

## Analysis of SPDEs and numerical methods for UQ

#### John Burkardt<sup>†</sup> & Clayton Webster\*

Thanks to Max Gunzburger & Guannan Zhang (FSU), Fabio Nobile (MOX), Raul Tempone (KAUST)

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#### Why stochastic models? A transition to non-deterministic simulations



# Many applications (especially those predicting future events) are affected by a relatively large amount of uncertainty in the input data such as model coefficients, forcing terms, boundary conditions, geometry, etc.

- An example includes forecasting financial markets where this may depend on the number of economic factors, number of underlying assets or the number of time points/time steps, human behaviors, etc.
- Important DOE examples include the enhancement of reliability of smart energy grids, development of renewable energy technologies, vulnerability analysis of water and power supplies, understanding complex biological networks, climate change estimation and design and licensing of current and future nuclear energy reactors (CASL simulation hub)
- The model itself may contain an incomplete description of parameters, processes or fields (not possible or too costly to measure)
- There may be small, unresolved scales in the model that act as a kind of background noise (i.e. macro behavior from micro structure)

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### Types of uncertainties



Stochastic models give quantitative information about uncertainty. In practice it is necessary to address the following types of uncertainties:

- Uncertainty may be *aleatoric* which means random and is due to the intrinsic variability in the system
   Remark: by variability we mean a type of uncertainty that is inherent and irreducible, e.g. turbulent fluctuations of a flow field around an airplane wing, permeability in an aquifer, etc.
   OR
- Uncertainty may be *epistemic* which means due to incomplete knowledge **Remark:** can be reduced by additional experimentation, improvements in measuring devices, etc., e.g. mechanic properties of many bio-materials, polymeric fluids, highly heterogeneous or composite materials, the action of wind or seismic vibrations on civil structures, etc.

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## Uncertainty quantification (UQ) attempts to quantitatively access the impact of uncertain data on simulation outputs



Let A denote the input data and B the output set such that  $u : A \mapsto B$ . There are various ways to describe the uncertainty in A with the goal of describing the uncertainty in some quantity of interest (Qol) Q(u):

Worst scenario approaches: Typically A is an  $\epsilon$ -ball around some nominal input data and the goal is to determine the worst case associated with the set relation B = u(A)

The range of the uncertainty of Q(u) is then defined by the interval I

$$I = \left[\underline{Q}(u), \overline{Q}(u)\right] = \left[\inf_{a \in A} Q(u(a)), \sup_{a \in A} Q(u(a))\right]$$

• the choice of the input set A is, in a large way, subjective and should be regarded as a working assumption

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#### Uncertainty quantification (UQ) Knowledge-based methods



**Fuzzy sets and possibility theory**: deterministic approach to UQ which generalizes classical set theory. Let  $C \subset A$ :

- for each  $x \in A$  set membership is defined by  $\mu_C : A \to [0, 1]$ , expressing the degree of truth of the statement "x belongs to C."
- define the  $\alpha$ -cut of C by  $C_{\alpha} \stackrel{\text{def}}{=} \{x \in A : \mu_C \ge \alpha\}$  which gives a set characterization of uncertainty

#### • the operator u then propagates the fuzziness in A into the fuzziness in B

**Evidence theory (Dempster-Shafer Theory)**: generalizing the probabilistic approach by defining the Belief Bel(C) (lower bound) and Plausibility Pl(C) (upper bound) functions, for the likelihood of an event C.

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$$\sum_{\varphi \in \Phi} m(\varphi) = 1, \quad m(\emptyset) = 0, \quad \text{however } \varphi_1 \subset \varphi_2 \not\Rightarrow m(\varphi_1) \leq m(\varphi_2)$$

$$Bel(C) = \sum_{\varphi \in \Phi, \varphi \subset C} m(\varphi), \quad Pl(C) = \sum_{\varphi \in \Phi, \varphi \cap C \neq \emptyset} m(\varphi)$$



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## Uncertainty quantification (UQ) Stochastic/Probabilistic methods



What is probability theory? Understood as a mathematical theory of a finite measure. Let  $(\Omega, \mathscr{F}, \mathbb{P})$  denote a (complete) probability space:  $\Omega$  is the event space,  $\mathscr{F} \subset 2^{\Omega}$  is the  $\sigma$ -algebra and  $\mathbb{P}$  is the probability measure, satisfying:

- $0 \leq P(A), \text{ if } A \in \mathscr{F} \text{ and } P(\Omega) = 1$
- ② A **positive measure**  $\mathbb{P}:\mathscr{F} o [0,1]$  which is countably additive

$$\mathbb{P}\left(\bigcup_{i=1}^{\infty} A_i\right) = \sum_{i=1}^{\infty} \mathbb{P}(A_i), \text{ for } \{A_i\}_{i=1}^{\infty} \in \mathscr{F} \text{ disjoint}$$

**Stochastic / probabilistic methods**: given a probability measure on the input data A the mapping u induces a probability measure on the output set  $B \implies$  SODEs/SPDEs (**Doob-Dynkin Lemma**)

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Models (linear or nonlinear) for a system  $\mathcal{M}$  may be stationary with state u, exterior loading or forcing f and random model description (realization)  $\omega \in \Omega$ , with probability measure  $\mathbb{P}$ :

$$\mathcal{M}(\omega; u) = f(\omega) \quad \text{a.e. in } D \subset \mathbb{R}^d$$

Evolution in time may be

• discrete (e.g. Markov chain), driven by discrete random process:

 $u_{n+1} = \mathcal{F}(\omega; u_n)$ 

 continuous (e.g. Markov process ≡ Stochastic Differential Equation), driven by random processes:

 $du = (\mathcal{M}(\omega; u) - f(\omega, t))dt + \mathcal{B}(\omega; u)dW(\omega, t) + \mathcal{P}(\omega; u)dQ(\omega, t).$ In this Itô evolution equation,  $W(\omega, t)$  is a Wiener process and  $Q(\omega, t)$  is a Poisson process.



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Motivation UQ Stochastic models RVs RFs SVD KLE Summary

### The Computational Stochastic PDE

Forward Problem: From Real World to Predictions to Decisions



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Forward Problem: From Real World to Predictions to Decisions



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Consider the following simple equation describing the pressure distribution in a porous medium:

$$\begin{cases} -\nabla \cdot (a(x) \nabla u(x)) &= f(x) & \text{ in } D \subset \mathbb{R}^d, \\ u(x) &= 0 & \text{ on } \partial D, \end{cases}$$

where f(x) is the source, a(x) describes the permeability and u(x) is the pressure distribution.

Q: What if the input data is random?

$$\left\{ \begin{array}{rl} -\nabla \cdot \left( a(\omega,x) \, \nabla u(\omega,x) \right) &= f(\omega,x) & \text{ in } \Omega \times D, \\ u(\omega,x) &= 0 & \text{ on } \Omega \times \partial D, \end{array} \right.$$

where  $f(\omega, x)$  is the stochastic source,  $a(\omega, x)$  describes the stochastic permeability and  $u(\omega, x)$  is the stochastic pressure distribution.

• Difficulties arise since instead of just asking for  $u(x_0)$  we instead want to know  $\mathbb{E}[u](x_0)$ ,  $\mathbb{V}ar[u](x_0)$  or even  $\mathbb{P}[u(\omega, x_0) \ge u_0]$ 



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An example problem Pressure distribution in a porous medium



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There have been many formulations and approaches to solve SPDEs:

- Statistical sampling methods:
  - Brute-force Monte Carlo (MC): convergence rate independent of the number of random variables, robust, embarrassingly parallel very slow convergence  $(1/\sqrt{\#\text{samples}})$
  - Quasi MC (QMC), Latin Hypercube Sampling (LHS), Lattice Rules
  - Variance reduction techniques: important, conditional and correlated sampling - limitations when confronted with large number of RVs
- Indirect methods (require closure approx.):
  - Moment methods: derive equations for the **moments** of the quantities of interest not applicable to nonlinear problems or non-Gaussian RVs
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- Direct methods: compute an approximation to u(ω, x) in a suitable subspace and use this solution to compute the desired statistics, e.g. Stochastic Galerkin, Stochastic Collocation.
  - Interval analysis : maximum bounds of output uncertainty can dramatically overestimated to uncertainties
  - Perturbation-based methods : Taylor expansion around a mean solution can only be used for linear Qols and when the variance in solution is small
  - Operator-based methods : compute the inverse of a given operator, if it exists, by using a Neumann series expansion or the weighted integral method restricted to small magnitude uncertainties and often limited to static problems
  - Stochastic polynomial approximations: spectral Galerkin, Wiener (polynomial) chaos, Karhunen-Loève, stochastic collocation - challenges include: determining the proper polynomial subspace, *curse of dimensionality*, adaptive and anisotropic refinement, low stochastic regularity and discontinuities, optimization and inverse problems



- Obtermine an accurate representation for the input stochastic (random) fields can be both simulated and analyzed. e.g. a Karhunen-Loève expansion
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The input data  $a(\omega, x)$ ,  $f(\omega, x)$  and the solution  $u(\omega, x)$  of the SPDE will (more than likely) be a random field defined by a set of random variables  $\mathbf{Y}(\omega) = (Y_1(\omega), \dots, Y_N(\omega))$ , i.e.  $u(\omega, x) = u(\mathbf{Y}(\omega), x)$ 

**1** How to deal with RVs  $Y_n(\omega)$ ?

2 How to represent RFs  $a(\mathbf{Y}(\omega), x)$ ,  $f(\mathbf{Y}(\omega), x)$  and  $u(\mathbf{Y}(\omega), x)$ ?

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A  $\mathbb{R}$ -valued random variable (RV)  $Y : (\Omega, \mathscr{F}, \mathbb{P}) \to \mathbb{R}$  is completely specified by a probability density function (pdf)  $\rho_Y$  and a cumulative distribution function (cdf) and  $F_Y$  s.t.  $\forall y \in \mathbb{R}$ :

$$\begin{split} F_Y(\lambda) &:= \mathbb{P}\left[\{Y(\omega) \le \lambda\}\right] = \int_{\{Y(\omega) \le \lambda\}} d\mathbb{P}(\omega) = \mathbb{E}\left[\chi_{\{Y(\omega) \le \lambda\}}\right] \\ &= \int_{-\infty}^{\lambda} \rho_Y(y) dy, \end{split}$$



where  $\int_{\mathbb{R}} \rho_Y(y) dy = 1$ 

#### Random variables Covariance, correlation and independence



For  $Y\in L^1_{\mathbb{P}}(\Omega)$  define the expected (mean) by

$$\overline{Y} = \mathbb{E}\left[Y\right] = \int_{\Omega} Y(\omega) d\mathbb{P}(\omega) = \int_{\mathbb{R}} y \rho_Y(y) dy$$

and fluctuating part by  $\widetilde{Y}=Y(\omega)-\overline{Y}(\omega),$  with  $\mathbb{E}[\widetilde{Y}]=0.$ 

- The variance  $\mathbb{V}ar[Y] = \mathbb{E}\left[\widetilde{Y} \otimes \widetilde{Y}\right] = \mathbb{E}\left[(\widetilde{Y})^2\right] = \mathbb{C}ov[Y,Y]$
- Let  $\mathbf{Y}(\omega) = (Y_n(\omega))_{n=1}^N$ ,  $N \in \mathbb{N}_+$  be a random vector, then the covariance and correlation of two RVs:

$$\mathbb{C}ov[Y_i, Y_j] := \mathbb{E}\left[\widetilde{Y}_i \otimes \widetilde{Y}_j\right], \quad \mathbb{C}orr = \frac{\mathbb{C}ov[Y_i, Y_j]}{\sqrt{\mathbb{V}ar[Y_i]}\sqrt{\mathbb{V}ar[Y_j]}}$$

- uncorrelated if  $\mathbb{C}ov[Y_i, Y_j] = 0$  (orthogonal), perfectly correlated if  $\mathbb{C}orr = 1$  and perfectly anti-correlated if  $\mathbb{C}orr = -1$
- independent if  $\forall \phi_1, \phi_2, \mathbb{E}[\phi_1(X_1)\phi_2(X_2)] = \mathbb{E}[\phi_1(X_1)]\mathbb{E}[\phi_2(X_2)]$

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We can now define a covariance matrix  $\mathbb{C}ov[\mathbf{Y}]$  whose (i, j) entry is the covariance between  $Y_i(\omega)$  and  $Y_j(\omega)$ :

 $\mathbb{C}ov[\mathbf{Y}]_{(i,j)} = \mathbb{C}ov[Y_i, Y_j]$ 

•  $\mathbb{C}ov[\mathbf{Y}]$  is symmetric, nonnegative definite, and has diagonal elements  $\mathbb{C}ov[\mathbf{Y}]_{(i,i)} = \mathbb{V}ar[Y_i]$ 

As before, the correlation matrix can be defined from the covariance matrix. Form a diagonal matrix  $\Sigma$  from the square roots of the variances, then compute the correlation matrix by:

 $\mathbb{C}orr[\mathbf{Y}] = \Sigma^{-1} \mathbb{C}ov[\mathbf{Y}]\Sigma$ 

- the diagonal entries of  $\mathbb{C}orr[\mathbf{Y}]$  are 1
- The Cauchy-Schwarz inequality guarantees that the off-diagonal elements lie between -1 and +1
- value of each covariance entry indicates the strength and direction of the correlation between the corresponding components





We can now define a covariance matrix  $\mathbb{C}ov[\mathbf{Y}]$  whose (i, j) entry is the covariance between  $Y_i(\omega)$  and  $Y_j(\omega)$ :

 $\mathbb{C}ov[\mathbf{Y}]_{(i,j)} = \mathbb{C}ov[Y_i, Y_j]$ 

•  $\mathbb{C}ov[\mathbf{Y}]$  is symmetric, nonnegative definite, and has diagonal elements  $\mathbb{C}ov[\mathbf{Y}]_{(i,i)} = \mathbb{V}ar[Y_i]$ 

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A random process/field (RP/RF)  $a(\omega, x)$  defined on a probability space  $(\Omega, \mathscr{F}, \mathbb{P})$  and indexed by a deterministic domain  $D \subset \mathbb{R}^d$ , returns a real value

- () a set of RVs indexed by  $x \in D$ . For every  $x \in D$ ,  $a(\cdot, x)$  is a RV on  $\Omega$
- ② a function-valued RV. For every  $\omega \in \Omega$ ,  $a(\omega, \cdot)$  is a random function a realization of x in the domain D

Often only second order information - mean and covariance are known



• Mean  $\overline{a}(x) = \mathbb{E}[a](x) = \int_{\Omega} a(\omega, \cdot) d\mathbb{P}(\omega)$  and  $\mathbb{V}ar[a](x) = \mathbb{E}[(\widetilde{a})^2](x)$  as a function of x with fluctuation part  $\widetilde{a}(\omega, x) = a - \overline{a}$ 

•  $\mathbb{P}\left[a \ge a_0\right] = \mathbb{P}\left[\left\{\omega \in \Omega \, : \, a(\omega, x) \ge a_0\right\}\right] = \mathbb{E}\left[\chi_{\{a \ge a_0\}}\right]$ 



The covariance may be considered at different spatial positions  $x \in D \subset \mathbb{R}^d$ :

 $\mathbb{C}ov[a](x_1, x_2) := \mathbb{E}\left[\widetilde{a}(\cdot, x_1)\widetilde{a}(\cdot, x_2)\right], \text{ for } x_1, x_2 \in D \times D$ 

- if  $\overline{a}(x) \equiv \overline{a}$  and  $\mathbb{C}ov[a](x_1, x_2) = C_a(x_1 x_2)$  then the process is homogeneous. Here representation through the spectrum using Fourier expansion if well known
- we will consider colored noise approximations using correlated second-order RFs: a(ω, x) ∈ L<sup>2</sup><sub>P</sub>(Ω; W(D)) where

$$L^2_{\mathbb{P}} := \left\{ a \text{ measurable and } \int_{\Omega} \|a(\omega, \cdot)\|^2_{W(D)} d\mathbb{P}(\omega) < +\infty \right\}$$

and W(D) a Banach space of functions  $a: D \to \mathbb{R}$ • we will not focus on white noise approximations which refers to uncorrelated RFs for which:

$$\overline{a}(x) = 0$$
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## Discretized white noise Piecewise constant approximations





Discretized white noise over a square subdivided into  $2, \ldots, 512$  triangles



- Ideally, this involves a combination of countably many independent RVs
  - $\Gamma_n \equiv Y_n(\Omega) \subset \mathbb{R}$  and  $\Gamma = \prod_{n=1}^N \Gamma_n \subset \mathbb{R}^N$ , the image of the random vector  $\mathbf{Y} = (Y_1, \dots, Y_N)$
  - Let  $\rho: \Gamma \to \mathbb{R}_+$ , with  $\rho \in L^{\infty}(\Gamma)$  be the joint probability density function (jPDF) of **Y**, then we want that:

$$ho(\mathbf{y}) = \prod_{n=1}^N 
ho_n(y_n), ext{ where } \mathbf{y} \in \Gamma ext{ and } orall n, y_n \in \Gamma_n$$

- The independence of the  $N\ {\rm RVs}$  allows to see each of them as the axis of a coordinate system (Doob-Dynkin lemma)
- The most popular approach: Karhunen-Loève (KL) expansion involves an ∞-dimensional expansion of the random field suitably truncated
- Challenge: this truncation in N RVs can be high-dimensional
- In the discrete case KL is similar to ROM and the SVD of a matrix



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Motivate the Karhunen-Loève expansion SVD: The Singular Value Decomposition



$$A = U S V^T$$

where

- U is an m by m orthogonal matrix  $(U^T U = I)$ ;
- S is an m by n diagonal matrix with nonnegative entries;
- V is an n by n orthogonal matrix;

The diagonal entries of S, called the singular values of A, are chosen to appear in descending order, and are equal to the square roots of the nonzero eigenvalues of  $AA^T$  or  $A^TA$ 

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- r, the number of nonzero diagonal elements in S, is the rank of A
  - very small nonzeros may indicate numeric singularities
- the *i*-th diagonal element of S is the *i*-th largest eigenvalue of  $AA^T$  (and also of  $AA^T$ ). Hence, we may write this value as  $\sqrt{\lambda_i}$ .
- Let  $u_i$  and  $v_i^T$  be the *i*-th columns of U and  $V^T$ . Then A maps the *i*-th column of  $V^T$  to the *i*-th column of U.
- The columns of U and V provide a singular value expansion of A:

$$A = \sum_{i=1}^{r} \sqrt{\lambda_i} \, u_i \, v_i^T$$



If we use all r terms, the singular value expansion is exact

Let  $A^k$  represent the sum of just the first k terms of the expansion

- $A^k$  is a matrix of rank k, the sum of k rank-1 outer products
- Of all rank k matrices,  $A^k$  is the best approximation to A in two senses:
- O Minimum L<sup>2</sup> norm:

$$\begin{split} \left\|A - A^k\right\|_{L^2} \equiv & \text{square root of maximum eigenvalue of } (A - A^k)^T (A - A^k) \\ \left\|A - A^k\right\|_{L^2}^2 = & s_{k+i}^2 = \lambda_{k+1} \end{split}$$

Ø Minimum Frobenius (sum of squares) norm:

$$\|A - A^{k}\|_{F} \equiv \sqrt{\sum_{i,j} (A_{i,j} - A_{i,j}^{k})^{2}}$$
$$\|A - A^{k}\|_{F}^{2} = \sum_{k+1}^{r} s_{i}^{2} = \sum_{k+1}^{r} \frac{\lambda_{i}}{\lambda_{i}}$$



 $\boldsymbol{U}$  and  $\boldsymbol{V}$  are natural bases for the input and output of  $\boldsymbol{A}$ 

In the natural bases, the SVD shows that multiplying by A is simply stretching the *i*-th component by  $s_i$ :

$$x = \sum_{i=1}^{r} v_i^T * c_i \implies y = A * x = \sum_{i=1}^{r} u_i * (s_i * c_i)$$

- The relative size of the singular values indicates the importance of each column
- $\bullet\,$  The singular value expansion produces an optimal, indexed family of reduced order models of  $A\,$



- SVD is the discrete version of the Karhunen-Loève (KL) expansion that is typically applied to RF that produce, for any time *t*, a field of values varying spatially with *x*
- Since it's easier to understand discrete problems, let's prepare for the KL expansion by looking at how the SVD is used with a set of data
- Let us re-imagine the columns of our discrete data as being n snapshots in discrete time indexed by j. Each snapshot will record m values in a "space" indexed by i



- If we pack our data into a single matrix A, then  $A_{i,j}$  means the measurement at position i and time j
- It is reasonable to expect correlation in this data; the "neighbors" of  $A_{i,j}$ , in either space or time, might tend to have similar values
- Moreover, the overall "shape" of the data for one time or one spatial coordinate might be approximately repeated elsewhere in the data
- This is exactly the kind of behavior the SVD can detect

			Space									
Time	1890	1	12	12	33	29	22	3	0			
	1891	0	31	23	44	18	13	1	0			
	1892	0	23	44	25	17	17	13	1			
	1893	1	30	49	37	15	23	10	1			
	1894	0	30	18	74	9	5	0	2			



We have a data file of the monthly snowfall in inches, over 121 winters at Michigan Tech. We'll think of the months as the "space" dimension.

Year	<u>Oct</u>	Nov	Dec	<u>Jan</u>	<u>Feb</u>	Mar	Apr	May	Tot
1890	1	12	12	33	29	22	3	0	112
1891	0	31	23	44	18	13	1	0	130
1892	0	23	44	25	17	17	13	1	140
1893	1	30	49	37	15	23	10	1	166
1894	0	30	18	74	9	5	0	2	138
2006	6	6	27	38	37	20	31	0	165
2007	0	21	40	55	32	24	14	0	186
2008	0	17	70	85	27	5	15	0	219
2009	3	4	87	39	19	0	0	0	152
2010	0	26	33	72	18	13	18	0	180

http://www.mtu.edu/alumni/favorites/snowfall/snowfall.html



To analyze our data, we consider each of the 121 snowfall records, starting with  $x^{1890}$ , as a column of 8 numbers, and form the m=8 by n=121 matrix A:

$$A = \left[x^{1890} | x^{1891} | \dots | x^{2010}\right]$$

- Now we determine the SVD decomposition  $A = USV^T$
- The columns of U are an orthogonal set of "spatial" behaviors or modes (typical behavior in a fixed year over a span of months)
- $\bullet\,$  The columns of V are typical behaviors or modes in a fixed month over a span of years.
  - In both cases, the most important behaviors are listed first
- $\bullet$  The diagonal matrix S contains the "importance" or "energy" or signal strength associated with each behavior



The S data shows the relative importance of the first two modes is:

$$\frac{s1^2}{\sqrt{\sum_{i=1}^8 s_i^2}} = 0.87 \quad \frac{s2^2}{\sqrt{\sum_{i=1}^8 s_i^2}} = 0.05$$

The first pair of modes,  $u_1$  and  $v_1$ , by itself, can approximate the entire dataset with a relative accuracy of 87%.









1 December/January High (DOMINANT)



3 February High, less January



2 More December, less later



J. Burkardt, http://www.sc.fsu.edu/~burkardt, C. Webster, http://www.csm.ornl.gov/~cgwebster — April 2-3, 2012







The same kind of approximating is occurring for all 121 sets of data!



The linear regression line suggests the "December/January High" pattern (upper left) is steadily gaining importance over the years.



J. Burkardt, http://www.sc.fsu.edu/~burkardt, C. Webster, http://www.csm.ornl.gov/~cgwebster — April 2-3, 2012



To see how heaviest snowfall is coming earlier, compare the 1890 January/February style snowfall with the 2008 December/January style:





Oata gathered at discrete places and times is easier to understand than the corresponding continuous cases

② The SVD shows how underlying patterns and correlations can be detected, and represented as a sum of the form

$$A = \sum_{i=1}^{r} \sqrt{\lambda_i} \, u_i \, v_i^T$$

where the  $\lambda$  values represent a strength, the u 's represent variation in space, and v variation in time

O The structure of the u and v vectors suggests something about the preferred modes of the system, and the size of the λ coefficients allows us to understand the relative important of different modes, and to construct reduced order models if we wish



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- Given that our data was stored in A, we may think of the matrices  $AA^T$  and  $A^TA$  as a form of a covariance matrix
- The singular values  $\sqrt{\lambda_i}$  are the square roots of eigenvalues of both these matrices
- U contains eigenvectors of the "spatial" covariance matrix  $AA^T$
- ${\ensuremath{\, \bullet }}\xspace V$  contains eigenvectors of the "temporal" covariance matrix  $A^TA$
- Very similar statements will hold for the continuous case



spectral expansion of its covariance function - other names Proper Orthogonal Decomposition (POD), Principle Component Analysis (PCA)

Let  $a(\omega,x)\in L^2_{\mathbb{P}}$  be a RF with continuous covariance  $\mathbb{C}_a:D\times D\to \mathbb{R}$ 



Examples of 1d covariance kernels for correlation lengths  $L_c = 1$  and  $L_c = 1/4$ 





- $\textbf{0} \ \mathbb{C}_a \text{ is symmetric, i.e. } \mathbb{C}_a(x_1,x_2) = \mathbb{C}_a(x_2,x_1), \, \forall x_1,x_2 \in D$
- **2**  $\mathbb{C}_a$  is non-negative definite, i.e. for any  $n = 1, \ldots$

$$\begin{split} &\sum_{i=1}^n \sum_{j=1}^n \mathbb{C}_a(x_i, x_j) v_i v_j \geq 0, \forall (x_1, \dots, x_n) \in D^n \text{ and } (v_1, \dots, v_n) \in \mathbb{R}^n.\\ &\text{In matrix notation: } v^T \mathbb{C}_a(x, x) v \geq 0, \forall v, x. \end{split}$$

Define the associated linear covariance operator  $T_{\mathbb{C}_a}: L^2(D) o L^2(D)$  s.t.:

$$[T_{\mathbb{C}_a}f](x_1) = \int_D \mathbb{C}_a(x_1, x_2)f(x_2)dx_2, \,\forall f \in L^2(D)$$

- it has a countable sequence of real eigenvalues  $\{\lambda_n\} \subset \mathbb{R}_+$ ,  $\lambda_n \to 0$
- corresponding eigenfunctions  $\{b_n(x)\}$  are  $L^2(D)$ -orthonormal





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Karhunen-Loève expansion Spectral representation of the kernel



# Theorem (Mercer, 1909)

Given  $\mathbb{C}_a$  continuous, symmetric, non-negative definite, then:

$$\lim_{N \to \infty} \max_{(x_1, x_2) \in D \times D} \left| \mathbb{C}_a(x_1, x_2) - \sum_{n=1}^N \frac{\lambda_n b_n(x_1) b_n(x_2)}{\lambda_n b_n(x_2)} \right| = 0$$

Mercer's spectral representation of the kernel:

$$\mathbb{C}_a(x_1, x_2) = \sum_{n=1}^{\infty} \lambda_n b_n(x_1) b_n(x_2)$$

Eigenvalues/eigenfunctions constructed from a 2nd-order Fredholm equation:

$$\begin{bmatrix} T_{\mathbb{C}_a} b_n \end{bmatrix}(x_1) = \int_D \mathbb{C}_a(x_1, x_2) b_n(x_2) dx_2 = \lambda_n b_n(x_1), \quad n = 1, \dots$$
  
with 
$$\int_D b_n(x_1) b_m(x_1) dx_1 = \delta_{nm}$$



Karhunen-Loève expansion



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Given  $\mathbb{C}_a$  continuous, symmetric, non-negative definite, then:

$$\lim_{N \to \infty} \max_{(x_1, x_2) \in D \times D} \left| \mathbb{C}_a(x_1, x_2) - \sum_{n=1}^N \lambda_n b_n(x_1) b_n(x_2) \right| = 0$$

Mercer's spectral representation of the kernel:

$$\mathbb{C}_a(x_1, x_2) = \sum_{n=1}^{\infty} \frac{\lambda_n}{b_n(x_1)} b_n(x_2)$$

Eigenvalues/eigenfunctions constructed from a 2nd-order Fredholm equation:

$$\begin{split} \left[T_{\mathbb{C}_a}b_n\right](x_1) &= \int_D \mathbb{C}_a(x_1, x_2)b_n(x_2)dx_2 = \lambda_n b_n(x_1), \quad n = 1, \dots \\ & \text{with} \quad \int_D b_n(x_1)b_m(x_1)dx_1 = \delta_{nm} \end{split}$$



### Karhunen-Loève expansion Approximating an ∞-dimensional RF



$$a(\omega, x) = \overline{a}(x) + \sum_{n=1}^{+\infty} b_n(x) Y_n(\omega),$$

with 
$$Y_n(\omega) = \int_D (a(\omega, x) - \overline{a}(x))b_n(x) dx$$

•  $(\lambda_n, b_n(x))$  are eigenpairs of  $T_{\mathbb{C}_a}$ ;  $Y_n(\omega)$  are centered, uncorrelated RVs:

$$\mathbb{E}[Y_n] = 0, \, \mathbb{C}ov[Y_n, Y_m] = \mathbb{E}[Y_n Y_m] = \delta_{nm}$$

but not necessarily independent, with  $\mathbb{V}ar[Y_n] = \lambda_n$ 

If the basis  $\{b_n\}$  has spectral approx. properties and the realizations of a are smooth, then  $\lambda_n = \mathbb{V}ar[Y_n] \to 0$  sufficiently fast as  $n \to \infty$  and we can truncate the series

$$a(\omega, x) \approx a_N(\omega, x) = \overline{a}(x) + \sum_{n=1}^N b_n(x) Y_n(\omega),$$

Rate of decay depends on the smoothness of  $\mathbb{C}_a$  and the corr. length  $L_c$ 

J. Burkardt, http://www.sc.fsu.edu/~burkardt, C. Webster, http://www.csm.ornl.gov/~cgwebster — April 2-3, 2012



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This truncated expansion corresponds to the Best N-term approximation

$$\min_{\substack{(\mathbf{Y}_n, b_n)\\ \int_D b_n b_m = \delta_{nm}}} \mathbb{E}\left[\int_D \left(a - \overline{a}(x) - \sum_{n=1}^N b_n \mathbf{Y}_n\right)^2\right]$$

 $\bullet~$  If we truncate using the N largest eigenvalues, we have an optimal - in variance - expansion in N~ RVs

• i.e. with  $\mathbb{C}_a$  continuous,  $a_N$  converges uniformly to a (Mercer's Thm)

$$\sup_{x \in D} E[(a-a_N)^2](x) = \sup_{x \in D} \left\{ \mathbb{C}_a(x) - \sum_{n=1}^N \lambda_n b_n^2(x) \right\} \to 0, \text{ as } N \to \infty$$

• KL expansion is the SVD of the map  $A : L^2(D) \to L^2_{\mathbb{P}}(\Omega)$ , where  $\mathbb{C}_a := A^*A$  (i.e. truncate at the N largest eigenvalues of  $A^TA$ )  $\Rightarrow$  finding a sparse representation (model reduction)





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## Karhunen-Loève expansion Example eigenfunctions





Distributions of the eigenfunctions of 1d exponential kernels,  $\mathbb{C}_a(x_1, x_2) = \sigma^2 e^{-\frac{|x_1-x_2|}{L_c}}$ , for correlation lengths  $L_c = 1$  and  $L_c = 1/4$ 



Eigenvalues values of the kernels for correlation lengths  $L_c = 1$  and  $L_c = 1/4$ 

- More modes required as the correlation decreases (noise level increases)
- In the asymptotic limit of white noise  $\Rightarrow$  infinity number of modes
- For a given  $L_c$ , the smoothness of the covariance kernel  $\mathbb{C}_a$  dictates the convergence rate of the eigenvalues



## Karhunen-Loève expansion



- the truncation error decreases monotonically with the number of terms in the expansion
- the convergence is inversely proportional to the correlation length and depends on the regularity of the covariance kernel

## Theorem (Schwab et al. 2005)

• If  $\mathbb{C}_a$  is piecewise analytic on  $D \times D$  with  $D \subset \mathbb{R}^d$  then:

$$0 \leq \lambda_n \leq c_1 \exp(-c_2 n^{1/d}), \quad \forall c_1, c_2 > 0 \text{ ind. of } n$$

• If  $\mathbb{C}_a$  is piecewise  $H^k \otimes L^2$  with k > 0 then:

$$0 \le \lambda_n \le c_3 \ n^{-k/d}, \quad \forall c_3 > 0$$

**Remark**: similar to SVD - if one wants the relative error (in the variance) less than some tolerance  $\delta$ , i.e.  $||a - a_N||^2 \le ||a|| \delta$ , then choose N s.t.  $\sum_{n=N+1}^{\infty} \lambda_n \le \delta \sum_{n=1}^{\infty} \lambda_n$ 





Exponential and Gaussian kernels with  $L_c = 1$  and N = 5 modes



Exponential and Gaussian kernels with  $L_c = 1/4$  and N = 20 modes





Although the  $Y_n$ 's are uncorrelated, in general they are not independent

• Gaussian RVs are uncorrelated  $\iff$  independent

A very convenient special case is the one where the RF  $a(\omega, x)$  is a Gaussian RF, defined by Gaussian random vector  $\mathbf{Y}(\mathbf{\Omega}) = \Gamma = \prod_{n=1}^{N} \Gamma_n$  s.t.

$$ho(\mathbf{y}) = \prod_n^N 
ho_n(y_n), ext{ where } y_n \in \Gamma_n$$

**Question:** Is  $a_N$  strictly positive? **Answer:** Not necessarily!

**Warning**: the truncated expansion might not be positive almost surely! Possible remedy: nonlinear transformations

$$\Rightarrow a(\omega, x) \approx a_N(\omega, x) = a_{\min} + e^{\gamma_N(\omega, x)} = a_{\min} + e^{b_0(x) + \sum_{n=1}^N b_n(x)Y_n(\omega)}$$

where  $\gamma_N(\omega, x)$  is a truncated Gaussian RF and  $a_N(\omega, x)$  is a lognormal RF





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• If  $a(\omega, x)$  is the given RF and if the marginal CDF  $F_a$  is known or can be approximated, then one can use a translation process, i.e. a nonlinear transformation of a stationary Gaussian field  $\gamma(\omega, x)$  with zero mean and unit variance:

$$a(\omega, x) = F_a^{-1} \circ \Phi(\gamma(\omega, x)),$$

where  $\Phi$  is the CDF of N(0,1)

• then we can approximate  $\gamma(\omega, x)$  using a truncated KL expansion in terms of Gaussian random parameters  $\{Y_n(\omega)\}_{n=1}^N$  s.t.

$$a_{N}(\omega, x) = F_{a}^{-1} \circ \Phi(\gamma_{N}(\omega, x)) = F_{a}^{-1} \circ \Phi\left(b_{0}(x) + \sum_{n=1}^{N} b_{n}(x)Y_{n}(\omega)\right)$$

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Summary What's coming next?



- Motivation, Probabilistic and Stochastic models, aleatoric and epistemic uncertainty
- Formulating a "Plan of attack" for solving stochastic problems
- An overview of numerical methods for solving SPDEs
- Random variables and random fields
- Stochastic representation of a random field:
  - Discrete case, Singular value decomposition
  - Spectral expansion, Karhunen-Loève expansion
- What's next? Formulating a well-posed SPDE, stochastic regularity
- How to compute a numerical solution  $u_N(\omega, x)$ ?
  - Monte Carlo FEM
  - Stochastic Galerkin FEM
  - Stochastic Collocation FEM
  - convergence analysis
- What happens when N becomes large?
  - curse of dimensionality and sparse representations