

Uncertainty quantification & approximation theory for parameterized (stochastic) PDEs

Part I: Background, motivation, and "tools of the trade"

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Supporting agencies: DOE (ASCR, BES), DOD (AFOSR, DARPA), NSF (CM)

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Summary of the course

Uncertainty quantification & approximation theory for parameterized (stochastic) PDEs $\,$



DAY 1

- Background, motivation, and "tools of the trade": Clayton Webster
- Methods and algorithms: Clayton Webster

DAY 2

- Orthogonal polynomials and best approximation: Hoang Tran
- Oiscrete least squares and compressed sensing techniques: Hoang Tran

DAY 3

- Sparse grid interpolation via global Lagrange polynomials: Clayton Webster
- Sparse grids interpolation via local hierarchical polynomials: Guannan Zhang

DAY 4

- Multilevel sampling and interpolation methods: Guannan Zhang
- Other topics, opportunities, and open discussion: Tran, Webster, and Zhang

Part I Outline

Background, motivation, and "tools of the trade"



- Why parameterized (stochastic) models?
- 2 An overview of uncertainty quantification (UQ)
- Stochastic models
- Tools of the trade: random variables
- 5 Tools of the trade: random processes and random fields
- 6 SVD: Discrete version of the Karhunen-Loève expansion
- The Karhunen-Loève expansion
- Summary of Part I

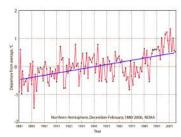
Why parameterized (stochastic) models? A transition to non-deterministic simulations



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.

Predicting future climate changes

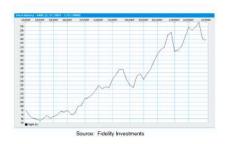


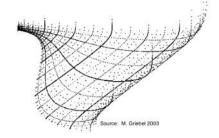


Global warming?



Forecasting financial markets

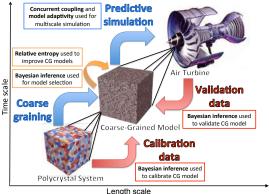




The amount of uncertainty may depend on the number of: economic factors, underlying assets, or the number of time points/time steps, as well as human behaviors, etc.



Modeling and predicting the behavior of large-scale engineered systems



- 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).

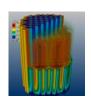
Additional examples

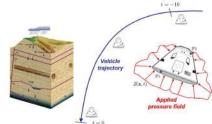


UQ examples of international importance:

- Enhancement of reliability of smart energy grids.
- Development of renewable energy technologies.
- Vulnerability analysis of water and power supplies.
- Understanding complex biological networks.
- Modeling unwanted vibrations re-entry vehicles experience when pierce the earths atmosphere.
- Design and licensing of current and future nuclear energy reactors.







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

Q 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)

Sational Laboratory

Worst scenario approaches

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 (QoI) denoted Q(u):

Worst scenario approaches: Typically A is an ϵ -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.

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

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

"Assumption has many advantages. Chiefly these are the same as those of theft over honest toil." – Bertrand Russel



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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 α -cut of C by $C_{\alpha} \stackrel{\mathrm{def}}{=} \{x \in A : \mu_C \geq \alpha\}$ which gives a set characterization of uncertainty
- ullet 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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$$\sum_{\varphi\in\Phi}m(\varphi)=1,\quad m(\emptyset)=0,\quad \text{however } \varphi_1\subset\varphi_2\not\Rightarrow m(\varphi_1)\leq m(\varphi_2)$$

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What is probability theory? Understood as a mathematical theory of a finite measure.

Let $(\Omega,\mathscr{F},\mathbb{P})$ denote a (complete) probability space: Ω is the event space, $\mathscr{F}\subset 2^{\Omega}$ is the σ -algebra and \mathbb{P} is the probability measure, satisfying:

- ② A positive measure $\mathbb{P}: \mathscr{F} \to [0,1]$ which is countably additive, i.e.,

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

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

- applies to aleatoric phenomena, i.e., frequencies of occurance.
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Uncertainty quantification (UQ)

Stochastic/Probabilistic methods



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

$$\mathcal{F}(\omega)[u(x)] = f(x,\omega)$$
 for a.e. $x \in D \subset \mathbb{R}^n$.

Evolution in time may be

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

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

$$du = (\mathcal{F}(\omega)[u] - f(\omega, x, t)) dt + \mathcal{B}(\omega)[u] dW(\omega, t) + \mathcal{P}(\omega)[u] dQ(\omega, t).$$

In this Itô evolution equation, $W(\omega,t)$ is a Wiener process and $Q(\omega,t)$ is a Poisson process.



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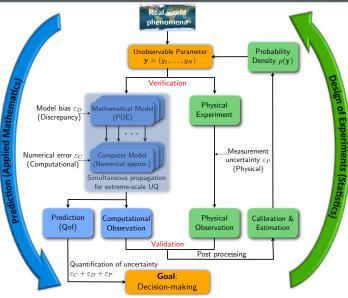
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An example stochastic partial differential equation (SPDE) Pressure distribution in a porous medium



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

$$\left\{ \begin{array}{rcl} -\nabla \cdot (a(x) \, \nabla u(x)) &= f(x) & \text{ in } D \subset \mathbb{R}^n, \\ u(x) &= 0 & \text{ on } \partial D, \end{array} \right.$$

where $n=1,2,3,\ f(x)$ is the source, a(x) describes the permeability and u(x) is the pressure distribution.

Q: What if the input data is random?

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Assumptions and well-posedness



① the solution has realizations in the Banach space $\mathcal{V}\equiv H^1_0(D)$, i.e., $u(\,\cdot\,,\omega)\in\mathcal{V}$ almost surely

$$||u(\cdot,\omega)||_{\mathcal{V}} \le C||f(\cdot,\omega)||_{\mathcal{V}^*},$$

where $\mathcal{V}^* = H^{-1}(D)$.

② the forcing term $f \in L^2_{\mathbb{P}}(\Omega; \mathcal{V}^*)$ is such that the solution u is unique and bounded in $L^2_{\mathbb{P}}(\Omega; \mathcal{V})$, i.e., Banach-valued functions that have finite second moments.

For $q \in \mathbb{N}_+$

$$L^q_{\mathbb{P}}(\Omega; \mathcal{V}) = \left\{ v : \Omega \to \mathcal{V} \mid v \text{ is measurable and } \int_{\Omega} \|v(\omega, \, \cdot\,)\|^q_{\mathcal{V}} \, dP(\omega) < +\infty \right\}.$$

 \bullet $a(\,\cdot\,,\omega)$ uniformly bounded and coercive, i.e., there exists $a_{\min},\,a_{\max}\in(0,+\infty)$ such that

$$\mathbb{P}\left[\omega \in \Omega: \ a(\omega, x) \in (a_{\min}, a_{\max}) \, \forall x \in \overline{D}\right] = 1.$$

Approximate \boldsymbol{u} or some statistical QoI depending on \boldsymbol{u}

Difficulties arise since instead of just asking for u(x), we instead want to approximate the entire stochastic solution $u(x,\omega)$ or some statistical QoI depending on u:

$$\phi_u := \mathbb{E}\left[\phi(u)\right] = \int_{\Omega} \int_{D} \psi(u(x,\omega), x, \omega) \, dx \, dP(\omega)$$

- Moments: $\overline{u} = \mathbb{E}[u](x)$ or $\mathbb{V}ar[u](x) = \mathbb{E}[\widetilde{u}^2](x)$, where $\widetilde{u} = u \overline{u}$
- $\bullet \ \ \mathbf{Probabilities:} \ \mathbb{P}\left\{u \geq u_0\right\} = \mathbb{P}\left[\left\{\omega \in \Omega \,:\, u(\omega) \geq u_0\right\}\right] = \mathbb{E}\left[\chi_{\left\{u \geq u_0\right\}}\right]$
- Statistics of functionals of u:

$$\phi(u) = \int_{\Sigma \subset D} u(x, \, \cdot \,) \, dx,$$

where Σ is a subdomain of interest.

Our goal: to develop highly efficient, robust, and scalable techniques that include uncertainty in the models, and allows us to quantify uncertainty in the outputs while providing reliable and verifiable predictions.



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- Determine an accurate representation for the input stochastic (random) fields can be both simulated and analyzed, e.g., a Karhunen-Loève expansion.
- @ Transform the stochastic problem into a **deterministic** parametric version in $\mathbb{R}^n \times$ an $\infty-$ dimensional space.
- Oesign an adaptive dimensional reduction procedure.
- Design an adaptive discretization procedure using sampling methods, polynomial methods, or a combination of both.
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- **©** Compare the convergence and **complexity** of the approach with existing methods.

Input data $a(x,\omega)$, $f(x,\omega)$, and the solution $u(x,\omega)$ of the SPDE will (more likely) be a random field defined by a set of random variables $\boldsymbol{y}(\omega)=(y_1(\omega),\ldots,y_d(\omega))$, s.t.

$$oldsymbol{y}(\omega):\Omega o \mathcal{U}=\prod_{i=1}^d \mathcal{U}_i\subset \mathbb{R}^d,$$

where $d \in \mathbb{N}$ that can be very large or even infinite.

WLOG assume
$$a(x,\omega)=a(x,{m y}),\ f(x,\omega)=f(x,{m y}),$$
 and the solution $u(x,\omega)=u(x,{m y})$

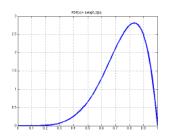
We need to answer the following questions:

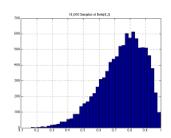
- **1** How to deal with random variables $y_k(\omega)$, $k=1,2,\ldots$?
- 2 How to represent random fields, $a(\cdot, y)$, $f(\cdot, y)$, and $u(\cdot, y)$?

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

$$\begin{split} F_y(\lambda) &:= \mathbb{P}\left[\{ y(\omega) \leq \lambda \} \right] = \int_{\{ y(\omega) \leq \lambda \}} d\mathbb{P}(\omega) = \mathbb{E}\left[\chi_{\{ y(\omega) \leq \lambda \}} \right] \\ &= \int_{-\infty}^{\lambda} \varrho_y(t) \, dt, \end{split}$$

where $\int_{\mathbb{R}} \varrho_y(t) dt = 1$





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

$$\overline{y} = \mathbb{E}[y] = \int_{\Omega} y(\omega) d\mathbb{P}(\omega) = \int_{\mathbb{R}} t \varrho_y(t) dt$$

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 $y(\omega) = (y_k(\omega))_{k=1}^d$, $d \in \mathbb{N}_+$ be a random vector, then the covariance and correlation of two RVs:

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- 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 $\forall \phi_1, \phi_2$, $\mathbb{E}[\phi_1(y_1)\phi_2(y_2)] = \mathbb{E}[\phi_1(y_1)]\mathbb{E}[\phi_2(y_2)]$

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

$$\overline{y} = \mathbb{E}[y] = \int_{\Omega} y(\omega) d\mathbb{P}(\omega) = \int_{\mathbb{R}} t \varrho_y(t) dt$$

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

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$$\mathbb{C}ov[\boldsymbol{y}]_{ij} = \mathbb{C}ov[y_i, y_j]$$

• $\mathbb{C}ov[y]$ is symmetric, nonnegative definite, and has diagonal elements $\mathbb{C}ov[y]_{ii} = \mathbb{V}ar[y_i]$

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

$$\mathbb{C}orr[\boldsymbol{y}] = \Sigma^{-1} \mathbb{C}ov[\boldsymbol{y}] \Sigma$$

- The diagonal entries of $\mathbb{C}orr[y]$ are 1.
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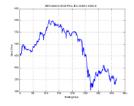
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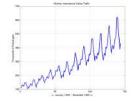
Random processes and fields: generalizations of random variables

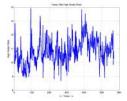
A random process/field (RP/RF) $a(x, y(\omega))$ defined on a probability space $(\Omega, \mathscr{F}, \mathbb{P})$ and indexed by a deterministic domain $D \subset \mathbb{R}^n$, returns a real value

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

Often only second order information - mean and covariance are known







- Mean $\overline{a}(x) = \mathbb{E}[a](x) = \int_{\Omega} a(\,\cdot\,, \boldsymbol{y}(\omega)) d\mathbb{P}(\omega)$ and $\mathbb{V}ar[a](x) = \mathbb{E}[(\widetilde{a})^2](x)$ as a function of x with fluctuation part $\widetilde{a}(x, \boldsymbol{y}(\omega)) = a \overline{a}$
- $\mathbb{P}[a \ge a_0] = \mathbb{P}[\{\omega \in \Omega : a(x, \boldsymbol{y}(\omega)) \ge a_0\}] = \mathbb{E}[\chi_{\{a \ge a_0\}}]$

The covariance covariance function describes the interaction between points 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: $a(x,y(\omega)) \in L^2_{\mathbb{P}}(\Omega;\mathcal{V}).$
- we will **not** focus on white noise approximations which refers to uncorrelated RFs for which: $\overline{a}(x) = 0$ and $\mathbb{C}ov[a](x_1, x_2) = \delta(x_1 x_2)$

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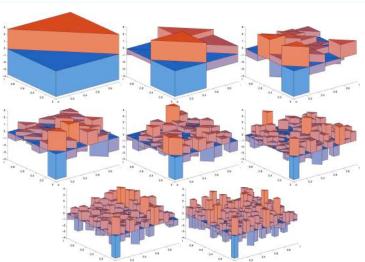
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Discretized white noise Piecewise constant approximations





Discretized white noise over a square subdivided into $2,\dots,512\ \mathrm{triangles}$



Given the mean and covariance of the RF we would like to construct a simple representation which captures this information and used for simulations.

- Ideally, this involves a combination of countably many independent RVs
- Let $\mathcal{U}_n \equiv y_n(\Omega) \subset \mathbb{R}$ be the image of the RV, i.e., $\mathcal{U}_n = [-1,1]$, and $\underset{}{\mathcal{U}} = \prod_{n=1}^d \mathcal{U}_n \subset \mathbb{R}^d$, the image of the random vector $\boldsymbol{y}(\omega) : \Omega \to \mathcal{U}$.
- Let $\varrho : \mathcal{U} \to \mathbb{R}_+$, with $\varrho \in L^{\infty}(\mathcal{U})$ be the joint probability density function (PDF) of $\mathbf{y} = (y_1, \dots, y_d)$, then we want that:

$$\varrho(\boldsymbol{y}) = \prod_{n=1}^d \varrho_n(y_n), \text{ where } \boldsymbol{y} \in \mathcal{U} \text{ and } y_n \in \mathcal{U}_n, \, \forall n$$

- ullet The independence of the d 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 d RVs can be high-dimensional
- In the discrete case KL is similar to ROM and the SVD of a matrix

Finite dimensional noise assumption

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Every (real) m by n matrix A has a singular value decomposition:

$$A = U S V^T$$

where

- U is an m by m orthogonal matrix $(U^T U = I)$;
- ullet 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



$$A = U S V^{T}$$

- ullet r, the number of nonzero diagonal elements in S, is the rank of A.
 - very small non-zeros 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.
- ullet 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:

- ullet 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:
- **1** Minimum L^2 norm:

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

Minimum Frobenius (sum of squares) norm:

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



U and V are natural bases for the input and output of 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.
- ullet The singular value expansion produces an optimal, indexed family of reduced order models of A.



- ullet SVD is the discrete version of the 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.



- ullet 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

SVD example Snowfall at Michigan Tech



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.

<u>Year</u>	<u>Oct</u>	Nov	<u>Dec</u>	<u>Jan</u>	<u>Feb</u>	Mar	<u>Apr</u>	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.atu.edu/elumei/feureitee/enumfell/enu									

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:

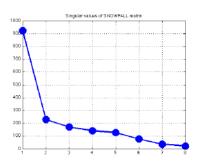
$$A = \left[x^{1890} | x^{1891} | \dots | x^{2010} \right]$$

- Now we determine the SVD decomposition $A = USV^T$.
- ullet 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.
- 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:

$$\frac{s_1^2}{\sqrt{\sum_{i=1}^8 s_i^2}} = 0.87 \quad \frac{s_2^2}{\sqrt{\sum_{i=1}^8 s_i^2}} = 0.05$$

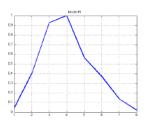
The first pair of modes, u_1 and v_1 , by itself, can approximate the entire dataset with a relative accuracy of 87%.

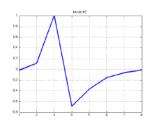


SVD example

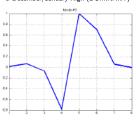
Four strongest snowfall modes for a year



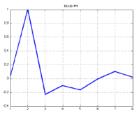




1 December/January High (DOMINANT)



 $2\ \mathsf{More}\ \mathsf{December},\ \mathsf{less}\ \mathsf{later}$

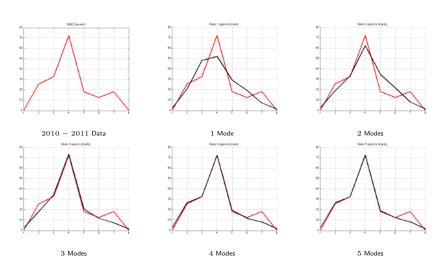


3 February High, less January

4 More November snow

SVD example Approximating 2010-2011 snowfall

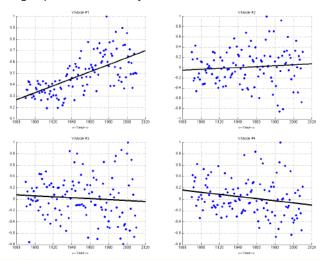




The same kind of approximating is occurring for all 121 sets of data!

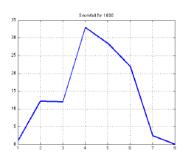


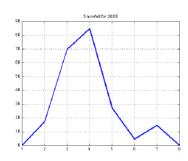
The linear regression line suggests the "December/January High" pattern (upper left) is steadily gaining importance over the years.





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





- Data gathered at discrete places and times is easier to understand than the corresponding continuous cases.
- 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 λ values represent a strength, the u's represent variation in space, and v variation in time.

② 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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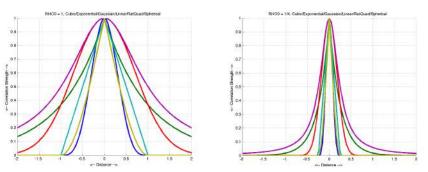
- Given that our data was stored in A, we may think of the matrices AA^T and A^TA as a form of a covariance matrix.
- The singular values $\sqrt{\lambda_i}$ are the square roots of eigenvalues of both these matrices.
- ullet . U contains eigenvectors of the "spatial" covariance matrix AA^T .
- ullet V contains eigenvectors of the "temporal" covariance matrix A^TA .
- Very similar statements will hold for the continuous case.

It will be helpful to keep the discrete case in mind as we briefly discuss the continuous case next.



The Karhunen-Loève expansion of a RF is a Fourier-type series based on the spectral expansion of its covariance function - other names Proper Orthogonal Decomposition (POD), Principle Component Analysis (PCA)

Let $a(x, \mathbf{y}(\omega)) \in L^2_{\mathbb{P}}$ 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 $L_c=1/4$

Properties of the covariance function



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

$$\sum_{i=1}^n \sum_{j=1}^n \mathbb{C}_a(x_i, x_j) v_i v_j \ge 0, \forall (x_1, \dots, x_n) \in D^n \text{ and } (v_1, \dots, v_n) \in \mathbb{R}^n.$$

In matrix notation: $v^T \mathbb{C}_a(x,x) v \geq 0, \ \forall v,x.$

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

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

Observation: If $T_{\mathbb{C}_a} f \in C^0(D) \ \forall f \in L^2(D)$, $\mathbb{C}_a \mapsto T_{\mathbb{C}_a}$ is injective, and $T_{\mathbb{C}_a}$ is compact, symmetric and non-negative definite, then:

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Mercer's (1909) spectral representation of the kernel:

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

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

$$\left[T_{\mathbb{C}_a}b_n\right](x_1)=\int_D\mathbb{C}_a(x_1,x_2)b_n(x_2)dx_2=\frac{\mathbf{\lambda}_n}{\mathbf{\lambda}_n}b_n(x_1),\quad n=1,\dots$$
 with
$$\int_Db_n(x_1)b_m(x_1)dx_1=\delta_{nm}$$

Karhunen-Loève expansion Spectral representation of the kernel

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



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

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

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

Theorem. [Mercer, 1909].

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

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



$$a(x,\omega) = \overline{a}(x) + \sum_{n=1}^{+\infty} b_n(x) \, y_n(\omega),$$
 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, i.e.,

$$\mathbb{E}[y_n] = 0, \, \mathbb{C}ov[y_n, y_m] = \mathbb{E}[y_n y_m] = \delta_{nm}$$

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

$$a(\omega, x) \approx a_d(\omega, x) = \overline{a}(x) + \sum_{n=1}^d b_n(x) y_n(\omega)$$

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



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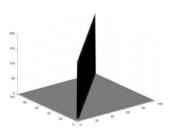
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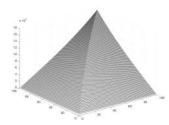
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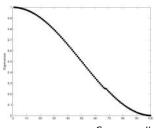
Karhunen-Loève expansion Example: Ornstein-Uhlenbeck process

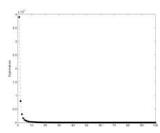






Peaked and smooth covariance functions





Corresponding KL eigenvalues



This truncated expansion corresponds to the Best d-term approximation:

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

- If we truncate using the d largest eigenvalues, we have an optimal in variance expansion in d random varibles.
- i.e., with \mathbb{C}_a continuous, a_d converges uniformly to a (Mercer's Theorem)

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

- "Karhunen-Loève 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 d largest eigenvalues of A^*A :
 - ⇒ Model redcution find a *d*-sparse representation

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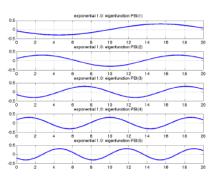
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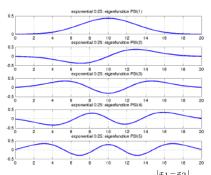
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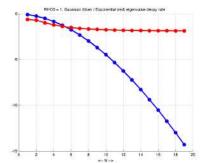
Distributions of the eigenfunctions of 1d exponential kernels, $\mathbb{C}_a(x_1,x_2)=\sigma^2 e^{-\frac{|x_1-x_2|}{L_c}}$, for correlation lengths $L_c=1$ and $L_c=1/4$

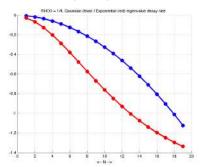
Karhunen-Loève expansion Examples of eigenvalues



$$\text{exp.: } \mathbb{C}_a(x_1,x_2) = \sigma^2 e^{-\frac{\|x_1-x_2\|_1}{L_c^2}}$$

Gaussian:
$$\mathbb{C}_a(x_1,x_2) = \sigma^2 e^{-\frac{\|x_1-x_2\|^2}{L_c^2}}$$





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 ⇒ 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.

Convergence of the spectrum

- 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:

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

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

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

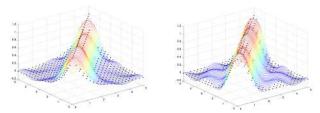
Remark: similar to SVD - if one wants the relative error (in the variance) less than some tolerance δ , i.e., $\|a-a_d\|^2 \leq \delta \|a\|$, then choose d s.t.

$$\sum_{n=d+1}^{\infty} \frac{\lambda_n}{\lambda_n} \le \delta \sum_{n=1}^{\infty} \frac{\lambda_n}{\lambda_n}.$$

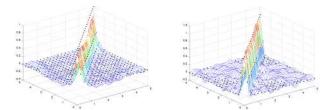
Karhunen-Loève expansion

Approximating 90% variance with Gaussian kernel





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



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

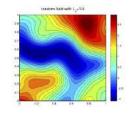


Although the y_n 's are uncorrelated, in general they are not independent:

ullet Gaussian RVs are uncorrelated \Longleftrightarrow independent, i.e., $\varrho({m y}) = \prod_{n=1}^d \varrho_n(y_n)$.

Gaussian random fields: For every $x \in D$, $a(x, \cdot) \sim N(\mu(x), \varrho(x, x'))$

- Characterized by mean $\mu(x)$ and covariance $\varrho(x,x')=\mathbb{E}[\hat{a}(x,\cdot)\hat{a}(x',\cdot)]$, where $\hat{a}(x,\cdot)=a(x,\cdot)-\mu(x)$
- e.g. exponential $\varrho(x,x')=\sigma^2 e^{-\frac{\left|x-x'\right|}{L_c}}$, Gaussian $\varrho(x,x')=\sigma^2 e^{-\frac{\left(x-x'\right)^2}{L_c}}$, etc.



Remark: The diffusion coefficient can not be a Gaussian field (finite probability becomes negative). Use nonlinear transformations, e.g., lognormal model

$$a(x,\omega) = a_{min} + e^{\gamma(x,\omega)}, \qquad \gamma \sim N(\mu(\cdot), \varrho(\cdot, \cdot)),$$

i.e., a truncated Gaussian random field

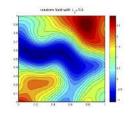


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Let $a(x,\omega)$ be a given stationary non-Gaussian random field with a given (or approximated) marginal CDF F_a , then:

lacktriangle one can translation process, i.e., a nonlinear transformation of a stationary Gaussian field $\gamma(x,\omega)$ with zero mean and unit variance:

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

where Φ is the CDF of N(0,1).

② we can approximate $\gamma(x,\omega)$ using a truncated Karhunen-Loève expansion in terms of Gaussian random parameters $\{y_n(\omega)\}_{n=1}^d$ s.t.

$$a_d(x,\omega) = F_a^{-1} \circ \Phi(\gamma_d(x,\omega)) = F_a^{-1} \circ \Phi\left(b_0(x) + \sum_{n=1}^d b_n(x)y_n(\omega)\right).$$



- Motivation, probabilistic, and stochastic models, aleatoric and epistemic uncertainty
- Formulating a "Plan of attack" for solving stochastic problems
- Random variables and random fields
- Stochastic representation of a random field:
 - Discrete case, Singular value decomposition
 - Spectral expansion, Karhunen-Loève expansion
- What's next? Well-posed (stochastic) parameterized PDEs, regularity, and numerical approximations:
- How to compute a numerical solution u(x, y), where $y \in \mathcal{U} \in \mathbb{R}^d$?
 - Monte Carlo FEMs
 - stochastic Galerkin FEMs
 - stochastic collocation FEMs
 - convergence and complexity analysis
- What happens when d becomes large?
 - curse of dimensionality and sparse representations