$q^{(k)}(\tilde{x}, \tilde{y}, \tilde{t}) = \sum_{i=0}^k c_i(\tilde{t}) \psi_i(\tilde{x}, \tilde{y}),$$

where the basis functions $\psi_i(\tilde{x}, \tilde{y})$ satisfy the boundary conditions on $\tilde{\Omega}$ and they need not be eigenfunctions but capture some essential features of the solution. We also require that all first- and second-order derivatives of $\psi_i(\tilde{x}, \tilde{y})$ are in $L_2(\tilde{\Omega})$.

$$\begin{aligned} \frac{\partial}{\partial \tilde{t}} q^{(k)}(\tilde{x}, \tilde{y}, \tilde{t}) - \tilde{\mathcal{L}} q^{(k)}(\tilde{x}, \tilde{y}, \tilde{t}) &:= R_e(k), \\ q(\tilde{x}, \tilde{y}, 0) - q^{(k)}(\tilde{x}, \tilde{y}, 0) &:= R_0(k). \end{aligned}$$

The Galerkin approach: weak solution to the PDE

The *orthogonality* condition of the Galerkin procedure:

$$\int_{\Omega} R_e(k) \psi_i(\tilde{x}, \tilde{y}) d\tilde{x} d\tilde{y} = 0, \quad \int_{\Omega} R_0(k) \psi_i(\tilde{x}, \tilde{y}) d\tilde{x} d\tilde{y} = 0, \quad i = 0, \dots, k,$$

which is equivalent to the weak formulation of the heat problem

$$\begin{aligned} \langle \partial_t q^{(k)}(\tilde{x}, \tilde{y}, \tilde{t}), \psi_i \rangle &= \langle \tilde{\mathcal{L}} q^{(k)}(\tilde{x}, \tilde{y}, \tilde{t}), \psi_i \rangle, \\ \langle q^{(k)}(\tilde{x}, \tilde{y}, 0), \psi_i \rangle &= \langle q(\tilde{x}, \tilde{y}, 0), \psi_i \rangle, \end{aligned}$$

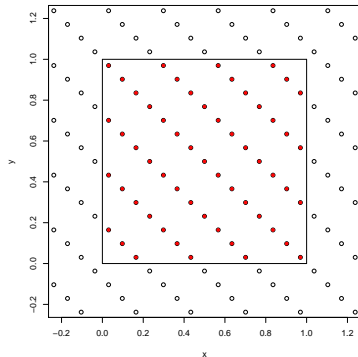
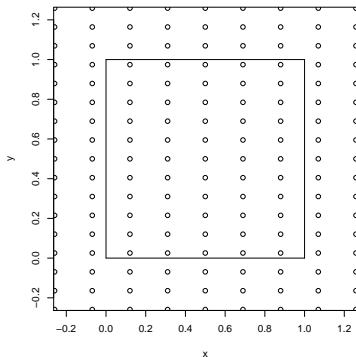
where $\langle \cdot, \cdot \rangle$ is the usual inner product in $L_2(\tilde{\Omega})$.

Basis element choice

$$\begin{aligned} \psi_i(\tilde{x}, \tilde{y}) = & \frac{1}{2\pi\tilde{\sigma}^2\sqrt{1-\tilde{\rho}^2}} \\ & \times \exp \left\{ -\frac{((\tilde{x} - \tilde{x}_i)^2 - 2\tilde{\rho}(\tilde{x} - \tilde{x}_i)(\tilde{y} - \tilde{y}_i) + (\tilde{y} - \tilde{y}_i)^2)}{2(1-\tilde{\rho}^2)\tilde{\sigma}^2} \right\} \\ & \times \tilde{x}(1-\tilde{x})\tilde{y}(1-\tilde{y}) \end{aligned}$$

for some parameters $(\tilde{\rho}, \tilde{\sigma})$ and a collection of nodes $\{(\tilde{x}_i, \tilde{y}_i)\}_{i=0}^k$ which form a grid over $\tilde{\Omega}$.

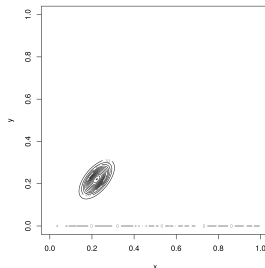
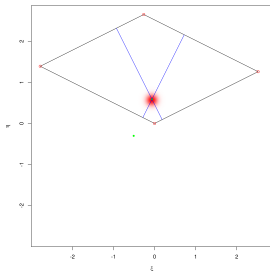
- Align basis elements along the principal axis of the differential operator,
- Spacing between nodes is l times the bandwidth $\tilde{\sigma}$ in each principal direction
- This scheme is aimed at better resolving the correlation in the solution



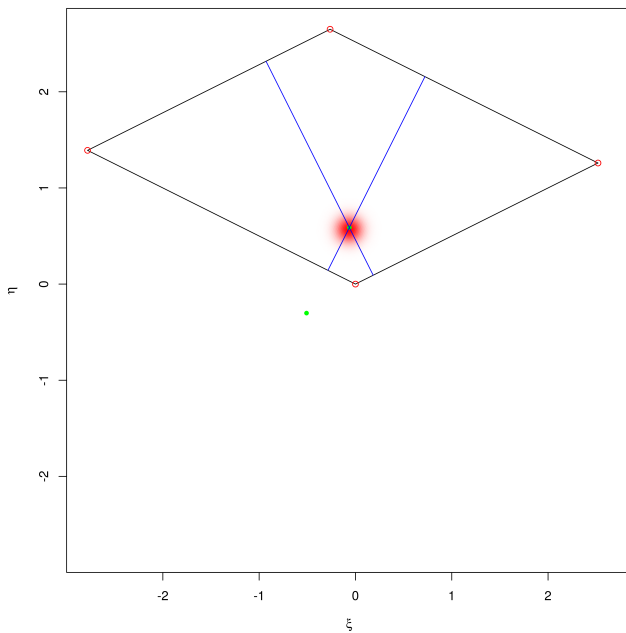
- A sample grid design for $l = 1$, $\tilde{\sigma} = 0.3$ and $\tilde{\rho} = 0.6$.
- **Left:** Grid along coordinates of observed prices within the computational domain (solid black square).
- **Right:** Grid along the major axes of the basis element kernels. The set of final node points $\{(x_i, y_i)\}_{i=0}^k$ is contained within the computational domain and is denoted by the red solid points.

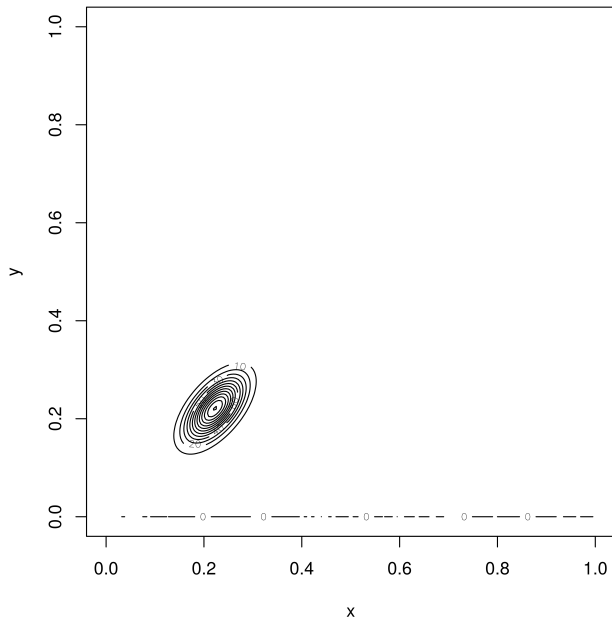
Small-time solution to the Fokker-Planck equation

- A semi-analytic solution for a small time \tilde{t}_ϵ allows us to project a smooth function onto our basis family instead of a δ -function
- This reduces the numerical error of the Galerkin solver



- **Left:** Computational domain under the transformation where diffusion parameters are both unity and correlation is zero.
- **Right:** Contour of small-time solution in original coordinate system.





Consistency Results

Denoting the approximate density due to the Galerkin solution as

$$f^{(k)}(x, y, a_x, b_x, a_y, b_y) = \frac{\partial^4 q^{(k)}(x, y, t)}{\partial a_x \partial b_x \partial a_y \partial b_y},$$

we show that the two integrals below converge in $L_2(\Omega)$:

Lemma (1)

$$\lim_{k \rightarrow \infty} \int_{a_x}^{b_x} \int_{a_y}^{b_y} \left(f^{(k)}(x, y, a_x, b_x, a_y, b_y) - f(x, y, a_x, b_x, a_y, b_y) \right)^2 dx dy = 0$$

Denoting the random variable

$$Z := (X(t), Y(t), m_X(t), M_X(t), m_Y(t), M_Y(t)),$$

we consider the distribution function $F(Z \leq z)$ generated by the true solution to the Fokker-Planck equation as well as the distribution $F^{(k)}(Z \leq z)$ generated by the Galerkin solution.

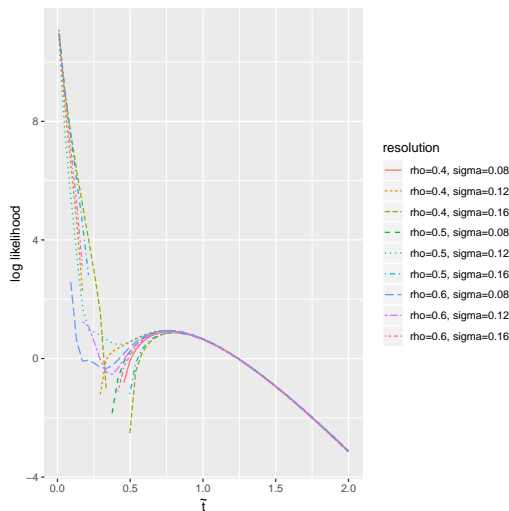
Lemma (Convergence in distribution)

For any $z \in Z$, $\lim_{k \rightarrow \infty} F^{(k)}(Z \leq z) = F(Z \leq z)$.

Assuming that the maximum likelihood estimate (MLE) (under $F^{(k)}$ for sufficiently large k) is continuous with respect to the data, we can use the convergence result above and Chebyshev's inequality to show that the MLE under $F^{(k)}$ is consistent as the number of basis elements k and number of data points n go to infinity.

Results and solution behavior for small \tilde{t}

Here we consider the parameters defining the basis family $(\tilde{\rho}, \tilde{\sigma})$



Simulation Study

We consider 50 simulated data sets with

$$\sigma_x = 1, \quad \sigma_y = 1, \quad \mu_x = 0, \quad \mu_y = 0$$

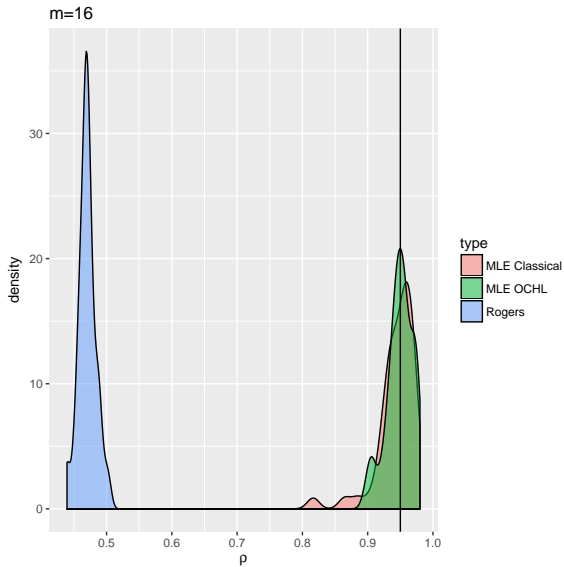
with increasing sample size over several values of ρ . For each case, we consider the ratio of mean-square error of the Galerkin solver compared to that generated by Gaussian likelihoods ignoring the boundaries:

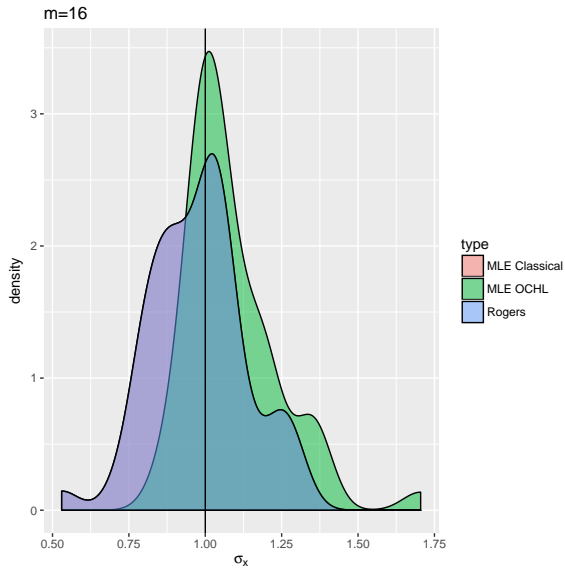
	$\rho = 0.95$			$\rho = 0.60$		
	$m = 4$	$m = 8$	$m = 16$	$m = 4$	$m = 8$	$m = 16$
$\hat{\sigma}_x$	0.475	1.238	1.304	0.203	0.127	0.232
$\hat{\sigma}_y$	0.593	1.040	1.088	0.111	0.120	0.260
$\hat{\rho}$	0.287	0.910	0.445	0.315	0.283	0.463

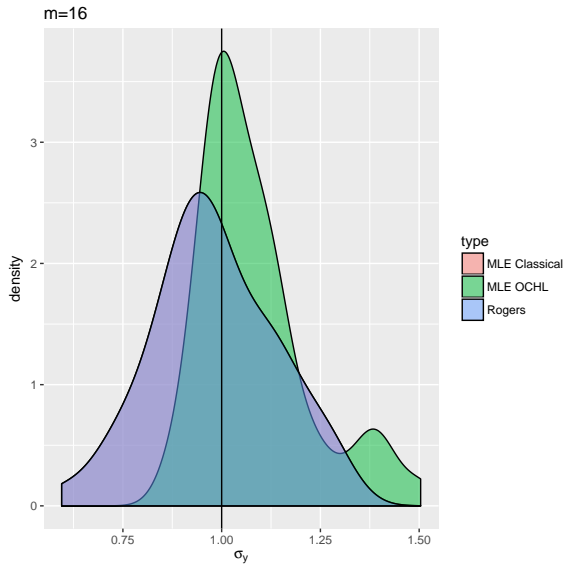
	$\rho = 0.0$		
	$m = 4$	$m = 8$	$m = 16$
$\hat{\sigma}_x$	0.137	0.243	0.167
$\hat{\sigma}_y$	0.189	0.171	0.107

Simulation study

- For data generated with $\rho = 0.95$, kernel-density approximations of the repeated-sampling densities of the MLEs are shown. Samples are obtained from the Galerkin likelihood (green) and the classical Gaussian likelihood (red) The data-generating parameters are denoted with the vertical solid line.
- The Rogers estimator (blue) is the only existing correlation estimator for bivariate OCHL data.

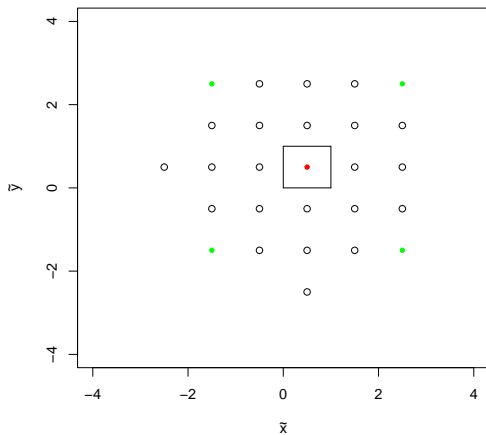






Open, Close, High, Low Prices: Small Time Solution and Gap-Fill Solution

Method of images and analytic differentiability



Symmetry constraint in constructing approximate solutions

A condition weaker than uniqueness which nonetheless restricts the solution space is the symmetry obeyed by the problem. We consider the transformation

$$\begin{aligned}x^{new} &= (a_x + b_x) - x^{old}, \\ y^{new} &= (a_y + b_y) - y^{old}.\end{aligned}$$

Performing the set of reflections

$$\{2, 4, 1, 3\} \cup \{2, 4, 3, 1\} \cup \{4, 2, 1, 3\} \cup \{4, 2, 3, 1\}.$$

produces a sum of images where only four elements are differentiable with respect to all four boundaries:

$$\frac{\partial^4 p_\epsilon(\tilde{x}, \tilde{y}, \tilde{t})}{\partial a_x \partial b_x \partial a_y \partial b_y} = \sum_{j'=1}^4 \frac{\partial^4 G(\tilde{x}, \tilde{y}, \tilde{t} | \tilde{x}_{(j')}, \tilde{y}_{(j')})}{\partial a_x \partial b_x \partial a_y \partial b_y}.$$

Calculation of the joint density

We can express the derivatives:

$$\begin{aligned}
 & \frac{\partial^4}{\partial a_x \partial b_x \partial a_y \partial b_y} G(\tilde{x}, \tilde{y}, \tilde{t} | \tilde{x}_{(j)}, \tilde{y}_{(j)}) = \\
 & G \cdot \mathcal{C}^4 \cdot \left(\partial_{a_x} \partial_{b_x} \partial_{a_y} \partial_{b_y} \right) \mathcal{P} \\
 & + G \cdot \mathcal{C}^3 \cdot \left(\partial_{a_x}^2 \partial_{b_x} \partial_{a_y} \partial_{b_y} + \partial_{a_x}^2 \partial_{a_y} \partial_{b_x} \partial_{b_y} + \partial_{a_x}^2 \partial_{b_x} \partial_{b_y} \partial_{a_y} + \right. \\
 & \quad \left. + \partial_{b_x}^2 \partial_{a_y} \partial_{a_x} \partial_{b_y} + \partial_{b_x}^2 \partial_{b_y} \partial_{a_x} \partial_{a_y} \partial_{a_y}^2 \partial_{b_y} \partial_{a_x} \partial_{b_x} \right) \mathcal{P} \\
 & + G \cdot \mathcal{C}^2 \cdot \left(\partial_{a_x}^3 \partial_{b_x} \partial_{a_y} \partial_{b_y} + \partial_{a_x}^2 \partial_{b_x} \partial_{a_y}^2 \partial_{b_y} + \partial_{a_x}^3 \partial_{b_x} \partial_{b_y} \partial_{a_y} + \right. \\
 & \quad \left. + \partial_{a_x}^3 \partial_{a_y} \partial_{b_x} \partial_{b_y} + \partial_{a_x}^2 \partial_{a_y} \partial_{b_y}^2 \partial_{b_x} + \partial_{b_x}^3 \partial_{a_y} \partial_{b_y} \partial_{a_x} + \partial_{b_x}^2 \partial_{a_y} \partial_{a_x} \partial_{b_y} \right) \mathcal{P} \\
 & + G \cdot \mathcal{C} \cdot \partial_{a_x}^4 \partial_{b_x} \partial_{a_y} \partial_{b_y} \mathcal{P}
 \end{aligned}$$

Analytic Derivative and higher order terms

Thinking of \tilde{t} as variable allows us to further simplify the expression. Since $\mathcal{C} = \mathcal{O}(1/\tilde{t})$, all three terms $G \cdot \mathcal{C}^3$, $G \cdot \mathcal{C}^2$, and $G \cdot \mathcal{C}$ are $o\left(G \cdot \mathcal{C}^4 \cdot \left(\partial_{a_x} \partial_{b_x} \partial_{a_y} \partial_{b_y}\right) \mathcal{P}\right)$, so that the $G \cdot \mathcal{C}^4$ order term in the derivative dominates the others for sufficiently small \tilde{t} .

$$\frac{\partial^4 p_\epsilon(\tilde{x}, \tilde{y}, \tilde{t})}{\partial a_x \partial b_x \partial a_y \partial b_y} \approx \sum_{j'=1}^4 G(\tilde{x}, \tilde{y}, \tilde{t} | \tilde{x}_{(j')}, \tilde{y}_{(j')}) \cdot \mathcal{C}^4 \cdot \left(\partial_{a_x} \partial_{b_x} \partial_{a_y} \partial_{b_y}\right) \mathcal{P}_{j'}.$$

Basis Expansion of Galerkin likelihood: RHS

The likelihood computed with the Galerkin solution as a function of \tilde{t} is of the form

$$\frac{\partial^4 q_{\text{Galerkin}}(\tilde{x}, \tilde{y}, \tilde{t})}{\partial a_x \partial b_x \partial a_y \partial b_y} = \sum_{k=1}^K e^{-\lambda_k \tilde{t}} p_k^{(4)}(\tilde{t}),$$

where $p_i^{(4)}(\tilde{t})$ is a fourth-order polynomial. This proceeds from the Galerkin solution being dependent on \tilde{t} only through the exponential term: the eigenfunctions of the solution are by design solely functions of (a_x, b_x, a_y, b_y) and (\tilde{x}, \tilde{y})

- This suggests a leading-order approximation that is fitted with the Galerkin solver via least squares:

$$f_{\text{Galerkin}}(\tilde{t}) = \tilde{t}^4 \left(\omega_1 e^{-\lambda_1 \tilde{t}} + \omega_2 e^{-\lambda_2 \tilde{t}} \right),$$

$$\log f_{\text{Galerkin}}(\tilde{t}) = \log(\omega_1) + 4 \log(\tilde{t}) - \lambda_1 \tilde{t} + \log \left(1 + \omega_2 / \omega_1 e^{-(\lambda_2 - \lambda_1) \tilde{t}} \right)$$

Leading order for small-time solution: LHS

The summand in the small-time likelihood with the greatest $\beta_{j'}$ contributes the most to the truncated small-time solution in the $\tilde{t} \leq 1$ region where the matched solution will be applied. Indexing j' such that $\beta_1 \geq \beta_2 \geq \beta_3 \geq \beta_4$, the small-time log-likelihood is

$$\begin{aligned} \log \left(\frac{\partial^4 p_\epsilon(\tilde{x}, \tilde{y}, \tilde{t})}{\partial a_x \partial b_x \partial a_y \partial b_y} \right) &\approx \log(K) - 4.5 \log(\tilde{t}) + \log(c_1) - \frac{\beta_1}{\tilde{t}} \\ &\quad + \log \left(1 + \sum_{j \neq 1} \frac{c_j}{c_1} \exp \left(-\frac{(\beta_j - \beta_1)}{\tilde{t}} \right) \right) \\ &\approx \log(K) - 4.5 \log(\tilde{t}) + \log(c_1) - \frac{\beta_1}{\tilde{t}} + \log(1 + \epsilon(\tilde{t})). \end{aligned}$$

- The proposed matched solution is of the form

$$\log f_{\text{gap}}(\tilde{t}) = \log(\omega(\tilde{t})) - \gamma(\tilde{t}) \log(\tilde{t}) - \frac{\beta(\tilde{t})}{\tilde{t}}$$

Gap-fill solution

$$\log f_{\text{Galerkin}}(\tilde{t}) \approx \log(\omega_1) + 4 \log(\tilde{t}) - \lambda_1 \tilde{t} + \log \left(1 + \omega_2/\omega_1 e^{-(\lambda_2 - \lambda_1)\tilde{t}} \right)$$

$$\log f_{\text{small-time}}(\tilde{t}) \approx \log(K) - 4.5 \log(\tilde{t}) - \frac{\beta_1}{\tilde{t}} + \log(c_1) + \log(1 + \epsilon(\tilde{t})),$$

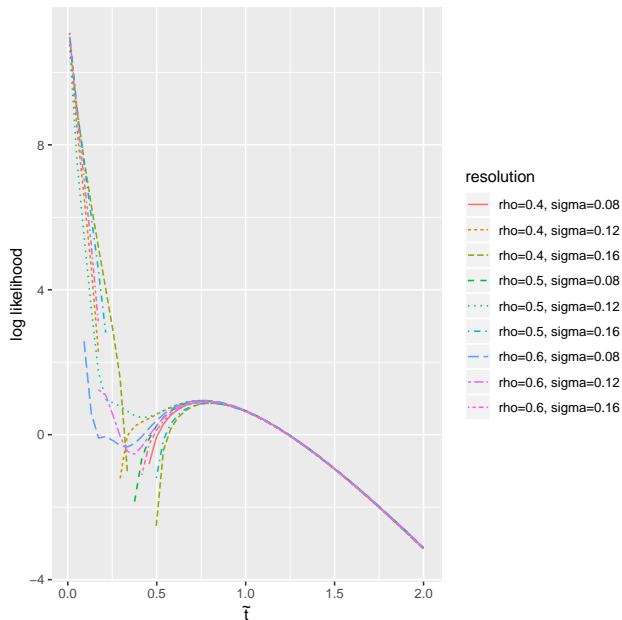
$$\log f_{\text{gap}}(\tilde{t}) = \log(\omega(\tilde{t})) - \gamma(\tilde{t}) \log(\tilde{t}) - \frac{\beta(\tilde{t})}{\tilde{t}}$$

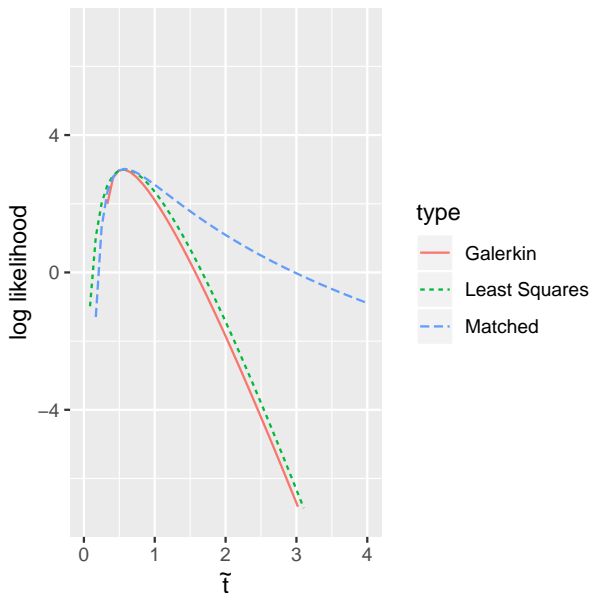
- At \tilde{t}^* , the left-hand side of the matching condition, the values for these parameters are defined such that they match the small-time solution

$$\omega(\tilde{t}^*) = K, \quad \gamma(\tilde{t}^*) = 4.5, \quad \beta(\tilde{t}^*) = \beta_1.$$

- At \tilde{t}_m , the right-hand side of the matching condition and the maximum of the LS solution, $\omega(\tilde{t})$, $\gamma(\tilde{t})$, and $\beta(\tilde{t})$ are chosen to match the value, first, and second derivatives of the logarithmic form of the Galerkin approximate log-likelihood.

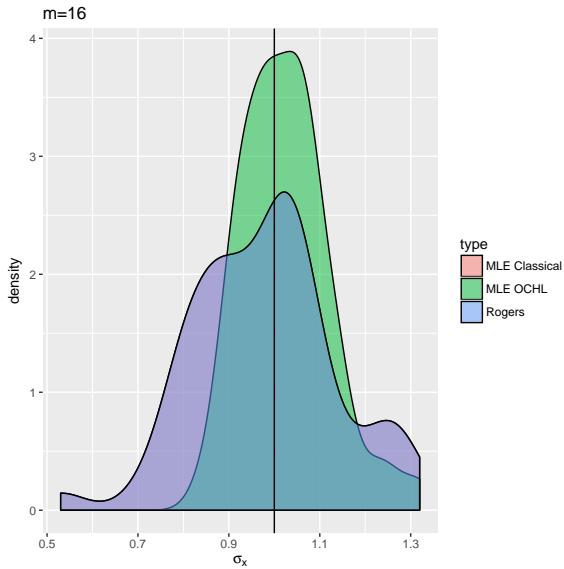
$$\begin{aligned}\omega(\tilde{t}) &= \omega(\tilde{t}^*)e^{-k(\tilde{t}-\tilde{t}^*)} + \omega(\tilde{t}_m) \left(1 - e^{-k(\tilde{t}-\tilde{t}^*)}\right), \\ \gamma(\tilde{t}) &= \gamma(\tilde{t}^*)e^{-k(\tilde{t}-\tilde{t}^*)} + \gamma(\tilde{t}_m) \left(1 - e^{-k(\tilde{t}-\tilde{t}^*)}\right), \\ \beta(\tilde{t}) &= \beta(\tilde{t}^*)e^{-k(\tilde{t}-\tilde{t}^*)} + \beta(\tilde{t}_m) \left(1 - e^{-k(\tilde{t}-\tilde{t}^*)}\right).\end{aligned}$$





Repeated Simulation Results

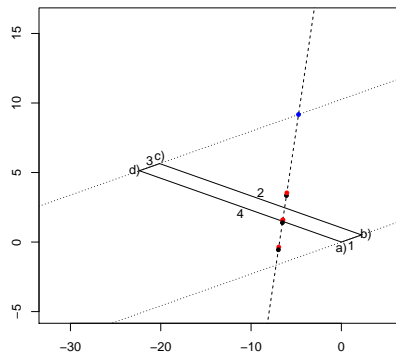
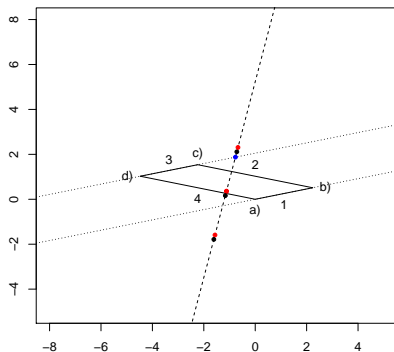
	$\rho = 0.95$		
	$m = 4$	$m = 8$	$m = 16$
$\hat{\sigma}_x$	0.124	0.429	0.318
$\hat{\sigma}_y$	0.310	0.147	0.365
$\hat{\rho}$	0.250	0.753	0.699



- Multivariate High-Frequency models
- Gaussian Process Emulator for interpolating the OCHL likelihood
- Particle Filter

Thank you

Illustration of proof for existence



Prior elicitation example I

Prior for $\theta_i(\Delta)$: The discrete-time autocorrelation coefficient $\theta(\Delta)$ of the volatility process is bounded above by 1 and below by 0 such that $\log(\sigma_j)$ is bounded as $j \rightarrow \infty$. Hence, we employ a truncated normal prior for $\theta(\Delta)$,

$$p(\theta(\Delta)) \propto N\left(a_\theta(\Delta), b_\theta^2(\Delta)\right) \mathbb{1}_{(\theta(\Delta) \in [0,1])},$$

Prior elicitation example II

which leads again to a tractable computational algorithm. Note that because of the truncation,

$$E[\theta(\Delta)] = a_\theta(\Delta) + \frac{\phi\left(-\frac{a_\theta(\Delta)}{b_\theta(\Delta)}\right) - \phi\left(\frac{1-a_\theta(\Delta)}{b_\theta(\Delta)}\right)}{\Phi\left(\frac{1-a_\theta(\Delta)}{b_\theta(\Delta)}\right) - \Phi\left(-\frac{a_\theta(\Delta)}{b_\theta(\Delta)}\right)} b_\theta(\Delta) \quad (1)$$

$$\begin{aligned} \text{Var}[\theta(\Delta)] = b_\theta^2(\Delta) & \left[1 + \frac{-\frac{a_\theta(\Delta)}{b_\theta(\Delta)}\phi\left(-\frac{a_\theta(\Delta)}{b_\theta(\Delta)}\right) - \frac{1-a_\theta(\Delta)}{b_\theta(\Delta)}\phi\left(\frac{1-a_\theta(\Delta)}{b_\theta(\Delta)}\right)}{\Phi\left(\frac{1-a_\theta(\Delta)}{b_\theta(\Delta)}\right) - \Phi\left(-\frac{a_\theta(\Delta)}{b_\theta(\Delta)}\right)} \right. \\ & \left. + \left\{ \frac{\phi\left(-\frac{a_\theta(\Delta)}{b_\theta(\Delta)}\right) - \phi\left(\frac{1-a_\theta(\Delta)}{b_\theta(\Delta)}\right)}{\Phi\left(\frac{1-a_\theta(\Delta)}{b_\theta(\Delta)}\right) - \Phi\left(-\frac{a_\theta(\Delta)}{b_\theta(\Delta)}\right)} \right\}^2 \right] \quad (2) \end{aligned}$$

where $\phi(\cdot)$ and $\Phi(\cdot)$ denote the density and the cumulative distribution functions of the standard normal distribution. Now, given the prior mean

Prior elicitation example III

$\hat{a}_{\hat{\theta}}$ and variance $\hat{b}_{\hat{\theta}}^2$ for $\hat{\theta}$, we choose the values of $a_{\theta}(\Delta)$ and $b_{\theta}(\Delta)$ so that the mean and variance of $\theta(\Delta)$ above are approximately equal to the mean and variance of $\exp\{-\hat{\theta}\Delta\}$. To simplify calculation of the moments of $\exp\{-\hat{\theta}\Delta\}$ we use a second-order Taylor expansion of $\exp\{-\hat{\theta}\Delta\}$ to approximate the first two moments of $\theta(\Delta)$ in terms of $\hat{a}_{\hat{\theta}}$ and $\hat{b}_{\hat{\theta}}^2$, an approach known as the Delta-Method (e.g., see [5]):

$$E\left[\exp\{-\hat{\theta}\Delta\}\right] \approx \exp(-\hat{a}_{\hat{\theta}}\Delta) \left(1 + \frac{1}{2}\hat{b}_{\hat{\theta}}^2\Delta^2\right), \quad (3)$$

$$E\left[\exp\{-2\hat{\theta}\Delta\}\right] \approx \exp(-2\hat{a}_{\hat{\theta}}\Delta) \left(1 + 2\hat{b}_{\hat{\theta}}^2\Delta^2\right). \quad (4)$$

Using (1), (2), (3), and (4), and by setting

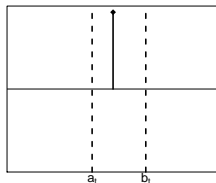
$$E[\theta(\Delta)] = E\left[\exp\{-\hat{\theta}\Delta\}\right], \quad \text{Var}[\theta(\Delta)] = \text{Var}\left[\exp\{-\hat{\theta}\Delta\}\right],$$

Prior elicitation example IV

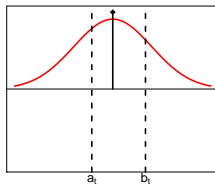
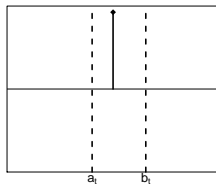
we obtain a system of two equations with two unknowns that can be solved numerically to find the values of $a_{\theta}(\Delta)$ and $b_{\theta}^2(\Delta)$ in terms of $\hat{a}_{\hat{\theta}}$, $\hat{b}_{\hat{\theta}}^2$, and Δ .

To elicit $\hat{a}_{\hat{\theta}}$ and $\hat{b}_{\hat{\theta}}^2$, recall that $\hat{\theta}$ is the inverse of the time scale of inertia for $\log(\hat{\sigma}_t)$ in the continuous-time formulation, which can be thought of as the characteristic time length, or unit of time, over which the process for the diffusion of $\log(\hat{\sigma}_t)$ “forgets” about an endogenous shock. The two hyper-parameters can be chosen so that the prior probability mass for $\hat{\theta}$ permits a reasonable range for the timescale of inertia.

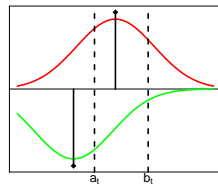
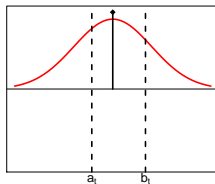
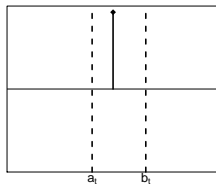
Solution by Method of Images: single dimension



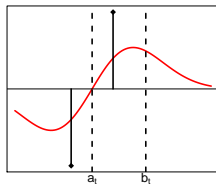
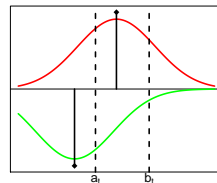
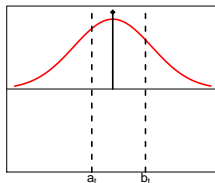
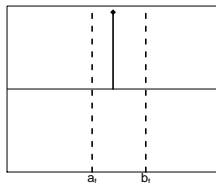
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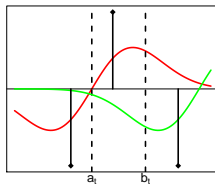
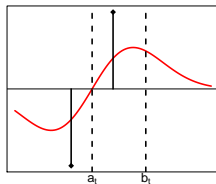
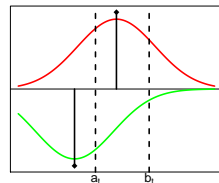
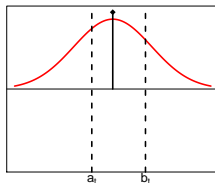
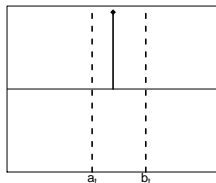
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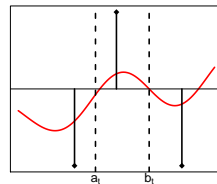
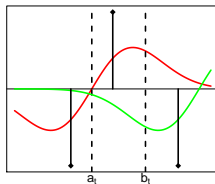
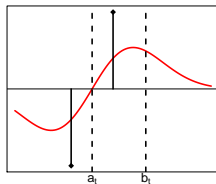
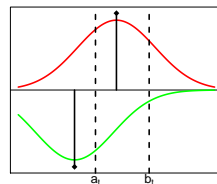
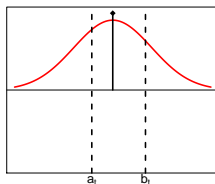
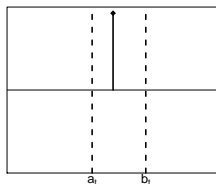
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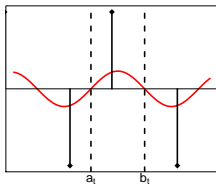
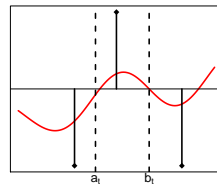
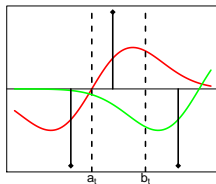
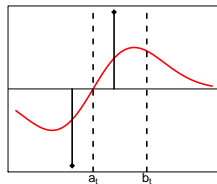
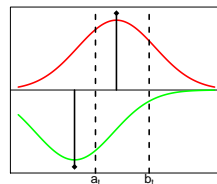
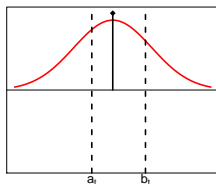
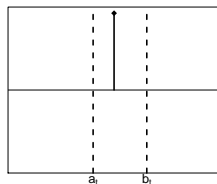
Solution by Method of Images: single dimension



Solution by Method of Images: single dimension

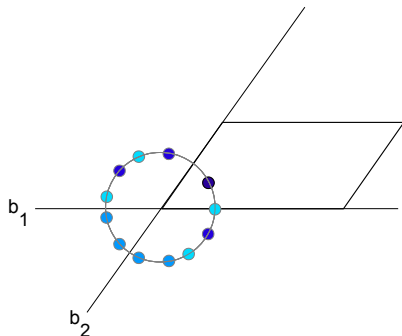
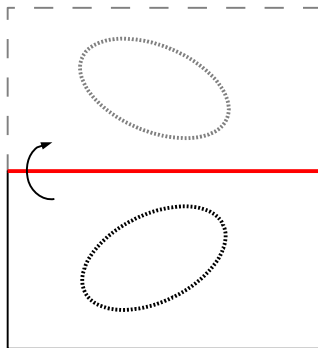


Solution by Method of Images: single dimension



Solution by Method of Images: bivariate case

Method of Images does not work in general.



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