# Introduction to stochastic numerics and computing ergodic limits

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# Plan of the talk

- What's for?
- Stochastics
  - SDE
  - Link to PDEs
- Numerics
  - Types of convergence
  - Mean-square approximations
  - Weak approximations
- Computing ergodic limits
  - Ensemble and time averaging and computational errors
  - Gradient system and 'non-Markovian' scheme
- Conclusions

- Molecular dynamics (MD) canonical ensemble (NVT) see Part II on Friday
- Financial Engineering (pricing and hedging derivatives; management of couterparty risk, etc.)
- Bayesian statistics
- SDEs  $\iff$  linear parabolic and elliptic PDEs
- Other (models in biology, chemistry, economics, etc.)

# Stochastics: SDEs

Assume that deterministic functions b(t, x) and  $b_r(t, x)$ , r = 1, ..., q, have some good analytical properties. The solution of the (Ito) **stochastic differential equation** 

$$dX = b(t, X)dt + \sum_{r=1}^{q} \sigma_r(t, X)dw_r(t), \quad X(t_0) = x,$$
(1)

is a stochastic process X(t) such that

$$X(t) = x + \int_{t_0}^t b(s, X(s)) ds + \sum_{r=1}^q \int_{t_0}^t \sigma_r(s, X(s)) dw_r(s)$$
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$$\int_{0}^{T} \sigma_r(t, X(t)) dw_r(t) = \lim_{\delta \to 0} \sum_{i=0}^{N-1} \sigma_r(t_i, X(t_i)) \times (w_r(t_{i+1}) - w_r(t_i)),$$

$$\int_{0}^{t} \sigma_{r}(t, X(t)) \circ dw_{r}(t) = \lim_{\delta \to 0} \sum_{i=0}^{N-1} \sigma_{r}\left(t_{i}, \frac{X(t_{i}) + X(t_{i+1})}{2}\right) (w_{r}(t_{i+1}) - w_{r}(t_{i}))$$

# Feynman-Kac formula

$$Lu(t,x) := \frac{\partial}{\partial t}u(t,x) + \frac{1}{2}\sum_{r=1}^{q}\sum_{i,j=1}^{d}\sigma_{r}^{i}\sigma_{r}^{j}\frac{\partial^{2}}{\partial x^{i}\partial x^{j}}u(t,x) \qquad (3)$$
$$+\sum_{i=1}^{d}b^{i}\frac{\partial}{\partial x^{i}}u(t,x)$$

The Cauchy problem for linear parabolic PDE:

$$Lu = 0, \quad t < T, \ x \in R^{d},$$
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$$u(T,x) = f(x), x \in \mathbb{R}^d.$$
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Then

$$u(t_0, x) = Ef(X_{t_0, x}(T)),$$
(6)

where  $X_{t_0,x}(t)$ ,  $t \ge t_0$ , is the solution of the Ito SDEs

$$dX = b(t, X)dt + \sum_{r=1}^{q} \sigma_r(t, X)dw_r(t), \quad X(t_0) = x.$$
(7)

#### Numerics

$$dX = b(t, X)dt + \sum_{r=1}^{q} \sigma_r(t, X)dw_r(t), \ X(0) = x,$$
(8)

Consider a numerical method for (8) based on the one-step approximation:

$$X_{t,x}(t+h) \simeq \bar{X}_{t,x}(t+h) = x + A(t,x,h;\xi), \ 0 \le t < t+h \le T, \quad (9)$$

where  $\xi$  is a random vector with moments of a sufficiently high order and A is a d-dimensional vector function. Introduce (for simplicity) the equidistant partition of the time interval [0, T] into N parts with the step h = T/N:  $0 = t_0 < t_1 < \cdots < t_N = T$ ,  $t_{k+1} - t_k = h$ .

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$$X_0 = x, \ X_{k+1} = X_k + A(t_k, X_k, h; \xi_{k+1}), \ k = 0, \dots, N-1,$$
 (10)

where  $\xi_1$  is independent of  $X_0$  and  $\xi_{k+1}$  for k > 0 is independent of  $X_0, \ldots, X_k, \xi_1, \ldots, \xi_k$ .

# Numerics - Example

Euler's method:

$$X_{k+1} = X_k + b_k h + \sum_{r=1}^{q} \sigma_{rk} \Delta_k w_r(h),$$
 (11)

where  $\Delta_k w_r(h) = w_r(t_{k+1}) - w_r(t_k)$ , and the index k at  $\sigma_r$  and b indicates that these functions are evaluated at the point  $(t_k, X_k)$ .

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$$(E(X(t_k) - X_k)^2)^{1/2} \le Ch^{1/2}, \qquad (12)$$

where C is a constant not depending on k and h.

# Numerics - mean-square methods

We usually distinguish two types of convergence of numerical methods for SDEs: mean-square (also called strong) and weak. Mean-square methods are used for direct simulation of SDEs' trajectories which, e.g., can give information on general behavior of a stochastic model.

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Definition

If for some method we would have

$$(E(X(t_k) - X_k)^2)^{1/2} \le Ch^p, \qquad (13)$$

then we say that **the mean-square order of accuracy** of the method is *p*.

[e.g. Milstein, T.; Springer, 2004]

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If an approximation  $\bar{X}$  is such that

$$|Ef(\bar{X}(T)) - Ef(X(T))| \le Kh^p$$
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for f from a sufficiently large class of functions, then we say that the weak order of accuracy of the approximation  $\bar{X}$  (the method  $\bar{X}$ ) is p.

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If a method converges with an order p in the mean-square sense, it also converges in the weak sense with order equal to or larger than p. The opposite is not true.

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The weak order of accuracy of Euler's method (11) is 1.

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A crude method like

$$X_{k+1} = X_k + b_k h + \sqrt{h} \sum_{r=1}^{q} \sigma_{rk} \eta_{rk} , \qquad (15)$$

where  $\eta_{rk}$ ,  $r = 1, \ldots, q$ ,  $k = 0, \ldots, N - 1$ , are independent random variables taking the values +1 and -1 with probabilities 1/2, also has first order of accuracy in the sense of weak approximation. We usually call (15) as the *weak Euler scheme* (Milstein (1978); see also Talay (1984), Milstein (1985) and Milstein&T, Springer, 2004).

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The main interest in weak approximations lies in the hope to obtain simpler methods and, in particular, methods not requiring modeling of complicated random variables.

# Numerics - Monte Carlo

The Monte Carlo technique:

$$u \equiv Ef(X(T)) \simeq \bar{u} \equiv Ef(X_N) \simeq \hat{u} \equiv \frac{1}{M} \sum_{m=1}^M f(X_N^{(m)}) , \qquad (16)$$

where  $X_N^{(m)}$ , m = 1, ..., M, are independent realizations of the random variable  $X_N$ .

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where  $X_N^{(m)}$ , m = 1, ..., M, are independent realizations of the random variable  $X_N$ .

The error of the Monte Carlo method in (16) is evaluated by

$$\bar{\sigma} = c \, \frac{\left[ Var\left\{ f(X_N) \right\} \right]^{1/2}}{M^{1/2}}$$

In practice:

$$\bar{u} \in (\hat{u} - \frac{c}{\sqrt{M}}\sqrt{\hat{v}}, \hat{u} + \frac{c}{\sqrt{M}}\sqrt{\hat{v}}), \qquad (17)$$
$$\hat{v} \equiv \frac{1}{M} \sum_{m=1}^{M} [f(_m X_N)]^2 - \hat{u}^2,$$

with probability 0.68 for c = 1, 0.95 for c = 2, and 0.997 for c = 3. [e.g. Milstein, T.; Springer, 2004]

# Ergodic limits

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$$dX = b(X)dt + \sum_{l=1}^{q} \sigma_{r}(X)dw_{r}(t), \ X(0) = x$$
(18)

We assume

- (A1) The coefficients of (18) are sufficiently smooth functions in  $\mathbf{R}^{d}$ .
- (A2) The solution of (18) is *regular*, i.e., it is defined for all  $t \ge 0$ .
- (A3) The process X(t) is *ergodic*, i.e., there exists a unique invariant measure  $\mu$  of X and independently of  $x \in \mathbf{R}^d$  there exists the limit

$$\lim_{t \to \infty} E\varphi(X(t;x)) = \int \varphi(x) \, d\mu(x) := \varphi^{\text{erg}}$$
(19)

for any function  $\varphi(x)$  with polynomial growth at infinity.

(A4) The Markov transition function P(t, x, dy) and the invariant measure  $\mu(dx)$  have sufficiently smooth densities p(t, x, y) and  $\rho(x)$ , respectively.

[Hasminskii 1980; 2012]

# Ergodic limits

**Ergodicity.** We are interested here in systems which solutions satisfy a stronger condition than (A3):

(A3e) The process X(t) is *exponentially ergodic*, i.e., for any  $x \in \mathbf{R}^d$  and any function  $\varphi$  with a polynomial growth we have the following strengthening of (19):

$$|E\varphi(X(t;x)) - \varphi^{erg}| \le Ce^{-\lambda t}, \quad t \ge 0,$$
(20)

where C > 0 and  $\lambda > 0$  are some constants.

It follows from (20) (and (19)) that for any  $\varepsilon>0$  there exists  $T_0>0$  such that for all  $T\geq T_0$ 

$$|E\varphi(X(T;x)) - \varphi^{erg}| \le \varepsilon.$$
(21)

Thus, the problem of computing the ergodic limit is reduced to evaluating the expectation  $E\varphi(X(T))$  at a finite (though usually large) time T.

# Example: Langevin equations

#### Example

The Langevin equations are very popular in, e.g., molecular simulation:

$$dP = f(Q)dt - \gamma Pdt + \sqrt{\frac{2\gamma}{\beta}}dw(t), \ P(0) = P_0 = p, \quad (22)$$
  
$$dQ = M^{-1}Pdt, \ Q(0) = Q_0 = q,$$

where P, Q,  $f(q) = -\nabla U(q)$  are *n*-dimensional column-vectors,  $\gamma > 0$  is a friction parameter, M is a non-singular, diagonal matrix, and  $w(t) = (w_1(t), \ldots, w_n(t))^{\top}$  and  $w_l(t)$  are independent standard Wiener processes.

Let  $U(q) \ge 0$  for all  $q \in \mathbf{R}^n$  and there exists an  $\alpha_1 > 0$  and  $0 < \alpha_2 < 1$ such that  $\frac{1}{2}(\nabla U(q), q) \ge \alpha_2 U(q) + \gamma^2 \frac{\alpha_2(2-\alpha_2)}{8(1-\alpha_2)}|q|^2 - \alpha_1$ . Then the solution of (22) is exponentially ergodic [Mattingly, Stuart, Higham (2002)] and its invariant measure is Gibbsian with the density

$$\rho(\mathbf{p}, \mathbf{q}) \propto \exp\left(-\beta \left\{\frac{1}{2}\mathbf{p}^{\top} \mathbf{M}^{-1} \mathbf{p} + U(\mathbf{q})\right\}\right)$$
(23)

# Ensemble averaging

$$|E\varphi(X(T;x)) - \varphi^{erg}| \leq \varepsilon.$$

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(24)

Monte Carlo estimate for the ergodic limit  $\varphi^{erg}$ :

$$\hat{\varphi}^{erg} = \frac{1}{M} \sum_{m=1}^{M} \varphi\left(\bar{X}^{(m)}(T; x)\right), \qquad (25)$$

where *M* is the number of independent approximate realizations and  $\bar{X}(T; x)$  is a weak approximation of X(T; x) with order p > 0.

The total error

$$R_{\hat{\varphi}^{erg}} := \hat{\varphi}^{erg} - \varphi^{erg} \tag{26}$$

consists of three parts: the error  $\varepsilon$  of the ergodic limit approximation; the error of numerical integration  $Kh^p$ ; and the Monte Carlo error; i.e.,

$${\cal R}_{\hatarphi^{
m erg}}\sim {\cal K} h^{
m p}+arepsilon+{\cal O}\left(rac{1}{\sqrt{M}}
ight).$$

[Milstein, T.; Physica D (2007)]

# Ensemble averaging

More precisely,

$$Bias(\hat{\varphi}^{erg}) = |E\hat{\varphi}^{erg} - \varphi^{erg}| \le Kh^p + \varepsilon$$
(27)

and the estimator's variance is

$$Var(\hat{\varphi}^{erg}) = O(1/M).$$
 (28)

# Time averaging

$$\lim_{t\to\infty}\frac{1}{t}\int_{0}^{t}\varphi(X(s;x))ds=\varphi^{erg} \quad a.s., \tag{29}$$

where the limit does not depend on x.

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Then by approximating a single trajectory, one gets for a sufficiently large  $\tilde{\mathcal{T}}$  :

$$\frac{1}{\tilde{T}}\int_{0}^{\tilde{T}}\varphi(X(s;x))ds\approx\check{\varphi}^{erg}=\check{\varphi}_{L}^{erg}:=\frac{1}{L}\sum_{l=1}^{L}\varphi(\bar{X}(t_{l};x)),\qquad(30)$$

where  $Lh = \tilde{T}$ . [Talay (1990); Talay, Tubaro (1990); Milstein, T (2007); Mattingly, Stuart, T. (2010)] Under some conditions on the SDE coefficients:

$$\operatorname{Bias}(\check{\varphi}_{L}^{\operatorname{erg}}) \leq K_{1}h^{p} + \frac{K_{2}}{\tilde{T}}; \qquad (31)$$

$$\operatorname{Var}(\check{\varphi}_{L}^{erg}) \leq \frac{K}{\tilde{T}};$$
 (32)

there exists a deterministic constant K so that for h sufficiently small, positive  $\varepsilon > 0$ , and  $\tilde{T}$  sufficiently large one has:

$$|\check{\varphi}_{L}^{erg} - \varphi^{erg}| \le Kh^{p} + \frac{C(\omega)}{\tilde{T}^{1/2-\varepsilon}}$$
 a.s., (33)

where  $C(\omega) > 0$  is an a.s. bounded random variable depending on  $\varepsilon$  and the particular  $\varphi$ .

There are three types of errors:

- 1 numerical integration error (estimated by Kh<sup>p</sup>);
- 2 the error due to the distance from the stationary distribution (i.e., the error due to the finite time of integration T estimated by K/T);

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- 1 numerical integration error (estimated by Kh<sup>p</sup>);
- 2 the error due to the distance from the stationary distribution (i.e., the error due to the finite time of integration T estimated by K/T);
- 3 the statistical error (  $\sim 1/\sqrt{ ilde{T}}$  ).

[Talay (1990); Mattingly, Stuart, T. (2010)]

# Time averaging: stats error

Run a long trajectory  $M\tilde{T}$  split into M blocks of a large length  $\tilde{T} = hL$  each. We evaluate the estimators  $_{m}\check{\varphi}_{L}^{erg}$ ,  $m = 1, \ldots, M$ , for each block. Since  $\tilde{T}$  is big and a time decay of correlations is usually fast,  $_{m}\check{\varphi}_{L}^{erg}$  can be considered as almost uncorrelated. We compute the sampled variance

$$\hat{D} = \frac{1}{M-1} \sum_{m=1}^{M} \left( {}_{m} \check{\varphi}_{L}^{erg} \right)^{2} - \left( \frac{1}{M} \sum_{m=1m}^{M} \check{\varphi}_{L}^{erg} \right)^{2}$$

For a sufficiently large  $\tilde{T}$  and M,  $\mathbf{E} \check{\varphi}_L^{erg}$  belongs to the confidence interval

$$\mathbf{E}\,\check{arphi}_{L}^{ ext{erg}} \in ig(\check{arphi}_{LM}^{ ext{erg}} - crac{\sqrt{\hat{D}}}{\sqrt{M}},\check{arphi}_{LM}^{ ext{erg}} + crac{\sqrt{\hat{D}}}{\sqrt{M}}ig) \;,$$

with probability, for example 0.95 for c = 2 and 0.997 for c = 3. Note that **E**  $\check{\varphi}_{I}^{erg}$  contains the two errors forming the bias.

[Leimkuhler, Matthews, T.; Proc. R. Soc. A (2014)]

$$d\mathbf{X} = \mathbf{a}(\mathbf{X})dt + \sigma d\mathbf{w}, \ \mathbf{X}(\mathbf{0}) = \mathbf{X}_{\mathbf{0}}, \tag{34}$$

$$a(x) := -\nabla U(x), \quad x \in \mathbf{R}^d, \tag{35}$$

 $\sigma = \sqrt{2/\beta}$ , d = 3n, and w(t) is a standard *d*-dimensional Wiener process.

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We use the following notation for the solution of (34):  $X(t) = X_{t_0,x}(t)$ when  $X(t_0) = x$ ,  $t \ge t_0$ , and also we will write  $X_x(t)$  when  $t_0 = 0$ .

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The solution X(t) of (34) is exponentially ergodic and

$$\rho(x) \propto \exp\left(-\frac{2}{\sigma^2}U(x)\right)$$

if there exist  $c_0 \in \mathbf{R}$  and  $c_1 > 0$  such that

$$(x, a(x)) \le c_0 - c_1 |x|^2.$$
 (36)

[Hasminskii (1980), Mattingly, Stuart, Higham (2002)]

The Euler scheme:

$$\mathbf{X}_{k+1} = \mathbf{X}_k + h\mathbf{a}(\mathbf{X}_k) + \sigma\sqrt{h}\xi_{k+1},$$
(37)

where  $\xi_k = (\xi_k^1, \dots, \xi_k^d)^\top$  and  $\xi_k^i$ ,  $i = 1, \dots, d$ ,  $k = 1, \dots$ , are i.i.d. random variables with the law  $\mathcal{N}(0, 1)$ .

The Euler scheme:

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Heun's scheme:

$$\hat{X}_{k+1} = X_k + ha(X_k) + \sigma\sqrt{h}\xi_{k+1},$$

$$X_{k+1} = X_k + \frac{h}{2} \left[ a(\hat{X}_{k+1}) + a(X_k) \right] + \sigma\sqrt{h}\xi_{k+1}.$$
(38)

The Euler scheme:

$$\mathbf{X}_{k+1} = \mathbf{X}_k + ha(\mathbf{X}_k) + \sigma \sqrt{h} \xi_{k+1},$$
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where  $\xi_k = (\xi_k^1, \dots, \xi_k^d)^\top$  and  $\xi_k^i$ ,  $i = 1, \dots, d$ ,  $k = 1, \dots$ , are i.i.d. random variables with the law  $\mathcal{N}(0, 1)$ .

Heun's scheme:

$$\hat{X}_{k+1} = X_k + ha(X_k) + \sigma\sqrt{h}\xi_{k+1},$$

$$X_{k+1} = X_k + \frac{h}{2} \left[ a(\hat{X}_{k+1}) + a(X_k) \right] + \sigma\sqrt{h}\xi_{k+1}.$$
(38)

Weak convergence:

$$|E\varphi(X(T)) - E\varphi(X_N)| \le Kh^p \tag{39}$$

N = T/h; Euler – p = 1, Heun – p = 2 [e.g. Milstein, T.; Springer 2004]

The Euler scheme:

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Ergodic limits [Talay (1990); Talay, Tubaro (1990); Mattingly, Stuart, Higham (2002); Milstein, T (2007); Mattingly, Stuart, T. (2010)]:

$$|\varphi^{\text{erg}} - E\varphi(X_N)| \le Kh^p + Ce^{-\lambda T}$$
(40)

#### Non-Markovian scheme

$$X_{k+1} = X_k + ha(X_k) + \sigma \frac{\sqrt{h}}{2}(\xi_k + \xi_{k+1}),$$
 (41)

where  $\xi_k = (\xi_k^1, \dots, \xi_k^i)^\top$  defined on  $(\Omega, P, \mathcal{F})$  and  $\xi_k^i$ ,  $i = 1, \dots, d$ ,  $k = 1, \dots$ , are i.i.d. random variables with the law  $\mathcal{N}(0, 1)$  [Leimkuhler, Matthews (2013)]

## Example

Let  $a(x) = -\alpha x$  with  $\alpha > 0$ , then X(t) from (34) is the Ornstein-Uhlenbeck process, which is Gaussian with  $\mathbf{E}X_x(t) = xe^{-\alpha t}$ ,  $Cov(X_x(s), X_x(t)) = \frac{\sigma^2}{2\alpha}(e^{-\alpha(t-s)} - e^{-\alpha(t+s)})$  for  $s \le t$  and  $Var(X_x(t)) = \frac{\sigma^2}{2\alpha}(1 - e^{-2\alpha t})$ .

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$$\begin{aligned} \mathbf{E}\mathbf{X}_{N} &= x_{0}(1-\alpha h)^{N} = x_{0}e^{-\alpha T}(1+\mathcal{O}(h)),\\ Var(\mathbf{X}_{N}) &= \frac{\sigma^{2}}{2\alpha}\frac{1-(1-\alpha h)^{2N}}{1+\alpha h}\\ &= \frac{\sigma^{2}}{2\alpha}(1-e^{-2\alpha T})-\frac{\sigma^{2}}{2}h+e^{-2\alpha T}\mathcal{O}(h)+\mathcal{O}(h^{2}), \ \alpha h<1, \end{aligned}$$

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**Assumption 1** There exist  $c_0 \in \mathbf{R}$  and  $c_1 > 0$  such that

$$(x, a(x)) \leq c_0 - c_1 |x|^2.$$

**Assumption 2** The potential  $U(x) \in C^7(\mathbf{R}^d)$ , its first-order derivatives grow not faster than a linear function at infinity and higher derivatives are bounded. The function  $\varphi(x) \in C^6(\mathbf{R}^d)$  and it and its derivatives grow not faster than a polynomial function at infinity.

The most restrictive condition in Assumption 2 is the requirement for  $a(x) = -\nabla U$  to be globally Lipschitz:

$$|a(x)|^2 \le K(1+|x|^2),$$
 (42)

where K > 0 is independent of  $x \in \mathbf{R}^d$ , which can be relaxed via [Milstein, T. (2005) and (2007): the concept of rejecting exploding trajectories]

Introduce the operator L

$$L:=\frac{\partial}{\partial t}+\mathcal{L},$$

where  ${\cal L}$ 

$$\mathcal{L} := \sum_{i=1}^{d} a^{i}(x) \frac{\partial}{\partial x^{i}} + \frac{\sigma^{2}}{2} \sum_{i=1}^{d} \frac{\partial^{2}}{\left(\partial x^{i}\right)^{2}}.$$
(43)

The function

$$u(t,x) = \mathbf{E}\varphi(\mathbf{X}_{t,x}(T))$$
(44)

satisfies the Cauchy problem for the backward Kolmogorov equation

$$Lu = 0, \qquad (45)$$
$$u(T, x) = \varphi(x).$$

#### Main theorem

Theorem (1; [Leimkuhler, Matthews, T.; Proc. R. Soc. A (2014)]) Let Assumptions 1-2 hold. Then the scheme (41) is first order weakly convergent and for all sufficiently small h > 0 its error has the form

$$\mathbf{E}\varphi(\mathbf{X}_{x}(T)) - \mathbf{E}\varphi(\mathbf{X}_{N}) = C_{0}(T, x)h + C(T, x)h^{2},$$
(46)

$$C_0(T,x) = \mathbf{E} \int_0^T B_0(t, \mathbf{X}_x(t)) \mathrm{d}t, \qquad (47)$$

$$B_{0}(t,x) = \frac{1}{2} \left[ \sum_{i,j=1}^{d} a^{j}(x) \frac{\partial a^{i}(x)}{\partial x^{j}} \frac{\partial u(t,x)}{\partial x^{i}} + \frac{\sigma^{2}}{2} \sum_{i,j=1}^{d} \frac{\partial a^{i}(x)}{\partial x^{j}} \frac{\partial^{2} u(t,x)}{\partial x^{i} \partial x^{j}} + \frac{\sigma^{2}}{2} \sum_{i,j=1}^{d} \frac{\partial^{2} a^{i}(x)}{(\partial x^{j})^{2}} \frac{\partial u(t,x)}{\partial x^{i}} \right],$$
$$|C(T,x)| \leq K(1+|x|^{\varkappa}e^{-\lambda T}),$$

for some K > 0,  $\varkappa \in \mathbf{N}$  and  $\lambda > 0$  independent of h and T.

# Corollary

Theorem 1:

$$\begin{split} \mathbf{E}\varphi(\mathbf{X}_{x}(T)) - \mathbf{E}\varphi(\mathbf{X}_{N}) &= C_{0}(T, x)h + C(T, x)h^{2}, \\ C_{0}(T, x) &= \mathbf{E}\int_{0}^{T}B_{0}(t, \mathbf{X}_{x}(t))\mathrm{d}t. \end{split}$$

# Corollary

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Theorem (2; [Leimkuhler, Matthews, T.; Proc. R. Soc. A (2014)]) Let Assumptions 1-2 hold. Then the coefficient  $C_0(T, x)$  from (47) goes to zero as  $T \to \infty$ :

$$|C_0(T,x)| \le K(1+|x|^{\varkappa})e^{-\lambda T}$$
(48)

for some constants K > 0,  $\varkappa \in \mathbf{N}$  and  $\lambda > 0$ , i.e., over a long integration time the scheme (41) is of order two up to exponentially small correction.

# Sketch of the proof

$$C_{0}(T,x) = \int_{0}^{T} \mathbf{E}B_{0}(t, \mathbf{X}_{x}(t)) dt = \int_{0}^{T} \int_{\mathbf{R}^{d}} B_{0}(t, y) p(t, x, y) dy dt \quad (49)$$
$$= \int_{0}^{T} \int_{\mathbf{R}^{d}} B_{0}(t, y) \rho(y) dy dt$$
$$+ \int_{0}^{T} \int_{\mathbf{R}^{d}} B_{0}(t, y) [p(t, x, y) - \rho(y)] dy dt,$$

where p(t, x, y) is the transition density for (34) and  $\rho(y)$  is the invariant density.

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$$\int_{\mathbf{R}^d} B_0(t, y) \exp\left(-\frac{2}{\sigma^2} U(y)\right) \mathrm{d}y = 0.$$
 (50)

We emphasize that the fact that the average of B<sub>0</sub>(t, x) with respect to the invariant measure is equal to zero is the reason why the scheme (41) is second order accurate in approximating ergodic limits.

- We emphasize that the fact that the average of B<sub>0</sub>(t, x) with respect to the invariant measure is equal to zero is the reason why the scheme (41) is second order accurate in approximating ergodic limits.
- 2 In the case of the Euler scheme (37) we get the same error expansion as (46) for the scheme (41) but with a different  $B_0(t,x) = B_0^E(t,x)$  (see [Milstein, T.; Springer 2004]):

$$B_0^E(t,x) = \frac{1}{2} \left[ \sum_{i,j=1}^d a^j \frac{\partial u}{\partial x^j} a^j \frac{\partial u}{\partial x^i} + \frac{\sigma^2}{2} \sum_{i,j=1}^d \frac{\partial^2 a^j}{\partial x^i} \frac{\partial u}{\partial x^j} + \frac{\sigma^2}{2} \sum_{i,j=1}^d a^j \frac{\partial^3 u}{\partial x^i (\partial x^j)^2} + \sigma^2 \sum_{i,j=1}^d \frac{\partial^2 a^j}{\partial x^i} \frac{\partial^2 u}{\partial x^j \partial x^i} + \frac{\sigma^4}{6} \sum_{i,j=1}^d \frac{\partial^4 u}{(\partial x^i)^2 (\partial x^j)^2} \right]$$

The average of  $B_0^E(t,x)$  with respect to the invariant measure is not equal to zero and, consequently, the Euler scheme (37) approximates ergodic limits with order one – the same order as its weak convergence over a finite time interval.

3 Let a one-step weak approximation X
<sub>t,x</sub>(t + h) of the solution X
<sub>t,x</sub>(t + h) of (34) generate a method of order p. The global error of the method:

$$R := \mathbf{E}\varphi(\mathbf{X}_{x}(T)) - \mathbf{E}\varphi(\bar{\mathbf{X}}_{x}(T))$$

$$= C_{0}(T, x)h^{p} + \dots + C_{n}(T, x)h^{p+n} + O(h^{p+n+1}),$$
(51)

where  $n \in \mathbf{N}$  and the functions  $C_0(T, x), \ldots, C_n(T, x)$  are independent of h which can be presented in the form

$$C_i(T,x) = \int_0^T \mathbf{E}B_i(s, \mathbf{X}_x(s)) \mathrm{d}s.$$

One can deduce from the proof of Theorem 2 that if the averages of  $B_i(s, x)$   $0 \le i \le n$ , with respect to the invariant measure are equal to zero then in the limit of  $T \to \infty$  the scheme has p + n order of accuracy in h.

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[Abdulle, Vilmart, Zygalakis 2014-15]

# Numerical experiments

**Anharmonic scalar model:** the one-dimensional potential energy U(x) = cos(x)



Figure: The error in computed distributions is plotted for each scheme.

[Leimkuhler, Matthews, T.; Proc. R. Soc. A (2014)]

#### Error in finite time



Figure: The lower plot shows the error in the distribution after time t, as computed using each scheme at h = 0.16. In the plots at the top, we compare the error growth with respect to stepsize h at multiples of t = 0.96.

# Conclusions

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- Some basics of stochastic numerics
- How to compute ergodic limits
- A 'non-Markovian' scheme tailored towards computing ergodic limits

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- How to compute ergodic limits
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#### Further numerical aspects for SDEs:

- Variance reduction/reduction of complexity
- Geometric integration and computing ergodic limits (see Part II on Friday)
- DEs with small noise, DEs with coloured noise
- Convergence/approximation of schemes in the case of nonglobally Lipschitz coefficients of SDEs
- Mean-square and weak approximation of SDEs in bounded domains (absorbtion, reflection)







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