# Testing finite activity against infinite activity for jumps，for high frequency observation：an overview 

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Our aim is to describe，in a somewhat informal way，a number of preliminary results： In this presentation we use quite strong assumptions，hoping to significantly weaken them in the near future，and we provide no proof，neither empirical or simulation studies，those being not complete by now．

The problem concerns processes having jumps，since it seems more and more clear that models in finance should take jumps into consideration．Traditionally，such models with jumps rely on Poisson or compound processes，as in［11］，［3］and［4］．However，more recently some financial models have been proposed，that allow for an infinite number of jumps in finite time intervals，such as the variance gamma model of［10］and［9］，the hyperbolic model of［7］，the pure jump model of［5］and the finite moment log stable process of［6］．These models can capture both small and frequent jumps，as well as large and infrequent ones．Since the qualitative properties and the mathematical analysis of models with finitely many jumps deeply differ from those of models with infinitely many jumps，it seems appropriate to develop some methods which can discriminate between the two types of models．

In an attempt to bring forth some contribution to these questions，we aim here to develop testing procedures to discriminate empirically between the two situations of finite and infinite number of jumps．This is an important information，which leads to qualita－ tively very different models．We have a 1－dimensional $X$ observed on a fixed time interval $[0, T]$ ，at discretely and regularly spaced times $i \Delta_{n}$ ．Assuming the observed path has jumps，we want to test whether there are a finite number of jumps or not，two properties commonly referred to as＂finite activity＂and＂infinite activity＂for the jump component of $X$ ．These properties are in fact satisfied，respectively，on two complementary subsets of the sample space $\Omega$ ：

$$
\begin{align*}
& \Omega_{T}^{f}=\left\{\omega: t \mapsto X_{t}(\omega) \text { has finitely many jumps in }[0, T]\right\} \\
& \Omega_{T}^{i}=\left\{\omega: t \mapsto X_{t}(\omega) \text { has infinitely many jumps in }[0, T] .\right\} \tag{1}
\end{align*}
$$

We need some assumptions on the process $X$ ．However，we wish to keep the solution as

[^0]nonparametric as possible, and in particular we do not want to specify the structure of the volatility, and as little as possible about the structure of the jumps.

1) Description of the Model: The underlying process $X$ is a 1-dimensional Itô semimartingale on some filtered space, which means that its characteristics ( $B, C, \nu$ ) are absolutely continuous with respect to Lebesgue measure. $B$ is the drift, $C$ is the quadratic variation of the continuous martingale part, and $\nu$ is the compensator of the jump measure $\mu$ of $X$. In other words, we have

$$
\begin{equation*}
B_{t}(\omega)=\int_{0}^{t} b_{s}(\omega) d s, \quad C_{t}(\omega)=\int_{0}^{t} \sigma_{s}(\omega)^{2} d s, \quad \nu(\omega, d t, d x)=d t F_{t}(\omega, d x) \tag{2}
\end{equation*}
$$

Here $b$ and $\sigma$ are optional process, and $F=F_{t}(\omega, d x)$ is a transition measure from $\Omega \times \mathbb{R}_{+}$ endowed with the predictable $\sigma$-field into $\mathbb{R} \backslash\{0\}$.

Equivalent to (2), one may write $X$ as
$X_{t}=X_{0}+\int_{0}^{t} b_{s} d s+\int_{0}^{t} \sigma_{s} d W_{s}+\int_{0}^{t} \int x 1_{\{|x| \leq 1\}}(\mu-\nu)(d s, d x)+\int_{0}^{t} \int x 1_{\{|x|>1\}} \mu(d s, d x)$.
where $W$ is a standard Wiener process. Writing $X$ in this way is more customary, but our assumptions mainly rely upon the ingredients coming in (2), and the process $\sigma_{t}$ which is the square-root of $c_{t}$.

Of course, having (2) is not enough, and we need assumptions on the coefficients. The assumptions on the "continuous part", that is on $b$ and $\sigma$, are as follows (for some of the results they could substantially be weakened, but in these notes it seems simpler to state a single set of assumptions):

Assumption 1. The volatility process $\sigma_{t}$ is an Itô semimartingale, that is it can be written (necessarily in a unique way) as

$$
\begin{equation*}
\sigma_{t}=\sigma_{0}+\int_{0}^{t} \widetilde{b}_{s} d s+\int_{0}^{t} \widetilde{\sigma}_{s} d W_{s}+N_{t}+\sum_{s \leq t} \Delta \sigma_{s} 1_{\left\{\left|\Delta \sigma_{s}\right|>1\right\}} \tag{3}
\end{equation*}
$$

where $N$ is a local martingale which is orthogonal to the Brownian motion $W$, and further the compensator of the process $[N, N]_{t}+\sum_{s \leq t} 1_{\left\{\left|\Delta \sigma_{s}\right|>1\right\}}$ is of the form $\int_{0}^{t} n_{s} d s$. We also assume that the processes $\widetilde{b}_{t}$ and $n_{t}$ are locally bounded (this means that there is a sequence ( $T_{n}$ ) of stopping times, increasing to infinity, and such that those processes, stopped at any time $T_{n}$, are bounded), and the processes $b_{t}$ and $\widetilde{\sigma}_{t}$ are càdlàg. Furthermore the process $\sigma$ is "non-degenerate" in the sense where $\int_{0}^{T} \sigma_{s}^{2} d s>0$ a.s. ( $T$ is the horizon of our observations).

The non-degeneracy above means that (almost) all possible paths have a non-vanishing contribution from the continuous martingale (or Wiener) part of $X$. In practice, all models
with finite activity for jumps satisfy this property, which is thus not a genuine restriction in the tests we are constructing here.

The assumptions on the jumps are quite restrictive, and in fact the same as in the paper [2] which is concerned with the estimation of the so-called Blumenthal-Getoor index. However, it is likely that these assumptions can be significantly weakened, at least for the test for which the null hypothesis is $\Omega_{T}^{f}$ :

Assumption 2. There are two constants $0 \leq \beta^{\prime}<\beta<2$ such that the Lévy measure $F_{t}=F_{t}(\omega, d x)$ is of the form

$$
\begin{equation*}
F_{t}(d x)=\frac{1}{|x|^{1+\beta}}\left(a_{t}^{(+)} 1_{\left\{0<x \leq z_{t}^{(+)}\right\}}+a_{t}^{(-)} 1_{\left\{-z_{t}^{(-)} \leq x<0\right\}}\right) d x+F_{t}^{\prime}(d x) \tag{4}
\end{equation*}
$$

where, for some locally bounded process $L_{t} \geq 1$,
(i) $a_{t}^{(+)}, a_{t}^{(-)}, z_{t}^{(+)}$and $z_{t}^{(-)}$are nonnegative predictable processes satisfying

$$
\begin{equation*}
\frac{1}{L_{t}} \leq z_{t}^{(+)} \leq 1, \quad \frac{1}{L_{t}} \leq z_{t}^{(-)} \leq 1, \quad A_{t}:=a_{t}^{(+)}+a_{t}^{(-)} \leq L_{t} \tag{5}
\end{equation*}
$$

(ii) $F_{t}^{\prime}=F_{t}^{\prime}(\omega, d x)$ is a signed measure, whose absolute value $\left|F_{t}^{\prime}\right|$ satisfies

$$
\begin{equation*}
\int\left(|x|^{\beta^{\prime}} \wedge 1\right)\left|F_{t}^{\prime}\right|(d x) \leq L_{t} \tag{6}
\end{equation*}
$$

This assumption is satisfied if the discontinuous part of $X$ is a stable process of index $\beta \in(0,2)$ : take $z_{t}^{+}=z_{t}^{-}=1$, and $a_{t}^{+}$and $a_{t}^{-}$are constants, and the residual measure $F_{t}^{\prime}$ is the restriction of the Lévy measure to the complement of $[-1,1]$, and (6) is satisfied with any $a \in(0,1)$. When the discontinuous part of $X$ is a tempered stable process the assumption is also satisfied with the same processes as above, but now the residual measure $F_{t}^{\prime}$ is not positive in general (although it again satisfies (6) with any $\beta^{\prime} \in(0, \beta)$.

This assumption also accounts for a stable or tempered stable with time varying intensity (when $z_{t}^{+}=z_{t}^{-}=1$ but $a_{t}^{+}$and $a_{t}^{-}$are genuine processes). It also accounts for any process of the form

$$
Y_{t}=Y_{0}+\int_{0}^{t} w_{s} d X_{s}
$$

as soon as $X$ satisfies the same assumption and $w_{t}$ is locally bounded and predictable.
It turns out that, under Assumption 2, the set $\Omega_{T}^{i}$ of (1) contains the set

$$
\begin{equation*}
\Omega_{T}^{i \beta}=\left\{\bar{A}_{T}>0\right\}, \quad \text { where } \bar{A}_{t}=\int_{0}^{t} A_{s} d s \tag{7}
\end{equation*}
$$

This is usually not equal to $\Omega_{T}^{i}$, though, so our test will in fact test the null hypothesis $\Omega_{T}^{f}$ against $\Omega_{T}^{i \beta}$, or vice-versa, instead of (1).

2 - Testing hypotheses which are subsets of the sample space: Here, we specify the notion of testing when the null and alternative hypotheses are families of possible outcomes. Suppose that we want to test the null hypothesis "we are in a subset $\Omega_{0}$ " of $\Omega$, against the alternative "we are in a subset $\Omega_{1}$ ", with of course $\Omega_{0} \cap \Omega_{1}=\emptyset$. We then construct a critical (rejection) region $C_{n}$ at stage $n$, that is when the time lag between observations is $\Delta_{n}$. This critical region is itself a subset of $\Omega$, which should depend only on the observed values of the process $X$ at stage $n$. We are not really within the framework of standard statistics, since the two hypotheses are themselves random.

We then take the following as our definition of the asymptotic size, for a given triple of coefficients:

$$
\begin{equation*}
a=\sup \left(\underset{n}{\lim \sup } \mathbb{P}\left(C_{n} \mid A\right): A \in \mathcal{F}, A \subset \Omega_{0}\right) \tag{8}
\end{equation*}
$$

Here $P\left(C_{n} \mid A\right)$ is the usual conditional probability knowing $A$, with the convention that it vanishes if $P(A)=0$. If $P\left(\Omega_{0}\right)=0$ then $a=0$, which is a natural convention since in this case we want to reject the assumption whatever the outcome $\omega$ is. Note that $a$ features some kind of "uniformity" over all subsets $A \subset \Omega_{0}$.

As for the asymptotic power, we define it as

$$
\begin{equation*}
P=\inf \left(\liminf _{n} \mathbb{P}\left(C_{n} \mid A\right): A \in \mathcal{F}, A \subset \Omega_{1}, \mathbb{P}(A)>0\right) \tag{9}
\end{equation*}
$$

Again, this is a number.
Before stating the results, we introduce some notation. First, the observed increments of $X$ are $\Delta_{i}^{n} X=X_{i \Delta_{n}}-X_{(i-1) \Delta_{n}}$. We take a sequence $u_{n}$ of positive numbers, which serve as thresholds and always go to 0 . There will be restrictions on this sequence, expressed by the following:

$$
\begin{equation*}
u_{n} / \Delta_{n}^{\rho} \rightarrow \infty \tag{10}
\end{equation*}
$$

for some $\rho>0$ : this condition becomes weaker when $\rho$ increases. Two specific values for $\rho$, in connection with the power $p \geq 2$ which will be used below, are of interest for us:

$$
\begin{equation*}
\rho_{1}(p)=\frac{p-2}{2 p}, \quad \rho_{2}(p)=\frac{2 p-4}{11 p-10} \tag{11}
\end{equation*}
$$

These quantities increase when $p$ increases, and $\rho_{1}(p)>\rho_{2}(p)>0$ when $p>2$. Finally, with any $p>0$ we associate the increasing processes

$$
\begin{equation*}
B\left(p, u_{n}, \Delta_{n}\right)_{t}=\sum_{i=1}^{\left[t / \Delta_{n}\right]}\left|\Delta_{i}^{n} X\right|^{p} 1_{\left\{\left|\Delta_{i}^{n} X\right| \leq u_{n}\right\}} \tag{12}
\end{equation*}
$$

consisting of the sum of the $p^{\text {th }}$ power of the increments of $X$, truncated at level $u_{n}$.

3 - The Finite Activity Null Hypothesis: We first set the null hypothesis to be finite activity, that is $\Omega_{0}=\Omega_{T}^{f}$, whereas the alternative is $\Omega_{1}=\Omega_{T}^{i}$. We choose an integer $k \geq 2$
and a real $p>2$. We consider the test statistics, which depends on $k$ and $p$, and also on the sequence $u_{n}$ going to 0 , and on the terminal time $T$, as follows:

$$
S_{n}=\frac{B\left(p, u_{n}, k \Delta_{n}\right)_{T}}{B\left(p, u_{n}, \Delta_{n}\right)_{T}}
$$

Theorem 1. a) Under Assumption 1 and if the sequence $u_{n}$ satisfies (10) with some $\rho<1 / 2$, we have

$$
\begin{equation*}
S_{n} \xrightarrow{\mathbb{P}} k^{p / 2-1} \quad \text { on the set } \Omega_{T}^{f} \tag{13}
\end{equation*}
$$

b) Under Assumptions 1 and 2, and if the sequence $u_{n}$ satisfies (10) with $\rho=\rho_{1}(p)$, we have

$$
\begin{equation*}
S_{n} \xrightarrow{\mathbb{P}} 1 \quad \text { on the set } \Omega_{T}^{i} \quad\left(\text { resp. } \Omega_{T}^{i \beta}\right) . \tag{14}
\end{equation*}
$$

To construct a reasonable test with a given level for finite samples, we need a central limit theorem associated with the convergence in (13), and a standardized version goes as follows $(\xrightarrow{\mathcal{L}-(s)}$ denotes the stable convergence in law, see for example [8] for this notion):
Theorem 2. Under Assumption 1, and if the sequence $u_{n}$ satisfies (10) with some $\rho<1 / 2$, we have

$$
\left(S_{n}-k^{p / 2-1}\right) / \sqrt{V_{n}} \xrightarrow{\mathcal{L - ( s )}} \mathcal{N}(0,1) \quad \text { in restriction to } \Omega_{T}^{f}
$$

where

$$
V_{n}=N(p, k) \frac{B\left(2 p, u_{n}, \Delta_{n}\right)_{T}}{\left(B\left(p, u_{n}, \Delta_{n}\right)_{T}\right)^{2}}
$$

and

$$
N(p, k)=\frac{1}{m_{2 p}}\left(k^{p-2}(1+k) m_{2 p}+k^{p-2}(k-1) m_{p}^{2}-2 k^{p / 2-1} m_{k, p}\right)
$$

and $m_{p, k}=\mathbb{E}\left(|U|^{p}|U+\sqrt{k-1} V|^{p}\right)$ and $m_{p}=\mathbb{E}\left(|U|^{p}\right)$ for $U$ and $V$ two independent $N(0,1)$ variables (with the notation of [1] we have $\left.N(p, k)=M(p, k) m_{p}^{2} / m_{2 p}\right)$.

We are now ready to exhibit a critical region. Denoting by $z_{a}$ the $a$-quantile of $N(0,1)$, that is $\mathbb{P}\left(U>z_{a}\right)=a$ where $U$ is $N(0,1)$, we set

$$
\begin{equation*}
C_{n}=\left\{S_{n}<k^{p / 2-1}-z_{a} \sqrt{V_{n}}\right\} \tag{15}
\end{equation*}
$$

Theorem 3. Under Assumptions 1 and 2, and if the sequence $u_{n}$ satisfies (10) with some $\rho<1 / 2$, the asymptotic level of the critical region defined by (15) for testing the null $\Omega_{T}^{f}$ against the alternative $\Omega_{T}^{i \beta}$ equals a, and the asymptotic power is 1 .

4 - The Infinite Activity Null Hypothesis: In the second case we set the null hypothesis to be infinite activity, that is $\Omega_{0}=\Omega_{T}^{i \beta}$, whereas the alternative is $\Omega_{1}=\Omega_{T}^{f}$. We choose three reals $\gamma>1$ and $p^{\prime}>p>2$. We define a family of test statistics as follows:

$$
S_{n}^{\prime}=\frac{B\left(p^{\prime}, \gamma u_{n}, \Delta_{n}\right)_{T} B\left(p, u_{n}, \Delta_{n}\right)_{T}}{B\left(p^{\prime}, u_{n}, \Delta_{n}\right)_{T} B\left(p, \gamma u_{n}, \Delta_{n}\right)_{T}}
$$

Theorem 4. Assume Assumptions 1 and 2.
a) If the sequence $u_{n}$ satisfies (10) with $\rho=\rho_{1}(p)$, we have

$$
\begin{equation*}
S_{n}^{\prime} \xrightarrow{\mathbb{P}} \gamma^{p^{\prime}-p} \quad \text { on the set } \Omega_{T}^{i \beta} \tag{16}
\end{equation*}
$$

b) If the sequence $u_{n}$ satisfies (10) with some $\rho \leq 1 / 2$, we have

$$
\begin{equation*}
S_{n}^{\prime} \xrightarrow{\mathbb{P}} 1 \quad \text { on the set } \Omega_{T}^{f} \tag{17}
\end{equation*}
$$

The associated central limit theorem is:
Theorem 5. Under Assumptions 1 and 2 with $\beta^{\prime}<\beta / 2$ and if the sequence $u_{n}$ satisfies (10) with $\rho=\rho_{2}(p)$, we have

$$
\left(S_{n}^{\prime}-\gamma^{p^{\prime}-p}\right) / \sqrt{V_{n}^{\prime}} \xrightarrow{\mathcal{L}-(s)} \mathcal{N}(0,1) \quad \text { in restriction to } \Omega_{T}^{i \beta}
$$

where

$$
\begin{aligned}
V_{n}^{\prime} & =\gamma^{2 p^{\prime}-2 p}\left(\frac{B\left(2 p, u_{n}, \Delta_{n}\right)_{T}}{\left(B\left(p, u_{n}, \Delta_{n}\right)_{T}\right)^{2}}+\left(1-2 \gamma^{-p}\right) \frac{B\left(2 p, \gamma u_{n}, \Delta_{n}\right)_{T}}{\left(B\left(p, \gamma u_{n}, \Delta_{n}\right)_{T}\right)^{2}}\right. \\
& +\frac{B\left(2 p^{\prime}, u_{n}, \Delta_{n}\right)_{T}}{\left(B\left(p^{\prime}, u_{n}, \Delta_{n}\right)_{T}\right)^{2}}+\left(1-2 \gamma^{-p^{\prime}}\right) \frac{B\left(2 p^{\prime}, \gamma u_{n}, \Delta_{n}\right)_{T}}{\left(B\left(p^{\prime}, \gamma u_{n}, \Delta_{n}\right)_{T}\right)^{2}} \\
& -2 \frac{B\left(p+p^{\prime}, u_{n}, \Delta_{n}\right)_{T}}{B\left(p, u_{n}, \Delta_{n}\right)_{T} B\left(p^{\prime}, u_{n}, \Delta_{n}\right)_{T}} \\
& \left.-2\left(1-\gamma^{-p}-\gamma^{-p^{\prime}}\right) \frac{B\left(p+p^{\prime}, \gamma u_{n}, \Delta_{n}\right)_{T}}{B\left(p, \gamma u_{n}, \Delta_{n}\right)_{T} B\left(p^{\prime}, \gamma u_{n}, \Delta_{n}\right)_{T}}\right)
\end{aligned}
$$

The critical region will therefore be

$$
\begin{equation*}
C_{n}^{\prime}=\left\{S_{n}^{\prime}<\gamma^{p^{\prime}-p}-z_{a} \sqrt{V_{n}^{\prime}}\right\} \tag{18}
\end{equation*}
$$

Theorem 6. Under Assumptions 1 and 2 with $\beta^{\prime}<\beta / 2$, and if the sequence $u_{n}$ satisfies (10) with $\rho=\rho_{2}(p)$, the asymptotic level of the critical region defined by (18) for testing the null $\Omega_{T}^{i \beta}$ against the alternative $\Omega_{T}^{f}$ equals $a$, and the asymptotic power is 1 .

Under the null hypothesis the rate of convergence is $1 / u_{n}^{\beta / 2}$ (contrary to the situation of Theorem 3, where the rate was $1 / \sqrt{\Delta_{n}}$ whatever $\beta$ and $u_{n}$ were). So, although asymptotically $u_{n}$ does not explicitly show in the test itself, one should probably take $u_{n}$ as small as possible (we have no choice as to $\beta$, of course). That is, one should take $\rho$ as large as possible, which in turn results in choosing $p$ as big as possible (recall (11)).

## References

[1] Ait Sahalia Y. and J. Jacod (2007): Testing for jumps in a discretely observed process. Annals of Statistics, forthcoming.
[2] Aït-Sahalia, Y., and J. Jacod (2008): Estimating the degree of activity of jumps in high frequency data. Annals of Statistics, forthcoming.
[3] Ball, C.A. and W.N. Torous (1983): A simplified jump process for common stock returns. J. Financial and Quantitative Analysis 18, 53-65.
[4] Bates, D.S. (1991): The crash of '87: whs it expected? The evidence form options markets. J. of Financ 46, 1009-1044.
[5] Carr, P.P., H. Geman, D.B. Madan and M. Yor (2002): The fine structure of asset returns: an empirical investigation. J. of Business 75, 305-332.
[6] Carr, P.P. and L. Wu (2003): The finite moment log stabke process and option pricing. J. of Finance 58, 753-777.
[7] Eberlein, E., U. Keller and K. Prause (1998): New insights into smile, mispricing and Value at risk. J. of Business 71, 371-405.
[8] Jacod, J. and A. Shiryaev (2003): Limit Theorems for Stochastic Processes, 2d ed., Springer-Verlag: Berlin.
[9] Madan, D.B., P.P. Carr and E.E. Chang (1998): The variance-gamma process and option pricing. European Finance Review 2, 79-105.
[10] Madan, D.B. and E. Seneta (1990): The variance-gamma (VG) model for share market returns. J. of Business 63, 511-524.
[11] Merton, R.C. (1976): Option pricing when underlying stock returns are discontinuous. J. Financial Economics 3, 125-144.


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