Ornstein-Uhlenbeck Process

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We first define several words. A stochastic process $\{Y_t : t \ge 0\}$ is

- stationary if, for all $t_1 < t_2 < \ldots < t_n$ and h > 0, the random *n*-vectors $(Y_{t_1}, Y_{t_2}, \ldots, Y_{t_n})$ and $(Y_{t_1+h}, Y_{t_2+h}, \ldots, Y_{t_n+h})$ are identically distributed; that is, time shifts leave joint probabilities unchanged
- Gaussian if, for all $t_1 < t_2 < \ldots < t_n$, the *n*-vector $(Y_{t_1}, Y_{t_2}, \ldots, Y_{t_n})$ is multivariate normally distributed
- Markovian if, for all $t_1 < t_2 < \ldots < t_n$, $P(Y_{t_n} \leq y | Y_{t_1}, Y_{t_2}, \ldots, Y_{t_{n-1}}) = P(Y_{t_n} \leq y | Y_{t_{n-1}})$; that is, the future is determined only by the present and not the past.

Also, a process $\{Y_t : t \ge 0\}$ is said to have **independent increments** if, for all $t_0 < t_1 < \ldots < t_n$, the *n* random variables $Y_{t_1} - Y_{t_0}$, $Y_{t_2} - Y_{t_1}$, ..., $Y_{t_n} - Y_{t_{n-1}}$ are independent. This condition implies that $\{Y_t : t \ge 0\}$ is Markovian, but not conversely. The increments are further said to be **stationary** if, for any t > s and h > 0, the distribution of $Y_{t+h} - Y_{s+h}$ is the same as the distribution of $Y_t - Y_s$. This additional provision is needed for the following definition.

A stochastic process $\{W_t : t \ge 0\}$ is a Wiener-Lévy process or Brownian motion if it has stationary independent increments, if W_t is normally distributed and $E(W_t) = 0$ for each t > 0, and if $W_0 = 0$. It follows immediately that $\{W_t : t > 0\}$ is Gaussian and that $Cov(W_s, W_t) = \theta^2 \min\{s, t\}$, where the variance parameter θ^2 is a positive constant. For concreteness' sake, we henceforth assume that $\theta = 1$. Almost all sample paths of Brownian motion are everywhere continuous but nowhere differentiable.

One technical stipulation is required for the following. A stochastic process $\{Y_t : t \ge 0\}$ is **continuous in probability** if, for all $u \in \mathbb{R}^+$ and $\varepsilon > 0$, $P(|Y_v - Y_u| \ge \varepsilon) \rightarrow 0$ as $v \to u$. This holds if $Cov(Y_s, Y_t)$ is continuous over $\mathbb{R}^+ \times \mathbb{R}^+$. Note that this is a statement about distributions, not sample paths.

Having dispensed with preliminaries, we turn to the central topic. A stochastic process $\{X_t : t \ge 0\}$ is an **Ornstein-Uhlenbeck process** or a **Gauss-Markov**

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process if it is stationary, Gaussian, Markovian, and continuous in probability [1, 2]. A fundamental theorem, due to Doob [3, 4, 5], ensures that $\{X_t : t \ge 0\}$ necessarily satisfies the following linear stochastic differential equation:

$$dX_t = -\rho(X_t - \mu)dt + \sigma \, dW_t$$

where $\{W_t : t \ge 0\}$ is Brownian motion with unit variance parameter and μ , ρ , σ are constants. We have moments

$$\mathbf{E}(X_t) = \mu, \qquad \operatorname{Cov}(X_s, X_t) = \frac{\sigma^2}{2\rho} e^{-\rho|s-t|}$$

in the unconditional (strictly stationary) case and

$$E(X_t \mid X_0 = c) = \mu + (c - \mu)e^{-\rho t}$$
$$Cov(X_s, X_t \mid X_0 = c) = \frac{\sigma^2}{2\rho} \left(e^{-\rho|s-t|} - e^{-\rho(s+t)} \right)$$

in the conditional (asymptotically stationary) case, where X_0 is initially constant. The latter case encompasses Brownian motion when $\mu = c = 0$, $\sigma = 1$ and $\rho \to 0^+$. The former case encompasses idealized **white noise** $\{dW_t/dt : t \ge 0\}$ when $\mu = 0$, $\sigma = \rho$ and $\rho \to \infty$.

Before proceeding, we note the following simple algorithm for generating a sample path of the Ornstein-Uhlenbeck process (also known as **colored noise**) over the time interval [0, T]. Let N be a large integer and let $z_0, z_1, ..., z_N$ be independent random numbers generated from a normal distribution with mean 0 and variance $\sigma^2/(2\rho)$. Define $x_0 = \mu + z_0$ for the unconditional case and $x_0 = c$ for the conditional case. Then define recursively

$$x_{n} = \mu + \kappa_{N}(x_{n-1} - \mu) + \sqrt{1 - \kappa_{N}^{2}} z_{n}$$

for $1 \leq n \leq N$, where $\kappa_N = \exp(-\rho T/N)$. The sequence $x_0, x_1, ..., x_N$ is called a first-order autoregressive sequence (a discrete analog of the OU process) with lag-one correlation coefficient κ_N . Finally, interpolate linearly the values $X(nT/N) = x_n$ for $0 \leq n \leq N$ to obtain the desired path [6, 7, 8]. More sophisticated simulation methods are found in [9, 10, 11].

For concreteness' sake, we henceforth assume that $\mu = 0$, $\rho = 1$ and $\sigma^2 = 2$. (Some authors take $\sigma^2 = 1$ instead; the decision becomes apparent in any paper by seeing whether $\text{Cov}(X_s, X_t)$ is $e^{-|s-t|}$ or $e^{-|s-t|}/2$.) The conditional probability

$$P(X_t \le x \mid X_0 = c) = \frac{1}{\sqrt{2\pi(1 - e^{-2t})}} \int_{-\infty}^x \exp\left(-\frac{(\xi - ce^{-t})^2}{2(1 - e^{-2t})}\right) d\xi$$

tends to the standard normal distribution, of course, as $t \to \infty$ (meaning that transients die out with time and don't affect long-term behavior). Likewise, $P(X_s \leq x)$ and $X_t \leq y | X_0 = c$) can be evaluated. One might believe that the solution of any problem involving the OU process would be similarly straightforward; the following sections serve, however, to eliminate such ideas [12, 13].

0.1. First-Passage Times. For $a \in \mathbb{R}$, we wish to find the length of time required for an OU process to cross the level x = a, given that it started at x = c. Define the first-passage time or hitting time $T_{a,c}$ by $T_{a,c} = \inf \{t \ge 0 : X_t = a \mid X_0 = c\}$. The random variable $T_{a,c}$ is 0 if and only if a = c. Let $f_{a,c}(t)$ denote the density function of $T_{a,c}$. In the special case when a = 0, it is known that [2, 12, 14, 15]

$$f_{0,c}(t) = \sqrt{\frac{2}{\pi}} \frac{|c|e^{-t}}{(1 - e^{-2t})^{3/2}} \exp\left(-\frac{c^2 e^{-2t}}{2(1 - e^{-2t})}\right)$$

but for $a \neq 0$, the formulas for $f_{a,c}(t)$ are more complicated (as we shall soon see). For a > 0 and c > 0, Thomas [16] and Ricciardi & Sato [17, 18] demonstrated that

$$E(T_{a,0}) = \sqrt{\frac{\pi}{2}} \int_{0}^{a} \left(1 + \operatorname{erf}\left(\frac{t}{\sqrt{2}}\right)\right) \exp\left(\frac{t^{2}}{2}\right) dt = \frac{1}{2} \sum_{k=1}^{\infty} \frac{\left(\sqrt{2}a\right)^{k}}{k!} \Gamma\left(\frac{k}{2}\right),$$
$$E(T_{0,c}) = \sqrt{\frac{\pi}{2}} \int_{-c}^{0} \left(1 + \operatorname{erf}\left(\frac{t}{\sqrt{2}}\right)\right) \exp\left(\frac{t^{2}}{2}\right) dt = \frac{1}{2} \sum_{k=1}^{\infty} (-1)^{k+1} \frac{\left(\sqrt{2}c\right)^{k}}{k!} \Gamma\left(\frac{k}{2}\right)$$

and, for example,

$$\begin{split} \mathbf{E}(T_{1,0}) &= 2.0934066496..., \\ \mathbf{E}(T_{2,0}) &= 10.4284093979..., \\ \end{split} \\ \mathbf{E}(T_{0,1}) &= 0.9019080126..., \\ \mathbf{E}(T_{0,2}) &= 1.4252045655.... \end{split}$$

The asymmetry in going from 0 to x, versus going from x to 0, is unsurprising: The process has mean 0, hence it tends to arrive at 0 more often than it departs from 0. For a > 0 and c > 0, we have [17, 18]

$$\begin{aligned} \operatorname{Var}(T_{a,0}) &= \sqrt{2\pi} \int_{0}^{a} \int_{-\infty}^{t} \int_{s}^{a} \left(1 + \operatorname{erf}\left(\frac{r}{\sqrt{2}}\right) \right) \exp\left(\frac{r^{2} + t^{2} - s^{2}}{2}\right) dr \, ds \, dt - \operatorname{E}(T_{a,0})^{2} \\ &= \operatorname{E}(T_{a,0})^{2} - \frac{1}{2} \sum_{k=1}^{\infty} \frac{\left(\sqrt{2a}\right)^{k}}{k!} \Gamma\left(\frac{k}{2}\right) \Psi\left(\frac{k}{2}\right), \end{aligned}$$

$$\operatorname{Var}(T_{0,c}) = \sqrt{2\pi} \int_{-c-\infty}^{0} \int_{s}^{t} \int_{0}^{0} \left(1 + \operatorname{erf}\left(\frac{r}{\sqrt{2}}\right)\right) \exp\left(\frac{r^{2} + t^{2} - s^{2}}{2}\right) dr \, ds \, dt - \operatorname{E}(T_{0,c})^{2}$$
$$= \frac{1}{2} \sum_{k=1}^{\infty} (-1)^{k} \frac{\left(\sqrt{2c}\right)^{k}}{k!} \Gamma\left(\frac{k}{2}\right) \Psi\left(\frac{k}{2}\right) - \operatorname{E}(T_{0,c})^{2}$$

where $\Psi(x) = \psi(x) - \psi(1)$ and $\psi(x)$ is the digamma function [19]. In particular, $\Psi(1) = 0$ and

$$\Psi(x) = \begin{cases} \sum_{j=1}^{x-1} \frac{1}{j} & \text{if } x \text{ is an integer} > 1\\ -2\ln(2) + 2\sum_{j=1}^{x-1/2} \frac{1}{2j-1} & \text{if } x \text{ is a half-integer} > 0. \end{cases}$$

For example,

$$Var(T_{1,0}) = 5.8420278024..., Var(T_{0,1}) = 0.8510837032..., Var(T_{2,0}) = 105.2752035488..., Var(T_{0,2}) = 1.0669454393....$$

To compute $f_{a,c}(t)$ exactly for arbitrary a and c, we would need to invert the following (Laplace transform) identity due to Darling & Siegert [20, 21, 22, 23]:

$$\mathbf{E}(e^{-\lambda T_{a,c}}) = \int_{0}^{\infty} f_{a,c}(t)e^{-\lambda t}dt = \begin{cases} \frac{D_{-\lambda}(-c)}{D_{-\lambda}(-a)}\exp\left(\frac{c^2-a^2}{4}\right) & \text{if } c < a\\ \frac{D_{-\lambda}(c)}{D_{-\lambda}(a)}\exp\left(\frac{c^2-a^2}{4}\right) & \text{if } c > a \end{cases}$$

where $D_{\nu}(x)$ is the **parabolic cylinder function** or Weber function [24]:

$$D_{\nu}(x) = \begin{cases} \sqrt{\frac{2}{\pi}} \exp\left(\frac{x^2}{4}\right) \int_{0}^{\infty} t^{\nu} \exp\left(-\frac{t^2}{2}\right) \cos\left(xt - \frac{\nu\pi}{2}\right) dt & \text{if } \nu > -1 \\ \frac{1}{\Gamma(-\nu)} \exp\left(-\frac{x^2}{4}\right) \int_{0}^{\infty} t^{-\nu-1} \exp\left(-\frac{t^2}{2} - xt\right) dt & \text{if } \nu < 0. \end{cases}$$

The two branches of this formula agree for $-1 < \nu < 0$. A differential equation

$$\frac{d^2y}{dx^2} - \left(\frac{x^2}{4} - \nu - \frac{1}{2}\right)y(x) = 0$$

is satisfied by $D_{\nu}(x)$ and, if ν is not an integer, independently by $D_{-\nu}(x)$. A series representation in terms of confluent hypergeometric functions [0.4] is also useful. Unfortunately a closed-form expression for the inverse Laplace transform seems not to be possible; only a numerical approach is feasible at present. Keilson & Ross [25] tabulated the distribution of $T_{a,c}$ for a number of values a and c. For example, the median time for an OU process X_t to reach a = 1, given that $X_0 = c = 0$, is 1.1892.... This corresponds to the 50th percentile of the distribution of $T_{1,0}$. The median of $T_{2,0}$, by contrast, is 7.2521....

We turn to a more complicated problem involving two (absorbing) boundaries rather than just one. Given a < c < b, what is the length of time required for the process to escape the interval (a, b), given that it started at x = c? Define $T_{a,b,c} = \inf \{t \ge 0 : X_t = a \text{ or } X_t = b | X_0 = c\}$ and let $f_{a,b,c}(t)$ denote the density function of $T_{a,b,c}$. Efforts have focused on the scenario in which -a = b > 0. The Laplace transform of $f_{-b,b,c}(t)$ satisfies [22]

$$E(e^{-\lambda T_{-b,b,c}}) = \frac{D_{-\lambda}(c) + D_{-\lambda}(-c)}{D_{-\lambda}(b) + D_{-\lambda}(-b)} \exp\left(\frac{c^2 - b^2}{4}\right)$$

assuming -b < c < b. From another table in [25], the median of $T_{-1,1,0}$ is found to be 0.4449.... The reason that this is less than 1.1892... is clear: Each direction of travel leads to a potential crossing. The median of $T_{-2,2,0}$ is 3.2439....

Keilson & Ross' approach to evaluating such probabilities was based on finding zeroes and residues in the complex plane of the parabolic cylinder functions. Alternative approaches for numerically computing $f_{a,c}(t)$ and $f_{-b,b,c}(t)$ include [26, 27, 28, 29]. We report on some related asymptotics in [0.4].

There is an obvious connection between first-passage times and extreme values of a process (in the conditional case). We simply summarize:

$$\left. \begin{array}{l} \Pr\left(\max_{0 \le t \le T} X_t \le a \middle| X_0 = c \right) & \text{if } c < a \\ \Pr\left(\min_{0 \le t \le T} X_t \ge a \middle| X_0 = c \right) & \text{if } c > a \end{array} \right\} = \Pr(T_{a,c} > T) = 1 - F_{a,c}(T)$$

and, if a < c < b,

$$P\left(\left.a \le \min_{0 \le t \le T} X_t \le \max_{0 \le t \le T} X_t \le b \right| X_0 = c\right) = P(T_{a,b,c} > T) = 1 - F_{a,b,c}(T)$$

where $F_{a,c}(t)$, $F_{a,b,c}(t)$ are the cumulative distribution functions of $T_{a,c}$, $T_{a,b,c}$. In the special case when -a = b > 0, the latter formula becomes a statement about $\max_{0 \le t \le T} |X_t|$, given $X_0 = c$. Also, the **range** of the process satisfies [22]

$$P\left(\max_{0 \le t \le T} X_t - \min_{0 \le t \le T} X_t \le r \middle| X_0 = c\right) = \int_0^r \int_{c-q}^c \left. \frac{\partial^2}{\partial a \,\partial b} F_{a,b,c}(T) \right|_{b=a+q} \, da \, dq$$

but no one apparently has calculated this probability.

0.2. Historical Maximums. If the condition $X_0 = c$ is discarded, what then can be said about $\max_{0 \le t \le T} X_t$ or $\max_{0 \le t \le T} |X_t|$? We focus solely on the former expression and write $M_T = \max_{0 \le t \le T} X_t$. It can be shown that [30, 31, 32]

$$P(M_T \le 0) = \frac{1}{\pi} \arcsin\left(e^{-T}\right)$$

which is a beautiful (but isolated) result. More generally [32],

$$\int_{0}^{\infty} \mathcal{P}(M_t \le y) e^{-\lambda t} dt = \frac{1}{\sqrt{2\pi}} \int_{-\infty}^{y} \frac{1}{\lambda} \left(1 - \frac{D_{-\lambda}(-x)}{D_{-\lambda}(-y)} \exp\left(\frac{x^2 - y^2}{4}\right) \right) \exp\left(-\frac{x^2}{2}\right) dx$$

for arbitrary y, or

$$\int_{0}^{\infty} g_t(y) e^{-\lambda t} dt = \frac{1}{\sqrt{2\pi}} \frac{D_{-\lambda-1}(-y)^2}{D_{-\lambda}(-y)^2} \exp\left(\frac{-y^2}{2}\right)$$

where $g_t(y)$ is the density function of M_t . For example, the median value of M_1 is 1.0393... and the median value of M_{10} is 2.2202.... It can be inferred from [0.3] that the median of M_T is asymptotically $\sqrt{2\ln(T)}$ as $T \to \infty$.

An alternative approach for numerically computing $P(M_t \leq y)$ via the Mellin transform is due to DeLong [33, 34, 35]. We hope to report on this later. An interesting application to computer science, involving the maximum size reached by a dynamic data structure over a long span of time, is described in [36].

0.3. Pickands' Constants. Assume that $\{Y_t : t \ge 0\}$ is a stationary Gaussian process with zero mean, unit variance and covariance function of the form

$$r(|s-t|) = Cov(Y_s, Y_t) = 1 - C |s-t|^{\alpha} + o(|s-t|^{\alpha})$$

as $|s - t| \to 0$, where $0 < \alpha \leq 2$ and C > 0 are constants. Assume further that $r(\tau) \ln(\tau) \to 0$ as $\tau \to \infty$. Pickands [37, 38, 39, 40, 41] demonstrated that $M_T = \max_{0 \leq t \leq T} Y_t$ has the Gumbel limiting distribution [42]

$$\lim_{T \to \infty} \mathbb{P}\left(\sqrt{2\ln(T)} \left(M_T - k_T\right) \le x\right) = \exp(-e^{-x})$$

where

$$k_T = \sqrt{2\ln(T)} + \frac{1}{\sqrt{2\ln(T)}} \left\{ \frac{2-\alpha}{2\alpha} \ln(\ln(T)) + \ln\left((2\pi)^{-\frac{1}{2}} 2^{\frac{2-\alpha}{2\alpha}} C^{\frac{1}{\alpha}} H_{\alpha} \right) \right\}$$

and H_{α} is a positive constant independent of C. It is known that $H_1 = 1$ (corresponding to the OU process) and $H_2 = 1/\sqrt{\pi}$. No other exact values for H_{α} are known. An alternative characterization of H_{α} is

$$H_{\alpha} = \lim_{T \to \infty} \int_{0}^{\infty} \mathcal{P}(\tilde{M}_{T} > y) e^{y} dy$$

where $\{\tilde{Y}_t : t \ge 0\}$ is a nonstationary Gaussian process with

$$\mathbf{E}(\tilde{Y}_t) = -|t|^{\alpha}, \quad \operatorname{Cov}(\tilde{Y}_s, \tilde{Y}_t) = |s|^{\alpha} + |t|^{\alpha} - |s - t|^{\alpha}$$

but this does not seem to help. Shao [43] and Debicki, Michna & Rolski [44] gave bounds on H_{α} ; for example,

$$0.009 \le H_{1/2} \le 715.94, \quad 0.208 \le H_{3/2} \le 3.04.$$

A conjecture that $H_{\alpha} = 1/\Gamma(1/\alpha)$ remains unproved. There is also a connection with the Gaussian correlation conjecture and with estimating small ball probabilities [45], topics which we hope to address later.

0.4. Upper Tail Asymptotics. We revisit the single-boundary first-passage time distribution and ask about the limiting value

$$\lambda(a) = \lim_{t \to \infty} \frac{1}{t} \ln \left\{ \Pr\left(T_{a,0} > t\right) \right\}$$

as a function of a > 0. In words, what can be said about the upper tail of the distribution of the first hitting time $T_{a,0}$ for an OU process X_t across the level x = a, given that $X_0 = 0$? Mandl [46, 47] and Beekman [48] demonstrated that $-1 < \lambda(a) < 0$ and that $\lambda(a)$ is the zero of $D_{-\lambda}(-a)$ closest to 0. Sample values include [17, 49, 50]

$$\lim_{a \to 0^+} \lambda(a) = -1, \qquad \lim_{a \to \infty} \lambda(a) \cdot \frac{\exp(a^2/2)}{a} = \frac{-1}{\sqrt{2\pi}},$$
$$\lambda(0.7649508673...) = -\frac{1}{2},$$
$$\lambda(1) = -0.3882382947... = 2(-0.1941191473...),$$
$$\lambda(2) = -0.0972745958... = 2(-0.0486372979...).$$

For the symmetric double-boundary first-passage time distribution, we examine

$$\lambda(-b,b) = \lim_{t \to \infty} \frac{1}{t} \ln \left\{ \Pr\left(T_{-b,b,0} > t\right) \right\}$$

as a function of b > 0. Breiman [51] proved that $-\infty < \lambda(-b, b) < 0$ and that $\lambda(-b, b)$ is the zero of $\Phi(\lambda/2, 1/2, b^2/2)$ closest to 0, where

$$\Phi(u, v, w) = 1 + \sum_{k=1}^{\infty} \frac{u(u+1)(u+2)\cdots(u+k-1)}{v(v+1)(v+2)\cdots(v+k-1)} \frac{w^k}{k!}$$

is the confluent hypergeometric function of the first kind. For simplicity, define $\mu(b) = \lambda(-b, b)$. Sample values include [50, 51, 52]

$$\begin{split} \lim_{b \to 0^+} \mu(b) &= -\infty, \quad \lim_{b \to \infty} \mu(b) \cdot \frac{\exp(b^2/2)}{b} = \frac{-1}{\sqrt{2\pi}}, \\ \mu(1) &= -2, \quad \mu(1.3069297277...) = -1, \quad \mu(1.6438001904...) = -\frac{1}{2}, \\ \mu\left(\sqrt{3} - \sqrt{6}\right) &= \mu(0.7419637843...) = -4, \\ \mu(2) &= -0.2429928807..., \quad \mu(3) = -0.0239463006..., \\ \mu\left(\sqrt{2}\right) &= -0.7984598320..., \quad \mu\left(2\sqrt{2}\right) = -0.0374612092.... \end{split}$$

The latter two values come from [52], where a different time scaling was chosen. Also, the constant $(3 - 6^{1/2})^{1/2}$ appears in [53, 54, 55] with regard to stopping rules in statistical sequential analysis.

For completeness' sake, here is the expression for $D_{-\lambda}(x)$ in terms of confluent hypergeometric functions:

$$D_{-\lambda}(x) = \frac{\sqrt{\pi}2^{-\lambda/2}}{\Gamma((1+\lambda)/2)} e^{-x^2/4} \Phi\left(\frac{\lambda}{2}, \frac{1}{2}, \frac{x^2}{2}\right) - 2\frac{\sqrt{\pi}2^{-(1+\lambda)/2}}{\Gamma(\lambda/2)} x e^{-x^2/4} \Phi\left(\frac{1+\lambda}{2}, \frac{3}{2}, \frac{x^2}{2}\right)$$

which gives rise to the values $\lambda(1)$, $\lambda(2)$ and $\lambda^{-1}(-1/2)$ listed earlier. The constant $\mu^{-1}(-1)$ is important in the study of sample path behavior of Brownian motion [50, 56, 57] and first appeared in [54], as far as is known. Some higher dimensional results are given in [50, 58]. Csáki [59, 60] recently outlined the distributional asymptotics of the maximum M_T , but we cannot discuss this topic further.

0.5. Addendum. New numerical transform inversion algorithms [61, 62, 63] make enhancement of the tables in [25, 32] possible. Also, the distribution of the L_2 -norm of X_t on [0, T] can be inferred from closed-form expressions in [64, 65]. We wonder about corresponding results for L_1 and L_{∞} -norms. The conjectured formula for H_{α} in terms of the gamma function is probably false [66, 67, 68, 69]; simulation-based point estimates $H_{3/2} \approx 0.77$ and confidence bounds $0.768 \leq H_{3/2} \leq 0.786$ do not carry over well to $H_{1/2}$ since the underlying algorithm becomes unreliable for $0 < \alpha < 1$.

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