Characteristic function estimation of Ornstein-Uhlenbeck-based stochastic volatility models

Emanuele Taufer, Nikolai Leonenko, Marco Bee
Characteristic function estimation of Ornstein-Uhlenbeck-based stochastic volatility models

Emanuele Taufer, Nikolai Leonenko, Marco Bee
DISA Working Papers
The series of DISA Working Papers is published by the Department of Computer and Management Sciences (Dipartimento di Informatica e Studi Aziendali DISA) of the University of Trento, Italy.

Editor
Ricardo Alberto MARQUES PEREIRA  ricalb.marper@unitn.it

Managing editor
Roberto GABRIELE  roberto.gabriele@unitn.it

Associate editors
Flavio BAZZANA  flavio.bazzana@unitn.it  Finance
Michele BERTONI  michele.bertoni@unitn.it  Financial and management accounting
Pier Franco CAMUSSONE  pierfranco.camussone@unitn.it  Management information systems
Luigi COLAZZO  luigi.colazzo@unitn.it  Computer Science
Michele FEDRIZZI  michele.fedrizzi@unitn.it  Mathematics
Andrea FRANCESCONI  andrea.francesconi@unitn.it  Public Management
Loris GAIO  loris.gaio@unitn.it  Business Economics
Umberto MARTINI  umberto.martini@unitn.it  Tourism management and marketing
Pier Luigi NOVI INVERARDI  pierluigi.noviinverardi@unitn.it  Statistics
Marco ZAMARIAN  marco.zamarian@unitn.it  Organization theory

Technical officer
Paolo FURLANI  paolo.furlani@unitn.it

Guidelines for authors
Papers may be written in English or Italian but authors should provide title, abstract, and keywords in both languages. Manuscripts should be submitted (in pdf format) by the corresponding author to the appropriate Associate Editor, who will ask a member of DISA for a short written review within two weeks. The revised version of the manuscript, together with the author’s response to the reviewer, should again be sent to the Associate Editor for his consideration. Finally the Associate Editor sends all the material (original and final version, review and response, plus his own recommendation) to the Editor, who authorizes the publication and assigns it a serial number.

The Managing Editor and the Technical Officer ensure that all published papers are uploaded in the international RepEc public-action database. On the other hand, it is up to the corresponding author to make direct contact with the Departmental Secretary regarding the offprint order and the research fund which it should refer to.

Ricardo Alberto MARQUES PEREIRA
Dipartimento di Informatica e Studi Aziendali
Università degli Studi di Trento
Via Inama 5, TN 38122  Trento  ITALIA
Tel +39-0461-282147  Fax +39-0461-282124
E-mail: ricalb.marper@unitn.it
Characteristic function estimation of Ornstein-Uhlenbeck-based stochastic volatility models

Emanuele Taufer  
Faculty of Economics, University of Trento - Italy

Nikolai Leonenko  
School of Mathematics, University of Cardiff - United Kingdom

Marco Bee  
Faculty of Economics, University of Trento - Italy

Abstract
Continuous-time stochastic volatility models are becoming increasingly popular in finance because of their flexibility in accommodating most stylized facts of financial time series. However, their estimation is difficult because the likelihood function does not have a closed-form expression. In this paper we propose a characteristic function-based estimation method for non-Gaussian Ornstein-Uhlenbeck-based stochastic volatility models. After deriving explicit expressions of the characteristic functions for various cases of interest we analyze the asymptotic properties of the estimators and evaluate their performance by means of a simulation experiment. Finally, a real-data application shows that the superposition of two Ornstein-Uhlenbeck processes gives a good approximation to the dependence structure of the process.

Keywords: Ornstein-Uhlenbeck process, Lévy process, stochastic volatility, characteristic function estimation.

MSC : 62F10, 62F12, 62M05.

1 Introduction

Since the pioneering work by Black and Scholes (1973) and Merton (1973), continuous-time stochastic models are the cornerstone of modern financial engineering. In the majority of cases, the behavior of financial variables is described by means of Geometric Brownian Motion (GBM). This
solution is adopted mostly because of its mathematical tractability, but several empirical studies have shown that the returns of certain security prices do not satisfy some of the properties implied by GBM. In particular, if the continuous-time log-price were a GBM, the log-returns would be normal, but at least for moderate to high-frequency data the normality assumption is typically rejected, because the data exhibit skewness, excess kurtosis, serial correlation, jumps and time-varying volatility; see Cont (2001) for a review. Thus, a lot of research has focused on the problem of finding continuous-time stochastic models that take into account one or more of the aforementioned issues. Barndorff-Nielsen and Shephard (2001a) introduced stochastic volatility models based on non-Gaussian Ornstein-Uhlenbeck (OU for short) processes; these models seem to be particularly well-suited, due to their flexibility in accommodating many of these problems.

Consider the following asset return process with stochastic conditional volatility of the log-asset price \( S(t) \)

\[
dS(t) = \{\mu + \beta X(t)\}dt + \sqrt{X(t)}dB(t),
\]

(1)

where \( B(t) \) is a standard Brownian motion independent of \( X(t) \), \( \mu \) is a drift and \( \beta \) is a risk premium parameter and \( \{X(t), t \geq 0\} \) is a non-negative Ornstein-Uhlenbeck type process, i.e. the solution of the stochastic differential equation

\[
dx(t) = -\lambda X(t)dt + d\hat{Z}(\lambda t).
\]

(2)

Here \( \lambda > 0 \) and \( \hat{Z}(t) \) is a homogeneous Lévy process, commonly referred to as the Background Driving Lévy Process (BDLP), for which \( \mathbb{E}[\log(1 + |\hat{Z}(1)|)] < \infty \). The specific timing adopted for \( \hat{Z}(\lambda t) \) in (2) implies that the marginal distribution of \( X(t) \) remains the same no matter what the value of \( \lambda \) is. The modelling via the use of general Lévy processes, other than Brownian motion, allows to introduce specific non-Gaussian distributions for the marginal law of \( X(t) \) and has received considerable attention in recent literature in an attempt to accommodate features such as jumps,
semi-heavy tails and asymmetry. For further details on these models and the subsequent discussion
the reader is referred to Barndorff-Nielsen and Shephard (2001a, 2001b, 2003), Barndorff-Nielsen
and Leonenko (2005).

In the present context the OU process $X(t)$ is positive, with no Gaussian component, moves
entirely by jumps and decays exponentially between two jumps while $S(t)$ remains continuous.

Suppose we observe the process $S(t)$ at fixed time instants $t_0 < t_1 < \cdots < t_n$ and define
\[ S_j = S(t_j) - S(t_{j-1}), \]
then $S_j|x_j^* \sim N(\mu(t_j - t_{j-1}) + \beta x_j^*, x_j^*)$ where $X_j^* = X^*(t_j) - X^*(t_{j-1})$ and
$X^*(t_j) = \int_{t_0}^{t_j} X(s)ds$ is the integrated volatility; $x_j^*$ is termed actual volatility. The terms $S_j|x_j^*$
$j = 1, \ldots, n$ are conditionally independent and we note that the distribution of $S_j$ will be a location
scale mixture of normals. By an appropriate design of the stochastic process for $X(t)$ one can allow
aggregate returns to be heavy tailed, skewed and exhibit volatility clustering. Typical choices for
the marginal distribution of $X(t)$ are the Inverse Gaussian and Gamma distributions, alternatively
one may model directly the BDLP of $X(t)$ obtaining a variety of models which can adapt very well
to practical situations.

The basic model can be extended in a number of ways in order to get a better fit to real
phenomena.

A leverage effect can be introduced by adding a term $\bar{Z}(t) = \hat{Z}(t) - E[\hat{Z}(t)]$ to (1) as discussed
in Barndorff-Nielsen and Shephard (2003), i.e.
\[ dS(t) = \{\mu + \beta X(t)\}dt + \sqrt{X(t)}dN(t) + \rho d\bar{Z}(\lambda t) \] (3)

Note that $\bar{Z}(t)$ is the centered version of the BDLP driving the stochastic volatility process $X(t)$
and that $S(t)$ can have jumps.

More sophisticated correlation structures for $X(t)$ can be introduced via the sum of $p$ independent
OU processes, i.e.

\[ X(t) = \sum_{j=1}^{p} X^j(t), \quad dX^j(t) = -\lambda_j X^j(t) dt + d\bar{Z}(\lambda_j t) \]  

(4)

where \( \bar{Z}(\lambda_j t) \) are independent copies of the BDLP \( \bar{Z} \). In most cases of interest the distribution of \( X^j \) is additive with respect to some of its parameter values obtaining that the distribution of \( X \) belongs to the same class as \( X^j \). Long-range dependence could be introduced by allowing \( p \) to go to infinity. For examples and discussion see Leonenko and Taufer (2005). Leverage can also be combined with superpositions by using \( \sum_{j=1}^{p} \rho_j d\bar{Z}(\lambda_j t) \).

Estimating the parameters of these models is difficult because of the inability to compute the appropriate likelihood function. Much work on estimation has been devoted to model-based estimation approaches based on MCMC methods as in Roberts et al. (2004), Gander and Stephens (2006, 2007), Griffin and Steel (2006); the idea of using MCMC methods was already set forth in Barndorff-Nielsen and Shephard (2001a). These methods seem to work well in practice even though simulation of OU and, more generally, Lévy processes, is difficult owing to their jump character and one usually has to resort to approximated numerical procedures that may be very slow in some instances. For further details and references on simulation of Lévy processes see Todorov and Tauchen (2006), Taufer and Leonenko (2009).

Alternatively, as suggested by Barndorff-Nielsen and Shephard (2002), one might consider non-model-based estimation approaches which exploits realized volatility, i.e. use the existence of high-frequency intraday data to directly estimate moments of integrated volatility; for extensions of this approach and further references to recent works one can consult Woerner (2007).

In this paper we consider estimation based on the empirical characteristic function (ch.f.). This technique is well established; in particular, it has been shown that arbitrarily high levels of efficiency can be obtained given its direct relation with the likelihood. In the present context this technique
looks quite interesting and promising since it requires neither discretization nor simulation and
general forms of the ch.f. are available for Lévy processes. Building on existing results we will be
able to provide explicit general formulas for m- variate ch.f. for models of practical interest.

A seminal paper on the subject is that of Feuerverger and McDunnough (1981) which discusses
in detail the i.i.d. case and the extension to dependent observations. Other references that are
relevant here are Madan and Seneta (1987), Feuerverger (1990), Knight and Yu (2002), Jiang and
Knight (2002), Yu (2004) who discusses in more depth empirical ch.f. estimation in a non i.i.d.
setting.

Ch.f.-based estimation has been applied also in estimating the Lévy density for subordinators
by Jongbloed et al. (2005) in a non-parametric setting and and Jongbloed and Van der Meulen
OU processes. The work of this paper generalizes and extends some of these ideas in the context of
stochastic volatility models.

2 Characteristic function estimation

2.1 Background

To enter the estimation problem, suppose that the law of $S(t)$ is a member of a family of distributions
indexed by a vector of parameters $\theta \in \Theta \subset \mathbb{R}^q$. In terms of model (1) this accounts to estimating $\alpha,
\beta$, the autoregression parameter(s) $\lambda$ and the parameters of the marginal distribution of the latent
OU process $X(t)$.

Estimation will be based on equispaced observations $S_1, \ldots, S_n$ with $\Delta = t_j - t_{j-1}, j = 1, \ldots, n.$
Given stationarity of $S_1, \ldots, S_n$, estimation will be based on the joint ch.f. of blocks of $m$ consecutive
observations, hence define $\zeta = (\zeta_1, \ldots, \zeta_m)'$ and let

$$\psi_\theta(\zeta) = E[\exp\{i[\zeta_1 S_1 + \cdots + \zeta_m S_m]\}]$$

(5)

with

$$\psi_n(\zeta) = \frac{1}{n-m+1} \sum_{j=1}^{n-m+1} \exp\{i\zeta_1 S_j + i\zeta_2 S_{j+1} + \cdots + i\zeta_m S_{j+m-1}\}$$

the corresponding empirical estimator. The ch.f. estimator of $\theta$, say $\hat{\theta}$, is the vector minimizing the objective function

$$Q_n(\theta) = \int \cdots \int_S |\psi_n(\zeta) - \psi_\theta(\zeta)|^2 dW(\zeta),$$

(6)

where $S \subset \mathbb{R}^m$ denotes the region of integration; in the one-dimensional case $m = 1$ we may restrict attention to the case $S \subset \mathbb{R}^+$. $W(\zeta)$ is to be considered a weighting function which may serve different purposes and characterize the estimation procedure. If it is chosen to be a step function we turn into a discrete setup which is much easier to implement from the computational point of view. Usually one chooses a grid of points $\zeta$ at which the ch.f. are evaluated and then minimizes a sum instead of an integral. It is known that the choice of the grid of points has effect on the efficiency of the estimation procedure; Feuerverger and McDunnough (1981) show that, either in the i.i.d. or dependent case, using a weight given by the Fourier transform of the score and a grid sufficiently fine and extended, the ch.f.-based estimation procedure is asymptotically equivalent to maximum likelihood estimation. In general the choice of an appropriate weight and grid may be impractical given that the required quantities are seldom available explicitly and so one resorts to second best choices.

If instead we adopt a continuous weight function, the choice is important either for computational reasons, as it aims at damping out the persistent oscillations of the objective function as $\zeta \to \infty$, and for efficiency reasons. As pointed out by Singleton (2001), the estimator reaches the ML efficiency
if the weight function $W(X_t, r)$ is chosen as

$$w^*(X_t, \zeta) = \frac{1}{2\pi} \int \exp(-i\zeta X_{t+1}) \frac{\partial \log f_\theta(X_{t+1}|X_t)}{\partial \theta} dX_{t+1}.$$ 

It can be seen that the optimal choice requires the computation of the conditional score function, which is generally impractical in our setup. Knight and Yu (2002) propose an exponentially decreasing function, Epps (2005) suggests using $dW(\zeta) = (|\phi_\theta(\zeta)|^2 / \int |\phi_\theta(v)|^2 dv)d\zeta$ and Jiang and Knight (2002) use the multivariate normal density. However, for non-Markovian processes, no general solution is available.

### 2.2 General forms of the ch.f.

We need some more notation at this point. For the stationary OU process $X(t)$ and the BDLP $\hat{Z}(t)$ define the respective ch.f. and cumulant functions as

$$\phi(\zeta) = E(e^{i\zeta X(t)}), \quad \kappa(\zeta) = \log \phi(\zeta),$$

$$\hat{\phi}(\zeta) = E(e^{i\zeta \hat{Z}(t)}), \quad \hat{\kappa}(\zeta) = \log \hat{\phi}(\zeta).$$

Recall (see Barndorff-Nielsen and Shephard (2001b)) that we have the relation

$$\hat{\kappa}(\zeta) = \zeta \kappa'(\zeta), \quad \text{with } \kappa'(\zeta) = \frac{d\kappa(\zeta)}{d\zeta}. \quad (7)$$

We will also need an expression for the cumulant generating functions which we will use in the form $k(\zeta) = \kappa(i\zeta)$ and $\hat{k}(\zeta) = \hat{\kappa}(i\zeta)$.

In order to implement ch.f.-based estimation we need to have an explicit form for the ch.f. of the above models to be used in (5) and (6). A fundamental result in this sense can be obtained from Barndorff-Nielsen and Shephard (2001b), Theorem 10.1, which gives the general form of the ch.f. of a function of $S$. For some integrable $f$ denote $f \cdot S = \int_0^\infty f(t) dS(t)$. The integral $f \cdot S$ is...
well defined under some conditions (see, for example Anh et al. (2002)) which are satisfied in all cases of interest here. For \( S(t) \) defined as in (3) it holds

\[
\log \mathbb{E}(e^{i\zeta f S}) = \lambda \int_0^\infty \left[ k(J e^{-\lambda s}) + \hat{k}(H(s)) \right] ds + i\zeta(\mu - \lambda \rho \xi) \int_0^\infty f(s) ds
\]

where \( \xi = \mathbb{E}\hat{Z}(1) \) is the expected value of \( \hat{Z}(1) \), \( J = G(0) \), \( H(s) = G(s) + i\zeta \rho f(s) \) with

\[
G(s) = \int_0^\infty \left[ \frac{1}{2} \zeta^2 f^2(s + u) - i\zeta^2 f(s + u) \right] e^{-\lambda s} du.
\]

It is easy to extend the result to the case of superpositions considered in (4) by using independence. For \( S(t) \) defined in (1) with \( X(t) \) defined by (4) we have

\[
C\{\zeta f \cdot S\} = \sum_{j=1}^p \lambda_j \int_0^\infty \left[ k(J(\lambda_j) e^{-\lambda_j s}) + \hat{k}(G(\lambda_j)(s)) \right] ds + i\zeta \mu \int_0^\infty f(s) ds.
\]

Note that dependence on \( \lambda_j \) of \( J \) and \( G(s) \) has been introduced in this case.

These formulas may be hard to implement, we need to develop general and tractable expressions of (8) for functions \( f \) of interest in estimation. We do it in two steps. First by providing some general forms which are useful in understanding the problems at hand; next by developing further expressions ready for numerical implementation.

The following Theorem obtains the joint ch.f. of the increments of the stochastic volatility process (1) in terms of the cumulant generating function of the BDLP \( \hat{Z}(1) \) of \( X(t) \).

**Theorem 1.** Let \( S(t) \) be defined by (1), the joint ch.f. of \( S_1, \ldots, S_m \) is

\[
\exp \left\{ i\mu \sum_{j=1}^m \zeta_j (t_j - t_{j-1}) \right\} \exp \left\{ \lambda \int_0^\infty k(J e^{-\lambda s}) ds \right\} \exp \left\{ \lambda \sum_{l=0}^m \int_{t_l}^{t_{l+1}} \hat{k}(G_l(s)) ds \right\},
\]

where

\[
J = \sum_{j=1}^m \left( \frac{1}{2} \zeta_j^2 - i\beta \zeta_j \right) \varepsilon_\lambda(t_{j-1}, t_j),
\]

\[
G_l(s) = \sum_{j=l+1}^m \left( \frac{1}{2} \zeta_j^2 - i\beta \zeta_j \right) \varepsilon_\lambda(t_{j-1}, t_j) e^{\lambda s} + \mathbb{I}_{(l>0)} \left( \frac{1}{2} \zeta_l^2 - i\beta \zeta_l \right) \varepsilon_\lambda(0, t_l - s).
\]

Here \( \varepsilon_\lambda(u, v) = \frac{1}{\lambda}(e^{-\lambda u} - e^{-\lambda v}) \). Moreover, we use the conventions \( t_{-1} = 0 \) and \( \sum_{j=m+1}^n f(j) = 0 \).
Proof of Theorem 1. To obtain our result, we need to apply formula (8) setting \( \zeta = 1 \) and \( f(s) = \sum_{j=1}^{m} \zeta_j \mathbb{I}_{(t_j-1,t_j]}(s) \). Note first of all that, since \((t_{j-1}, t_j] \) and \((t_{j'-1}, t_{j'}] \) do not overlap, for \( j \neq j' \), then \( f^2(s) = \sum_{j=1}^{m} \zeta_j^2 \mathbb{I}_{(t_j-1,t_j]}(s) \), this allows to derive straightforwardly expression (10). To get the expression for \( H(s) \) we write

\[
\mathbb{I}_{(t_j-1,t_j]}(u + s) = \mathbb{I}_{(t_j-1-s,t_j-s]}(u) \mathbb{I}_{(0,t_j-1]}(s) + \mathbb{I}_{(0,t_j-s]}(u) \mathbb{I}_{(t_j-1,t_j]}(s)
\]

to obtain

\[
H(s) = \sum_{j=1}^{m} \left( \frac{1}{2} \zeta_j^2 - i \beta \zeta_j \right) \left[ \varepsilon_\lambda(t_{j-1} - s, t_j - s) \mathbb{I}_{(0,t_j-1]}(s) + \varepsilon_\lambda(0, t_j - s) \mathbb{I}_{(t_j-1,t_j]}(s) \right].
\]

In order to compute \( \int k(G(s))ds \) it is convenient to separate the components of \( G(s) \) over the disjoint intervals \((t_{j-1}, t_j] \) and proceed to integrate over separate regions. To this end, write \( \mathbb{I}_{(0,t_j-1]}(s) = \sum_{l=0}^{j-1} \mathbb{I}_{(u_{l-1},u_l]}(s) \) with the convention that \( t_{-1} = 0 \) and rearrange terms around common factors to obtain (11).

In the case of leverage we state it as follows:

**Theorem 2.** Let \( S(t) \) be defined by (3), the joint ch.f. of \( S_1, \ldots, S_m \) is

\[
\exp \left\{ i(\mu - \lambda \xi \rho) \sum_{j=1}^{m} \zeta_j (t_j - t_{j-1}) \right\} \exp \left\{ \lambda \int_0^\infty k(Je^{-\lambda s})ds \right\} \exp \left\{ \lambda \sum_{l=0}^{m} \int_{t_{l-1}}^{t_l} \dot{k}(H_l(s))ds \right\},
\]

where

\[
H_l(s) = G_l(s) + \rho \zeta_l \mathbb{I}_{(t_{l-1},t_l]}(s), \quad l = 0, \ldots, m, \quad \zeta_0 = 0,
\]

and the quantities \( J, G_l(s) \) and \( \varepsilon_\lambda(u,v) \) are defined as in Theorem 1.

Finally, when the process \( X(t) \) is the sum of \( p \) independent components, the ch.f. of the stochastic volatility model is easily obtained by composing the different contributions. The following theorem gives the form of the ch.f. of the process.
Theorem 3. Let $S(t)$ be defined by (1) with $X(t)$ as in (4), the joint ch.f. of $S_1, \ldots, S_m$ is

$$
\exp \left\{ \sum_{j=1}^{p} \lambda_j \left[ \int_0^\infty \tilde{k}_j(J^{(\lambda_j)}e^{-\lambda_js})ds + \sum_{l=0}^{m} \int_{t_{l-1}}^{t_l} \tilde{k}_j(G^{(\lambda_j)}(s))ds \right] + i\mu \sum_{j=0}^{m} \int_0^\infty \zeta_j \mathbb{I}(t_{j-1}, t_j)(s)ds \right\} = \exp \left\{ \sum_{j=1}^{p} \lambda_j \left[ \int_0^\infty \tilde{k}_j(J^{(\lambda_j)}e^{-\lambda_js})ds + \sum_{l=0}^{m} \int_{t_{l-1}}^{t_l} \tilde{k}_j(G^{(\lambda_j)}(s))ds \right] + i\mu \sum_{j=0}^{m} \int_{t_{j-1}}^{t_j} \zeta_j ds \right\}.
$$

(13)

2.3 Consistency and asymptotic distribution

Consistency and asymptotic normality of the estimators presented here follow from Knight and Yu (2002), who discuss ch.f. estimation-based methods for dependent observations, if their conditions A1-A8 are satisfied. These assumptions concern identifiability and regularity issues as well as specific convergence conditions on the process that generated the sequence of observations; the reader is referred to the aforementioned paper for a detailed listing and for the formulas of the asymptotic variance of the estimators.

Turning to our case, identifiability and regularity conditions will generally hold for the estimators because of linear independence of the functions $e^{i\zeta_j x}$.

Among the remaining assumptions, the hardest one to verify is condition A7, which postulates mean square convergence of a sequence of zero-mean martingale differences constructed on the process. However, according to the remark of Knight and Yu (2002, p. 696), A7 holds under suitable mixing conditions which are satisfied in our case.

In fact, for the stochastic volatility model (1) $\alpha$-mixing of the volatility process $X$ implies $\alpha$-mixing of the observed return process $S$, and this is actually true for models (3) and (4) as well. To see this, one needs to consider the $\sigma$-algebra $\mathcal{F}_t$ generated by the background driving Lévy process $\hat{Z}$ and then, by previous conditioning on $\mathcal{F}_t$, to follow the same way of reasoning as in Lemma 6.3 in Sørensen (2000). That $X$ is $\alpha$-mixing follows from Jongbloed et al. (2005), who proved that the
OU process is $\beta$-mixing (and hence also $\alpha$-mixing) under the condition
\[
\int_2^{\infty} \log(x)\nu(dx) < \infty,
\]
where $\nu$ is the Lévy measure of $X$. Note that we have $\beta$-mixing with exponential rate either for OU processes and their superpositions since $\sigma$- algebras are equivalent. Conditions for $\beta$-mixing holds for the Tempered Stable and Gamma cases discussed here.

2.4 Computational issues

If we define the quantities
\[
c(l, m) = \sum_{j=l}^{m} \left( \frac{1}{2} \xi_j^2 - i\beta \xi_j \right) \varepsilon_\lambda(t_{j-1}, t_j) \quad \text{and} \quad c(l) = \frac{1}{2} \xi_l^2 - i\beta \xi_l,
\]
the functions $G_l(s)$ defined in Theorems 1 to 3, can be written in a form more suitable for computations by rearranging terms not involving the variable of integration $s$, i.e.
\[
G_l(s) = c(l + 1, m)e^{\lambda s} + c(l)\mathbb{I}_{(l>0)}\varepsilon_\lambda(0, t_l - s) = \frac{1}{\lambda}c(l)\mathbb{I}_{(l>0)} + \left[ c(l + 1, m) - \frac{1}{\lambda}c(l)\mathbb{I}_{(l>0)}e^{-\lambda t_l} \right] e^{\lambda s} = c_1(l) + c_2(l, m)e^{\lambda s}, \quad l = 0, \ldots, m.
\]

If a leverage term is present, then, see Theorem 2, it can be simply incorporated into the term $c_1(l)$ and will not change the discussion below. Therefore, overall computation of the ch.f. of $S_1, \ldots, S_m,$ for $t_j - t_{j-1} = \Delta, \; j = 1, \ldots, m$ can be reduced to the following form
\[
\phi_{\mu, \lambda, \xi, \rho}(\xi_1, \ldots, \xi_m) = \exp \left\{ i(\mu - \lambda \xi \rho)\Delta \sum_{j=1}^{m} \xi_j \right\} \exp \left\{ A + \sum_{l=0}^{m} B_l \right\}, \quad \text{(14)}
\]
where the terms $A$ and $B_l, \; l = 0, \ldots, m$ can be computed by integrals of the form
\[
A = \lambda \int_0^{\infty} \hat{k}(c_2(0, m)e^{-\lambda s}) ds \quad \text{and} \quad B_l = \lambda \int_{t_{l-1}}^{t_l} \hat{k} \left( c_1(l) + c_2(l, m)e^{\lambda s} \right) ds.
\]
Note that exploiting (7) we have \( \tilde{k}(\zeta) = \zeta k'(\zeta) \); hence for \( A \) we have the estimate

\[
A = \lambda \int_0^\infty c_2(0, m)e^{-\lambda s}k'(c_2(0, m)e^{-\lambda s})ds = -k(c_2(0, m)e^{-\lambda s}) \bigg|_0^\infty = k(c_2(0, m)).
\] (15)

Similarly for the terms \( B_l \), using the same devices, we have the estimate

\[
B_l = \lambda \int_{l_{l-1}}^{l_l} c_1(l)k'(c_1(l) + c_2(l, m)e^{\lambda s})ds + \lambda \int_{l_{l-1}}^{l_l} c_2(l, m)e^{\lambda s}k'(c_1(l) + c_2(l, m)e^{\lambda s})ds = \lambda \int_{l_{l-1}}^{l_l} c_1(l)k'(c_1(l) + c_2(l, m)e^{\lambda s})ds + k(c_1(l) + c_2(l, m)e^{\lambda t_l}) - k(c_1(l) + c_2(l, m)e^{\lambda t_{l-1}}).
\]

Overall, by noting that \( c_1(0) = 0 \), some further simplifications can be done when \( l = 0 \). Finally, one obtains the following expression to be substituted for \( A + \sum_{l=0}^m B_l \) in (14):

\[
A + \sum_{l=0}^m B_l = k(c_2(0, m)e^{\lambda t_0}) + \sum_{l=1}^m \left[ k(c_1(l) + c_2(l, m)e^{\lambda t_l}) - k(c_1(l) + c_2(l, m)e^{\lambda t_{l-1}}) \right] + \sum_{l=1}^m \lambda \int_{l_{l-1}}^{l_l} c_1(l)k'(c_1(l) + c_2(l, m)e^{\lambda s})ds.
\] (16)

As we see, to completely explicit the formula we need to perform one last integration which depends on the functional form of \( k'(\cdot) \); the above forms imply that the computational complexity of the ch.f. of Theorems 1 to 3 is the same for any given \( m \)-value. Overall, the given formula is well apt for numerical integration if necessary.

3 Applications

3.1 Marginal models for \( X(t) \)

To get more specific into our modelling framework we present two important models for the marginal distribution of the volatility process which allow to obtain a variety of behaviours for the process \( S(t) \). A general reference for the material of this section is Schoutens (2003).

We have provided expressions for the ch.f. of \( S_1, \ldots, S_m \) in terms of \( \tilde{k}(\zeta) \) or \( k(\zeta) \), hence one can model directly from the BDLP or alternatively, if it is preferable, directly by the marginal distribution for the process \( X(t) \). The two approaches are linked by the relation (7).
A first important model for $X$ is the so-called Tempered Stable which we denote by $X \sim TS(\nu, \delta, \gamma)$, for $0 < \nu < 1$ and $\delta, \gamma > 0$. The ch.f. takes the form

$$\phi_{TS(\nu, \delta, \gamma)}(\zeta) = \exp \left\{ \delta \gamma - \delta (\gamma^{1/\nu} - 2i\zeta)^\nu \right\}.$$  

When $\nu = 1/2$ we have the important sub-case of the Inverse Gaussian distribution which we denote by $X \sim IG(\delta, \gamma)$. For the $TS$ case, exploiting (7) we obtain

$$\hat{k}(\zeta) = -2\zeta \nu \delta^{2\nu}(\gamma^{1/\nu} + 2\zeta)^{-\nu}. \quad (17)$$

Another important case is the Generalized Inverse Gaussian ($GIG$) model, which we denote by $X \sim GIG(\nu, \delta, \gamma)$, with $\delta \geq 0, \gamma > 0$ if $\nu > 0$, $\delta > 0, \gamma > 0$ if $\nu = 0$ and $\delta > 0, \gamma \geq 0$ if $\nu < 0$. This model has ch.f.

$$\phi(\zeta)_{GIG(\nu, \delta, \gamma)} = \left( \frac{\gamma^2}{\gamma^2 - 2i\zeta} \right)^{\nu/2} \frac{K_\nu(\delta \sqrt{\gamma^2 - 2i\zeta})}{K_\nu(\delta \gamma)},$$

where $K_\nu(x)$ denotes the modified Bessel function of the third kind. The appropriate form of $\hat{k}(\zeta)$ or $k'(\zeta)$ for the $GIG$ model, not shown here, can be recovered by (7). Note that

$$\frac{\partial K_\nu(x)}{\partial x} = -\frac{1}{2}(K_{\nu-1}(x) + K_{\nu+1}(x)).$$

In the limiting case $\nu > 0, \delta = 0$ it reduces to the density of a Gamma distribution $\Gamma(\gamma^2, \nu/2)$; for $\nu < 0, \gamma = 0$ one gets those of a reciprocal Gamma distribution. For the Gamma case we have the simple form

$$k(\zeta) = \frac{\nu \zeta}{2(\gamma^2 + \zeta)}.$$

Notably the IG distribution is common either to the $TS$ and to the $GIG$ distribution, being $IG(\delta, \gamma) \sim TS(1/2, \delta, \gamma) \sim GIG(-1/2, \delta, \gamma)$.

In general the $GIG$ is not closed under convolution and so, apart from some special cases, cannot be used to model superpositions of OU processes.
3.2 Example: the Tempered Stable scenario

Figure 1: Contour plots of the objective function (6) for model (1) with \( X(t) \) following an \( IG-OU \) process with \( \delta = 2, \gamma = 1, \lambda = 0.1 \). Length of series \( n = 1000 \). Row 1: univariate ch.f’s are used; Row 2: bivariate; Row 3: trivariate. Column 1: contour plots for \( (\delta, \gamma) \); Column 2: contour plots for \( (\delta, \lambda) \); Column 3: contour plots for \( (\gamma, \lambda) \). In all cases, the third parameter is set to its true value.
In this section we provide computational formulas and some simulation results on the important case of the Tempered Stable scenario for $X(t)$.

From (17) we get
\[ k'(\zeta) = \frac{\hat{k}(\zeta)}{\zeta} = -2\nu \delta^{2\nu} (\gamma^2 + 2\zeta)^{\nu-1}. \] (18)

Using (18) in the last addend of (16), we note that the integral to be computed is of the form $\int (1 + ce^{\lambda s})^{\nu-1} ds$; by appropriate transformations it can be given explicit form in terms of the hypergeometric function
\[ F(a, b, c, z) = \frac{\Gamma(c)}{\Gamma(a)\Gamma(b)} \sum_{n=0}^{\infty} \frac{\Gamma(a+n)\Gamma(b+n)}{\Gamma(c+n)} \frac{z^n}{n!}. \]

For details about convergence conditions and integral representations of this function we refer the reader to Abramowitz and Stegun (1972), pp. 556-558. In more detail for our case a solution can be obtained by
\[ \int (1 + ce^{\lambda s})^{\nu-1} ds = \frac{1}{(\nu - 1)\lambda} \left( \frac{1 + e^{\lambda s}/c}{1 + ce^{\lambda s}} \right)^{1-\nu} F(1 - \nu, 1 - \nu, 2 - \nu, -e^{-\lambda s}/c). \]

When $\nu = 1/2$ one obtains the IG stochastic volatility model. This case can be considerably simplified; in particular, we report below the cumulant function of $S_1, \ldots, S_m$ for $m = 1$. Letting $AT$ denote the ArcTanh function and $t_j - t_{j-1} = \Delta$, we have
\[ \log(E(e^{i\zeta S_1})) = \frac{2\delta \left( \frac{\zeta^2}{2} - i\zeta \beta \right)}{\lambda} \left[ AT \left( (B^{-1}) \sqrt{1 + \frac{2(1-e^{-\Delta\lambda})(|\zeta^2/2-i\zeta\beta|)}{\gamma^2\lambda}} \right) - AT(B^{-1}) \right], \]
where $B = \sqrt{1 + \frac{2(\zeta^2/2-i\zeta\beta)}{\gamma^2\lambda}}$. As we see, although cumbersome, the above function is straightforwardly applicable in numerical procedures.

When $m > 1$, the ch.f. of $S_1, \ldots, S_m$ can be obtained by means of (14). Note that the last integral in (16) is given by
\[ \int_{t_{i-1}}^{t_i} c_1(l)k' \left( c_1(l) + c_2(l, m)e^{\lambda s} \right) ds = -\delta c_1(l) \int_{t_{i-1}}^{t_i} (\gamma^2 + 2(c_1(l) + c_2(l, m)e^{\lambda s}))^{-1/2} ds, \]
whose explicit solution is equal to

\[-\delta c_1(l) \left( -2AT \left( \frac{(\gamma^2 + 2c_1(l) + 2c_2(l,m)e^{\lambda s})^{1/2}}{(\gamma^2 + 2c_1(l))^{1/2}} \right) \right) / \lambda(\gamma^2 + 2c_1(l))^{1/2} \bigg|_{t_{l-1}}.

Here we give a brief report on some results with simulated data by estimating the parameters for model (1) with \( \mu = \beta = 0 \) and where \( X(t) \) follows an \( IG(\delta, \gamma) \) OU process with \( \delta = 2, \gamma = 1 \) and auto-regression parameter \( \lambda = 0.1 \). We generated a series of observations of length 1000 of the volatility process \( X(t) \) as described in Taufer and Leonenko (2009) and successively used it in a Euler scheme to generate \( S \). Figure 1 depicts the contour plots of the objective function. Each row indicates use of the \( m \)-variate ch.f., \( m = 1, 2, 3 \); each column shows a possible pair of the three target parameters.

Although this example is based on a single series, the results are quite promising, also in the light of the little time used in the minimizing procedure. Inspection of the contour plots reveals quite a regular behavior of the objective function, with very low values in the neighborhood of the true parameter values.

4 Numerical results

The preceding sections provide an analysis of the most relevant results and open problems concerning ch.f.-based estimation of non-Gaussian OU processes. There is in particular one issue that requires further investigation. We have discussed consistency and asymptotic normality of the estimators; however, no information is available about the rate of convergence of the estimators, so that the actual validity of these results in finite samples has to be investigated numerically. In this and the next section we limit ourselves to the case of the Inverse Gaussian distribution. We first consider the case \( m = 2, p = 1 \) (no superposition) and then extend the analysis to the case \( m = 2, p = 2 \), corresponding to the superposition of two OU processes.
As pointed out in section 2.1, finding the optimal weighting function $W(\cdot)$ in (6) is not trivial in a framework like the present one, where the observations are not i.i.d. As no general solution is available (Jiang and Knight 2002, pag. 206), we will adopt a strategy that has frequently been used in applications and has proved to have reasonable properties in similar frameworks, namely an exponentially decreasing weight function.

4.1 Simulation experiment 1: no-superposition

The results of the no-superposition case are shown in Table 1 and Figure 2. At each replication we simulate the process with $T = 1500$ and discard the first 500 observations, in order to eliminate the effect of the starting point. Then we minimize the objective function (6); the integral is computed numerically by means of the R function `adapt`, based on the algorithm developed by Berntsen et al. (1991). The minimization is performed by the R function `constrOptim`, which employs the Nelder-Mead simplex method. Because of the heavy computational burden we limit ourselves to the simulation of $B = 100$ replications of the process (1) with $X(t)$ given by (2). The values of the parameters used to simulate the process are $\delta = 2$, $\gamma = 1$, $\lambda = 0.1$ and $\beta = 0.03$. Table 1 shows the sample mean, standard deviation, skewness and kurtosis computed across the 100 replications. Figure 2 displays the simulated distributions of the estimators.

<table>
<thead>
<tr>
<th></th>
<th>Mean</th>
<th>Standard deviation</th>
<th>Skewness</th>
<th>Kurtosis</th>
</tr>
</thead>
<tbody>
<tr>
<td>$\delta$</td>
<td>1.96</td>
<td>0.259</td>
<td>0.810</td>
<td>4.551</td>
</tr>
<tr>
<td>$\gamma$</td>
<td>1.060</td>
<td>0.172</td>
<td>0.390</td>
<td>4.484</td>
</tr>
<tr>
<td>$\lambda$</td>
<td>0.177</td>
<td>0.093</td>
<td>1.176</td>
<td>5.746</td>
</tr>
<tr>
<td>$\beta$</td>
<td>0.030</td>
<td>0.030</td>
<td>-0.434</td>
<td>3.881</td>
</tr>
</tbody>
</table>

Despite the small number of replications, the results are quite interesting. In particular, the estimators seem to be approximately unbiased and the simulated distributions are bell-shaped.
Figure 2: Simulation results with $p = 1$
4.2 Simulation experiment 2: superposition of two OU processes

The simulation performed in the preceding subsection is repeated here when $X(t)$ is the sum of 2 independent OU processes (see (4)). The remaining settings of the experiment are the same as before, with the exception of the parameter $\beta$, which is now set to zero. Results are shown in Table 2 and Figure 3.

**Table 2**: Simulation results with $p = 2$ (true values: $\delta_1 = 1.5$, $\delta_1 = 0.5$, $\gamma = 1$, $\lambda_1 = 1.5$, $\lambda_2 = 0.5$).

<table>
<thead>
<tr>
<th></th>
<th>Mean</th>
<th>Standard deviation</th>
<th>Skewness</th>
<th>Kurtosis</th>
</tr>
</thead>
<tbody>
<tr>
<td>$\delta_1$</td>
<td>1.48</td>
<td>0.24</td>
<td>0.15</td>
<td>5.39</td>
</tr>
<tr>
<td>$\delta_2$</td>
<td>0.35</td>
<td>0.19</td>
<td>1.19</td>
<td>7.31</td>
</tr>
<tr>
<td>$\gamma$</td>
<td>0.93</td>
<td>0.13</td>
<td>-0.23</td>
<td>3.01</td>
</tr>
<tr>
<td>$\lambda_1$</td>
<td>1.76</td>
<td>0.51</td>
<td>1.80</td>
<td>10.42</td>
</tr>
<tr>
<td>$\lambda_2$</td>
<td>0.6</td>
<td>0.33</td>
<td>1.61</td>
<td>6.78</td>
</tr>
</tbody>
</table>

The results are comparable to the case $p = 1$, with just a small loss of precision due to the larger number of parameters. Only the estimators of $\lambda_1$ and $\lambda_2$ seem to be less stable than the estimator of $\lambda$ in the no-superposition setup: the standard deviation of $\hat{\lambda}_1$ is more than five times larger than the standard deviation when $p = 1$.

4.3 A real-data application: the S&P’s 500 volatility

In this subsection we fit model (1) to the daily log-prices of the Standard & Poor’s index over the 5-years period from December 1, 2003 to December 1, 2008. We consider either the case where the volatility process $X(t)$ is a single Inverse-Gaussian OU process ($p = 1$) and the case where it is the superposition of two independent Inverse-Gaussian OU processes ($p = 2$). The first two panels of figure 4 display respectively the time series of levels and of log-returns; the remaining two panels show the autocorrelation function of the returns and of the squared returns.
Figure 3: Simulation results with $p = 2$
Figure 4: Some descriptive statistics about the S&P’s 500 index
Table 3 reports some basic descriptive statistics. Most well-known stylized facts typically observed in financial data are evident: in particular, there are clusters of volatility and the data are strongly leptokurtic. Some other features of the data are strictly related to the downward trend of the last three months, caused by the financial crisis that took place in the second half of the year 2008. There is little evidence of autocorrelation of the log-returns series, with a couple of negative values at lags 1 and 2; on the contrary, the squared log-returns are strongly correlated, with significant values approximately up to lag 40. In addition, the negative skewness is probably a consequence of the large negative returns often observed in the last three months of the series.

<table>
<thead>
<tr>
<th>Mean</th>
<th>Standard deviation</th>
<th>Skewness</th>
<th>Kurtosis</th>
</tr>
</thead>
<tbody>
<tr>
<td>-0.0002</td>
<td>0.0129</td>
<td>-0.4672</td>
<td>20.8052</td>
</tr>
</tbody>
</table>

Table 3: Basic descriptive statistics

Estimation of parameters was performed as in sections 4.1 and 4.2 with bivariate ch.fs and an exponentially decreasing weight function. Table 4 shows the estimates and the value of the objective function (6) at convergence obtained when implementing model (1) respectively with \( p = 1 \) and \( p = 2 \) superpositions.

Table 4: Estimation results

<table>
<thead>
<tr>
<th>( \delta_1 )</th>
<th>( \delta_2 )</th>
<th>( \gamma )</th>
<th>( \lambda_1 )</th>
<th>( \lambda_2 )</th>
<th>( \beta )</th>
<th>obj. function</th>
</tr>
</thead>
<tbody>
<tr>
<td>( p = 1 )</td>
<td>0.657</td>
<td>-0.615</td>
<td>0.278</td>
<td>-0.055</td>
<td>0.0017</td>
<td></td>
</tr>
<tr>
<td>( p = 2 )</td>
<td>0.327</td>
<td>0.309</td>
<td>0.589</td>
<td>0.065</td>
<td>0.801</td>
<td>0.056</td>
</tr>
</tbody>
</table>

The most relevant difference observed when moving from \( p = 1 \) to \( p = 2 \) concerns the parameters of the volatility processes, namely \( \lambda_1 \) and \( \lambda_2 \). It can indeed be seen that \( \lambda_2 \) is much larger than \( \lambda \), whereas \( \lambda_1 \) is very small. A plausible interpretation is that the two processes of the case \( p = 2 \) roughly correspond respectively to small, frequent and fast-moving shocks (the process with larger
λ) and to large, rare and persistent shocks (the process with smaller λ). Analogous results about these parameters were obtained by Griffin and Steel (2006, p. 627). When \( p = 1 \), an intermediate value of \( \lambda \) is obtained because a single process has to account for all the movements. Note that since the sum of \( \delta_1 \) and \( \delta_2 \) in the case \( p = 2 \) is approximately equal to \( \delta_1 \) when \( p = 1 \) and \( \gamma \) is nearly unchanged, the marginal distribution of the volatility remains essentially the same for the two models. The risk-premium coefficient \( \beta \) is positive in both cases, as expected from financial theory; its relatively large value is in line with the high level of risk in the period considered in the present application.

With the aim of evaluating the goodness-of-fit in the cases \( p = 1 \) and \( p = 2 \), we now compare the theoretical ch.f.s and autocorrelation functions of the estimated processes with their empirical counterparties.

Consider first the ch.f.s. Figure 5 shows three functions: each of them is the sum of the squares of the real and imaginary parts of a ch.f.. The first one (solid line) is the empirical ch.f.; the second one (dashed line) is the estimated ch.f. with \( p = 1 \); the third one (dotted line) is the same function with \( p = 2 \). In all of these cases we let \( \zeta_1 \) vary in the interval \([0, 5]\) and fix \( \zeta_2 = 0 \). The two estimated functions are virtually indistinguishable and in good agreement with the empirical ch.f. This seems to suggest that the marginal distribution is estimated quite precisely and that the approach works well for either value of \( p \).

Next we turn to the dependence structure of the processes: especially when the squared log-returns exhibit a strong autocorrelation, the analysis of the autocorrelation of the observed and estimated processes is an important measure of goodness-of-fit. Notice that the theoretical autocorrelation function of the squared log-returns can be obtained using the results in Barndorff-Nielsen and Shephard (2001a, sect. 4.2). In this case the evidence is quite different. Figure 6 shows the
Figure 5: Empirical and theoretical characteristic functions
Figure 6: Empirical and theoretical autocorrelation functions of squared log-returns
theoretical autocorrelation functions of the estimated processes for $p = 1$ and $p = 2$: for comparison purposes, we also reported the empirical autocorrelation function of the observed squared returns. It is clear that with $p = 2$ the autocorrelation function decreases much more slowly and therefore captures more precisely the autocorrelation in the data. This allows to conclude that adding a second component improves considerably the fit to the dependence structure of the returns, while, as seen above, the marginal distribution is estimated approximately in the same way by the two models.

5 Conclusions

In this paper we have considered the problem of estimating the parameters of a continuous-time stochastic volatility model where the latent volatility follows an OU process. These processes are well suited for modelling some typical stylized facts observed in financial data. However, their estimation is quite difficult because usually the likelihood function cannot be written in closed form. Here we have developed an estimation method that makes use of the ch.f.; this technique is quite appealing because it is asymptotically equivalent to maximum likelihood.

The main contribution of the paper consists in extending some results first given by Barndorff-Nielsen and Shephard (2001a) in order to obtain the explicit expression of the ch.f. for various choices of the marginal distribution of the volatility process. We have investigated, by means of some simulation experiments, the small sample behavior of the estimators. Finally the method has been applied to the estimation of the volatility process of the daily time series of the S&P’s 500 equity index over the period 2003-2008.

From a theoretical point of view the method proposed in this article is convenient because it avoids both discretization and simulation. The results of the Monte Carlo experiments confirmed
the validity of the approach. As for the real-data analysis, the autocorrelation function of the squared log-returns decays very slowly. In accordance to other previous studies, in this case the superposition of two OU processes provides a better fit than a single OU process, allowing to estimate more precisely the dependence structure of the process.

A possible extension of these models is via the use of continuous superpositions where one could also introduce long-memory in the volatility process (see Barndorff-Nielsen and Leonenko (2005) for details). The technique of estimation adopted here seems feasible in these cases even though it may be difficult to obtain explicit expressions; instead one would be able to get expressions for cumulants and the approach the problem of estimation by means of GMM.

References


Singleton, K.J. (2001). Estimation of affine asset pricing models using the empirical characteristic
function. *J. Econometrics* 102, 111-141.


Characteristic function estimation of Ornstein-Uhlenbeck-based stochastic volatility models

Emanuele Taufer, Nikolai Leonenko, Marco Bee