

# Inverse Problems

## Sommersemester 2022

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## Practical matters

- Lectures on Mondays at 10.15-12.00 in T9/046 by Vesa Kaarnioja.
- Exercises on Mondays at 12.15-14.00 in T9/046 by Vesa Kaarnioja starting next week.
- Weekly exercises published after the lecture. Please return your written solutions to Vesa either by email ([vesa.kaarnioja@fu-berlin.de](mailto:vesa.kaarnioja@fu-berlin.de)) or at the beginning of the exercise session in the following week.
- The course grade is determined as a weighted average of the exercise points (25%) and the course exam (75%).
- 50% completion of all tasks ensures a passing grade, 90% completion of all tasks ensures the best grade.

## Course contents

- The first part of the course will cover classical variational regularization methods. We will follow Chapters 1–4 in
  - J. Kaipio and E. Somersalo (2005). *Statistical and Computational Inverse Problems*. Springer, New York, NY.
- Second part of the course will cover Bayesian inverse problems. We will follow the texts
  - D. Sanz-Alonso, A. M. Stuart, and A. Taeb (2018). *Inverse Problems and Data Assimilation*. <https://arxiv.org/abs/1810.06191>
  - J. Kaipio and E. Somersalo (2005). *Statistical and Computational Inverse Problems*. Springer, New York, NY.
  - D. Calvetti and E. Somersalo (2007). *Introduction to Bayesian Scientific Computing: Ten Lectures on Subjective Computing*. Springer, New York, NY.

# What is an inverse problem?

- **Forward problem:** Given known causes (initial conditions, material properties, other model parameters), determine the effects (data, measurements).
- **Inverse problem:** Observing the effects (noisy data), recover the cause.

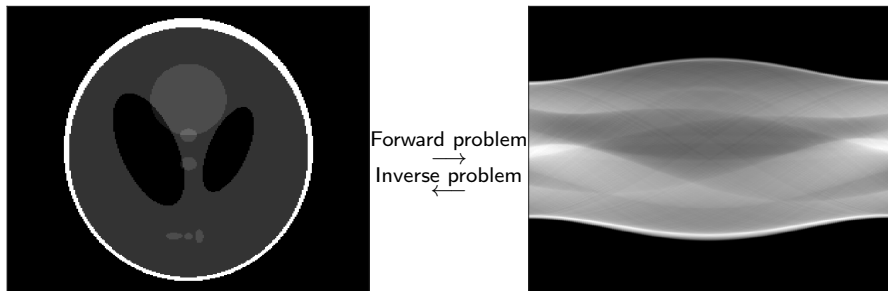


Figure: Computerized tomography (CT)



Figure: Image deblurring (deconvolution)

$$y = (K * f)(x) = \int_{\mathbb{R}^2} K(x - x')f(x') dx'$$

## Introduction: What is an inverse problem?

We consider the indirect measurement of an unknown physical quantity  $x \in X$ . The measurement  $y \in Y$  is related to the unknown by a physical or mathematical *model*

$$y = F(x), \quad (1)$$

where  $F: X \rightarrow Y$  is called the *forward mapping*.

- Computing  $y$  for a given  $x$  is called the *forward problem*.
- Finding  $x$  for a given measurement  $y$  (the *data*) is called the *inverse problem*.

The inverse problem is often ill-posed, making it more difficult than the corresponding direct problem.

A problem is called *well-posed* (according to Hadamard), if

- (a) a solution exists,
- (b) the solution is unique, and
- (c) the solution depends continuously on the data.

If one or more of these conditions are violated, the problem is called *ill-posed*.

Some examples of ill-posed inverse problems are X-ray tomography, image deblurring, the inverse heat equation, and electrical impedance tomography (EIT).

The ill-posedness of an inverse problem poses a challenge because usually, errors are present in the measurements. Incorporating these into model (1) in the form of additive *noise*  $\eta$  leads to a more realistic model

$$y = F(x) + \eta.$$

The violation of the above conditions leads to various difficulties.

- If condition (a) is violated, i.e., if the image  $\text{Ran}(F)$  of  $F$  does not cover the whole space  $Y$ , then there may not exist a solution to  $F(x) = y$  for noisy data  $y = F(x^\dagger) + \eta$  created by a ground truth  $x^\dagger$ , although a solution exists for noise free data  $y = F(x^\dagger)$ , since  $\eta$  does not need to lie in  $\text{Ran}(F)$ .
- If condition (c) is violated, then the solution to  $F(x) = y$  for noisy data  $y = F(x^\dagger) + \eta$  may be far away from the solution for noise free data  $y = F(x^\dagger)$ , even if  $F$  is invertible and the noise  $\eta$  is small, due to the discontinuity of  $F^{-1}$ .



### Example.

The deblurring (or deconvolution) problem of recovering an input signal  $x$  from an observed signal  $y$  (possibly contaminated by noise) occurs in many imaging as well as image and signal processing applications. The mathematical model is

$$y(t) = \underbrace{\int_{-\infty}^{\infty} a(t-s)x(s)ds}_{=:(a*x)(t)},$$

where the function  $a$  is known as the blurring kernel.

If  $\hat{a}$  is “nice”, we can use the Fourier transform together with the convolution theorem to solve the problem analytically:

$$\begin{aligned} y(t) = (a * x)(t) &\Leftrightarrow \hat{y}(\xi) = \hat{a}(\xi)\hat{x}(\xi) \Leftrightarrow \hat{x}(\xi) = \frac{\hat{y}(\xi)}{\hat{a}(\xi)} \\ \Leftrightarrow x(t) = \mathcal{F}^{-1}\left\{\frac{\hat{y}}{\hat{a}}\right\}(t) &= \frac{1}{2\pi} \int_{-\infty}^{\infty} e^{it\xi} \frac{\hat{y}(\xi)}{\hat{a}(\xi)} d\xi. \end{aligned}$$

Let  $x_{\text{exact}}$  denote the solution to this problem with exact, noiseless data.

However, if we can only observe noisy measurements, we must consider

$$y(t) = (a * x)(t) + \eta(t) \quad \Leftrightarrow \quad \hat{y}(\xi) = \hat{a}(\xi)\hat{x}(\xi) + \hat{\eta}(\xi).$$

The solution formula from the previous slide gives (in the Fourier side)

$$\hat{x}(\xi) = \frac{\hat{y}(\xi)}{\hat{a}(\xi)} = \hat{x}_{\text{exact}}(\xi) + \frac{\hat{\eta}(\xi)}{\hat{a}(\xi)};$$

then we apply the inverse Fourier transform on both sides. However, this reconstruction might not be well-defined and it is typically not stable, i.e., it does not depend continuously on the data  $y$ . The kernel  $a$  usually decreases exponentially (or has compact support). A typical example is a Gaussian kernel

$$a(t) = \frac{1}{2\pi\alpha^2} \exp\left(-\frac{t^2}{2\alpha^2}\right)$$

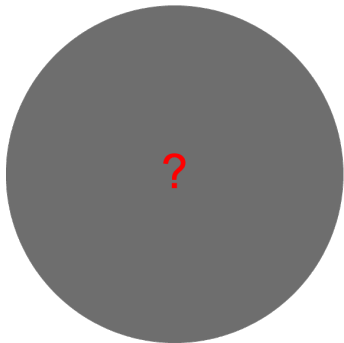
for some  $\alpha > 0$ .

By the Plancherel theorem,  $\hat{a} \in L^2(\mathbb{R})$  and

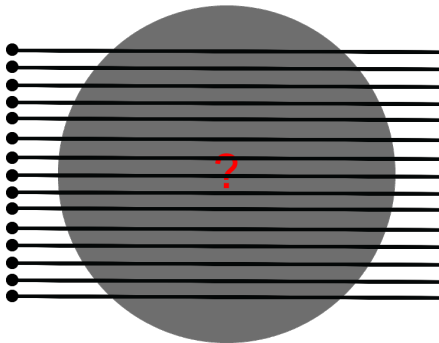
$$\int_{-\infty}^{\infty} |a(t)|^2 dt = \int_{-\infty}^{\infty} |\hat{a}(\xi)|^2 d\xi$$

if  $a \in L^2(\mathbb{R})$ . This implies in particular that  $\hat{a}(\xi) \rightarrow 0$  as  $|\xi| \rightarrow \infty$ . As a consequence, high frequencies  $\hat{\eta}(\xi)$  of the noise get amplified arbitrarily strong in  $\hat{x}$ . Thus, even the presence of small noise can lead to large changes in the reconstruction.

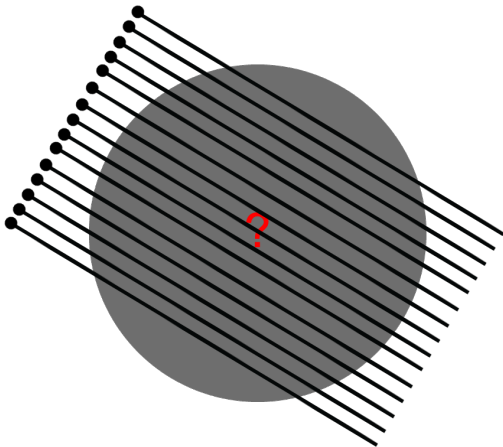
## Case study: parallel-beam X-ray tomography



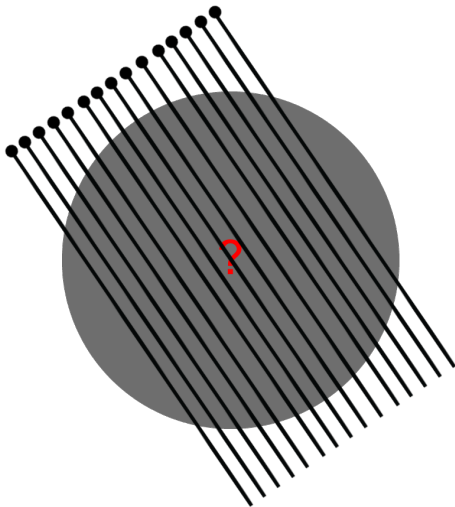
## Case study: parallel-beam X-ray tomography



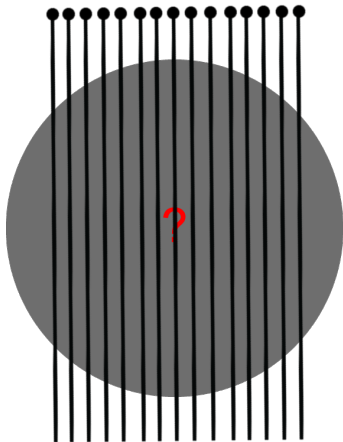
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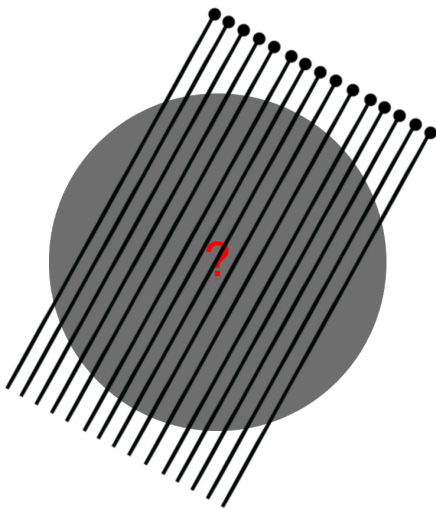


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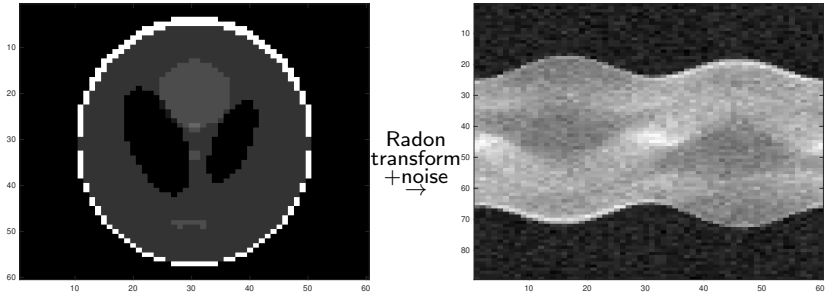




## Case study: parallel-beam X-ray tomography



Let us consider the following phantom (bottom left), which we use to simulate measurements taken from 60 angles contaminated with 5 % Gaussian noise (sinogram on the bottom right). Inverse problem: use the sinogram data (X-ray images taken from the different directions) to reconstruct the internal structure of the physical body (i.e., the phantom).

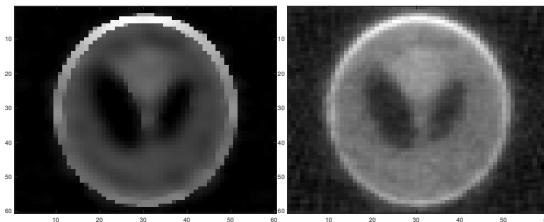


Technical (but important) note: to avoid the so-called inverse crime, the measurements for the inversion on the following page were generated using a higher resolution phantom.

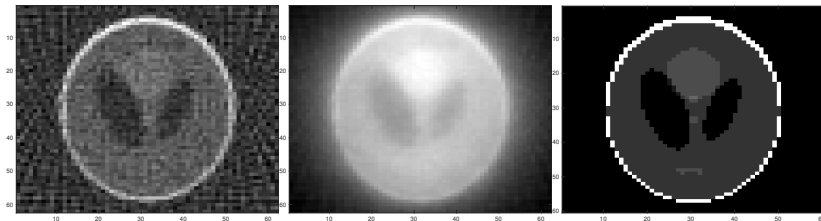
Formation of a CT sinogram (Samuli Siltanen):

[https://www.youtube.com/watch?v=q7Rt\\_OY\\_7tU](https://www.youtube.com/watch?v=q7Rt_OY_7tU)

Reconstructions  $\arg \min_x \{ \|Ax - m\|^2 + \mathcal{R}(x) \}$  from noisy measurements  $m$  with some selected penalty terms  $\mathcal{R}$  are given immediately below.



Left: reconstruction with total variation regularization. Right: same with Tikhonov regularization. Some other reconstructions for comparison (and the target phantom).



Left: filtered back projection. Middle: unfiltered back projection. Right: ground truth.

- Successful solution of inverse problems requires specially designed algorithms that can tolerate errors in measured data.
- How to incorporate all possible prior and expert knowledge about the possible solutions when solving inverse problems?
- The statistical approach to inverse problems aims to quantify how uncertainty in the data or model affects the solutions obtained in problems.

## Separable Hilbert space

A vector space  $H$  is called a *real inner product space* if there exists a mapping  $\langle \cdot, \cdot \rangle: H \times H \rightarrow \mathbb{R}$  satisfying

- $\langle x, y \rangle = \langle y, x \rangle$  for all  $x, y \in H$ ;
- $\langle ax_1 + bx_2, y \rangle = a\langle x_1, y \rangle + b\langle x_2, y \rangle$  for all  $x_1, x_2, y \in H$  and  $a, b \in \mathbb{R}$ ;
- $\langle x, x \rangle \geq 0$ , where equality holds iff  $x = 0$ .

Moreover,  $H$  is a *Hilbert space* if

- $H$  is *complete* with respect to the norm  $\|x\| = \sqrt{\langle x, x \rangle}$ ,  $x \in H$ ,

and it is said to be *separable* if, additionally,

- there exists a *countable orthonormal basis*  $\{\phi_n\}$  of  $H$  with respect to the inner product  $\langle \cdot, \cdot \rangle$ , that is,

$$\langle \phi_j, \phi_k \rangle = \delta_{j,k} \quad \text{and} \quad x = \sum_n \langle x, \phi_n \rangle \phi_n \quad \text{for all } x \in H.$$

Let  $H_1$  and  $H_2$  be real Hilbert spaces and let  $A: H_1 \rightarrow H_2$  be a continuous linear operator.

### Definition

The kernel (or null space) of operator  $A$  is defined as

$$\text{Ker}(A) := \{x \in H_1 \mid Ax = 0\}.$$

The range (or image) of operator  $A$  is defined as

$$\text{Ran}(A) := \{y \in H_2 \mid y = Ax, x \in H_1\}.$$

- A linear operator  $A$  is continuous if and only if there exists a constant  $C > 0$  such that  $\|Ax\| \leq C\|x\|$  for all  $x \in H_1$ .
- There exists a unique adjoint operator  $A^*: H_2 \rightarrow H_1$  defined by

$$\langle Ax, y \rangle = \langle x, A^*y \rangle \quad \text{for all } x \in H_1, y \in H_2.$$

(This is a consequence of the Riesz representation theorem.)

- $\text{Ker}(A)$  is a *closed* subspace of  $H_1$ , and  $\text{Ran}(A)$  is a subspace of  $H_2$ .

Let  $H$  be a real Hilbert space.

### Definition

Two elements  $x, y \in H$  are said to be *orthogonal* if  $\langle x, y \rangle = 0$ .

Let  $X \subset H$  be a subset. The orthogonal complement of  $X$  in  $H$  is defined as

$$X^\perp := \{y \in H \mid \langle x, y \rangle = 0 \text{ for all } x \in X\}.$$

- For any set  $X \subset H$ ,  $X^\perp$  is a closed subspace of  $H$  and  $X \subset (X^\perp)^\perp$ .
- If  $X$  is a non-closed subspace, then  $(X^\perp)^\perp = \overline{X}$ .
- If  $X$  is a closed subspace, then  $X = (X^\perp)^\perp$ . In this case, there exists the *orthogonal decomposition*

$$H = X \oplus X^\perp,$$

which means that every element  $y \in H$  can be uniquely represented as

$$y = x + x^\perp, \quad x \in X, \quad x^\perp \in X^\perp.$$

- If  $X \subset H$  is a closed subspace, the mapping  $P_X: H \rightarrow X, y \mapsto x$ , is an *orthogonal projection*, i.e.,  $P_X^2 = P_X$  and  $\text{Ran}(P_X) \perp \text{Ran}(I - P_X)$ .

## Lemma

Let  $X \subset H$  be a closed subspace. The orthogonal projection  $P_X: H \rightarrow X$  satisfies the following properties:

- $P_X$  is linear;
- $P_X$  is self-adjoint:  $P_X^* = P_X$ ;
- $\|P_X\| = 1$  if  $X \neq \{0\}$ ;
- $I - P_X = P_{X^\perp}$ ;
- $\|y - P_X y\| \leq \|y - z\|$  for all  $z \in X$ ;
- $z = P_X y$  iff  $z \in X$  and  $y - z \in X^\perp$ .



## Proposition

Let  $H_1$  and  $H_2$  be Hilbert spaces and  $A: H_1 \rightarrow H_2$  a continuous linear operator. Then

$$H_1 = \text{Ker}(A) \oplus (\text{Ker}(A))^\perp = \text{Ker}(A) \oplus \overline{\text{Ran}(A^*)},$$
$$H_2 = \overline{\text{Ran}(A)} \oplus (\text{Ran}(A))^\perp = \overline{\text{Ran}(A)} \oplus \text{Ker}(A^*).$$

*Proof.*  $H_1 = \overline{\text{Ker}(A) \oplus (\text{Ker}(A))^\perp}$  and  $H_2 = \overline{\text{Ran}(A) \oplus (\text{Ran}(A))^\perp} = \overline{\text{Ran}(A)} \oplus (\text{Ran}(A))^\perp$  follow immediately from the previous discussion.<sup>†</sup> The claim

$$(\text{Ran}(A))^\perp = \text{Ker}(A^*) \quad (2)$$

follows immediately by observing that  $x \in \text{Ker}(A^*)$  iff

$$0 = \langle A^*x, y \rangle = \langle x, Ay \rangle \quad \text{for all } y \in H_1.$$

The claim  $(\text{Ker}(A))^\perp = \overline{\text{Ran}(A^*)}$  follows by applying (2) with  $A$  replaced by  $A^*$ . □

<sup>†</sup>We use the fact that  $\overline{X^\perp} = X^\perp$  for any subspace  $X$  of  $H$ . The inclusion  $\overline{X^\perp} \subset X^\perp$  is trivial. To see the other direction, let  $x \in X^\perp \Rightarrow \langle x, y \rangle = 0$  for all  $y \in X$ . Let  $z \in \overline{X^\perp}$  be arbitrary. We can write  $z = \lim_{n \rightarrow \infty} z_n$  for a sequence  $\{z_n\} \subset X^\perp$ . Hence,  $\langle x, z \rangle = \langle x, \lim_{n \rightarrow \infty} z_n \rangle = \lim_{n \rightarrow \infty} \langle x, z_n \rangle = 0$ . Since  $z \in \overline{X^\perp}$  was arbitrary,  $\overline{X^\perp} = X^\perp$ .

# Fredholm equation

Let us formalize the problem that we will concentrate on during the first part of the course.

Let  $H_1$  and  $H_2$  be separable real Hilbert spaces and let  $A: H_1 \rightarrow H_2$  be a *compact*<sup>†</sup> linear operator. We are interested in finding  $x \in H_1$  such that

$$y = Ax,$$

where  $y \in H_2$  is given.

## Examples:

- $H_1 = H_2 = L^2(a, b)$ .
- $H_1 = \mathbb{R}^n$ ,  $H_2 = \mathbb{R}^m$ , and  $A \in \mathbb{R}^{m \times n}$ .

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<sup>†</sup>A continuous (bounded) linear operator  $A: H_1 \rightarrow H_2$  between Hilbert spaces  $H_1$  and  $H_2$  is said to be compact if the sets  $\overline{A(U)} \subset H_2$  are compact for every bounded set  $U \subset H_1$ .

## Singular value decomposition of a compact operator

Let us assume that  $H_1$  and  $H_2$  are separable real Hilbert spaces and let  $A: H_1 \rightarrow H_2$  be a compact linear operator.

Then there exist (possibly countably infinite) orthonormal sets of vectors  $\{v_n\} \subset H_1$  and  $\{u_n\} \subset H_2$ , and a sequence of positive numbers  $\lambda_1 \geq \lambda_2 \geq \lambda_3 \geq \dots \geq 0$  with  $\lim_{n \rightarrow \infty} \lambda_n = 0$  in the countably infinite case such that

$$Ax = \sum_n \lambda_n \langle x, v_n \rangle u_n \quad \text{for all } x \in H_1. \quad (3)$$

In particular,

$$\overline{\text{Ran}(A)} = \overline{\text{span}\{u_n\}} \quad \text{and} \quad (\text{Ker}(A))^\perp = \overline{\text{span}\{v_n\}}.$$

The system  $(\lambda_n, v_n, u_n)$  is called a *singular system* of  $A$ , and (3) is a *singular value decomposition* (SVD) of  $A$ .

## Singular value decomposition of matrices: $H_1 = \mathbb{R}^n$ and $H_2 = \mathbb{R}^m$

Let  $H_1 = \mathbb{R}^n$  and  $H_2 = \mathbb{R}^m$ , meaning that

$$y = Ax$$

is a matrix equation with  $A \in \mathbb{R}^{m \times n}$ ,  $x \in \mathbb{R}^n$ , and  $y \in \mathbb{R}^m$ .

Since this operator has finite rank ( $\text{rank}(A) := \dim \text{Ran}(A) < \infty$ ), we have

$$Ax = \sum_{j=1}^p \lambda_j (x^T v_j) u_j, \quad p := \text{rank}(A) \leq \min\{n, m\},$$

where  $\{v_j\}_{j=1}^p \subset \mathbb{R}^n$  and  $\{u_j\}_{j=1}^p \subset \mathbb{R}^m$  are sets of orthonormal vectors and  $\{\lambda_j\}_{j=1}^p$  are positive numbers such that  $\lambda_j \geq \lambda_{j+1}$ .

It is possible to complete the sequences of (orthonormal) singular vectors  $\{v_j\}_{j=1}^p \subset \mathbb{R}^n$  and  $\{u_j\}_{j=1}^p \subset \mathbb{R}^m$  with complementary orthonormal vectors  $\{v_j\}_{j=p+1}^n$  and  $\{u_j\}_{j=p+1}^m$  such that  $\{v_j\}_{j=1}^n$  forms an orthonormal basis for  $\mathbb{R}^n$  and  $\{u_j\}_{j=1}^m$  forms an orthonormal basis for  $\mathbb{R}^m$ . This can be done, e.g., using the Gram–Schmidt process.

Define the matrices

$$V = [v_1, \dots, v_n] \in \mathbb{R}^{n \times n},$$
$$U = [u_1, \dots, u_m] \in \mathbb{R}^{m \times m}.$$

Due to the orthonormality of  $\{v_j\}$  and  $\{u_j\}$ , the matrices  $V$  and  $U$  are orthogonal:

$$V^T V = V V^T = I \quad \text{and} \quad U^T U = U U^T = I.$$

Next, we define the matrix  $\Lambda \in \mathbb{R}^{m \times n}$  as follows:

$$\Lambda = \left( \begin{array}{ccc|c} \lambda_1 & & & \\ & \ddots & & \\ & & \lambda_m & \\ \hline & & & O_{m \times (n-m)} \end{array} \right) \quad \text{if } m < n,$$

$$\Lambda = \left( \begin{array}{ccc} \lambda_1 & & \\ & \ddots & \\ & & \lambda_n \\ \hline & & & O_{(m-n) \times n} \end{array} \right) \quad \text{if } m > n,$$

and  $\Lambda = \text{diag}(\lambda_1, \dots, \lambda_m)$  if  $m = n$ .

It is simple to check that

$$Ax = \sum_{j=1}^p \lambda_j u_j v_j^T x = U \Lambda V^T x \quad \text{for all } x \in \mathbb{R}^n,$$

which yields the matrix singular value decomposition (SVD)

$$A = U \Lambda V^T.$$

*This is the decomposition MATLAB returns with the command `svd`.*

Note that in the matrix SVD, the singular values  $\{\lambda_j\}_{j=1}^{\min\{m,n\}}$  are *non-negative* and

$$\text{Ran}(A) = \text{span}\{u_j \mid 1 \leq j \leq p\},$$

$$\text{Ker}(A) = \text{span}\{v_j \mid p + 1 \leq j \leq n\},$$

$$(\text{Ran}(A))^\perp = \text{span}\{u_j \mid p + 1 \leq j \leq m\},$$

$$(\text{Ker}(A))^\perp = \text{span}\{v_j \mid 1 \leq j \leq p\},$$

where  $p = \text{rank}(A) = \max_{1 \leq k \leq \min\{m,n\}} \{k \mid \lambda_k > 0\}$ .

## Solvability of $y = Ax$

Let us assume that  $H_1$  and  $H_2$  are separable real Hilbert spaces and let  $A: H_1 \rightarrow H_2$  be a compact linear operator. Let  $P: H_2 \rightarrow \overline{\text{Ran}(A)}$  be an orthogonal projection. This can be represented using the singular system of  $A$  as

$$Py = \sum_n \langle y, u_n \rangle u_n.$$

### Theorem

The equation  $y = Ax$  has a solution iff

$$y = Py \quad \text{and} \quad \underbrace{\sum_n \frac{1}{\lambda_n^2} |\langle y, u_n \rangle|^2}_{\text{"Picard criterion"}} < \infty.$$

In this case, the solution is of the form

$$x = x_0 + \sum_n \frac{1}{\lambda_n} \langle y, u_n \rangle v_n \quad \text{for arbitrary } x_0 \in \text{Ker}(A).$$



*Proof.* “ $\Rightarrow$ ” Suppose that  $y = Ax$  has a solution  $x \in H_1$ . This implies that  $y \in \text{Ran}(A)$  (thus  $y = Py$ ) and, moreover,

$$\begin{aligned}\langle y, u_j \rangle &= \langle Ax, u_j \rangle = \langle x, A^* u_j \rangle = \lambda_j \langle x, v_j \rangle \\ \Rightarrow \sum_n \frac{1}{\lambda_n^2} |\langle y, u_n \rangle|^2 &= \sum_n |\langle x, v_n \rangle|^2 \stackrel{\text{Bessel inequ.}}{\leq} \|x\|^2 < \infty.\end{aligned}$$

“ $\Leftarrow$ ” Next, suppose that  $y = Py$  and the Picard criterion hold and define  $x := x_0 + \sum_n \lambda_n^{-1} \langle y, u_n \rangle v_n$ , where  $x_0 \in \text{Ker}(A)$  is arbitrary. We obtain

$$Ax = Ax_0 + \sum_n \frac{1}{\lambda_n} \langle y, u_n \rangle Av_n = \sum_n \langle y, u_n \rangle u_n = Py = y. \quad \square$$

*Remark.* In the above proof, it is helpful to note that if  $A$  has the SVD

$$Ax = \sum_n \lambda_n \langle x, v_n \rangle u_n,$$

then its adjoint  $A^*$  has the SVD

$$A^*y = \sum_n \lambda_n \langle y, u_n \rangle v_n.$$

Note that for any  $x \in H_1$ , we have

$$\|Ax - y\|^2 = \|Ax - Py\|^2 + \|(I - P)y\|^2 \geq \|(I - P)y\|^2.$$

Hence, if  $y$  has a nonzero component in the subspace orthogonal to the range of  $A$ , the equation  $Ax = y$  cannot be satisfied exactly. Thus, the best we can do is to solve the projected equation

$$Ax = PAx = Py.$$

However, there is in general no guarantee that the Picard criterion

$$\sum_n \frac{1}{\lambda_n^2} |\langle Py, u_n \rangle|^2 < \infty$$

is satisfied for a general  $y \in H_2$  if  $\text{rank}(A) = \dim \text{Ran}(A) = \infty$ .

## Truncated singular value decomposition (TSVD)

Let us define a family of finite-dimensional orthogonal projections by

$$P_k: H_2 \rightarrow \text{span}\{u_1, \dots, u_k\}, \quad y \mapsto \sum_{n=1}^k \langle y, u_n \rangle u_n.$$

By the orthogonality of  $\{u_n\}$ ,

$$P(P_k y) = \sum_n \langle P_k y, u_n \rangle u_n = \sum_{n=1}^k \langle y, u_n \rangle u_n = P_k y$$

and

$$\sum_n \frac{1}{\lambda_n^2} |\langle P_k y, u_n \rangle|^2 = \sum_{n=1}^k \frac{1}{\lambda_n^2} |\langle y, u_n \rangle|^2 < \infty.$$

Note that  $k \leq \text{rank}(A)$  if  $\text{rank}(A) < \infty$ .

It follows that the problem

$$Ax = P_k y \quad (4)$$

is always solvable. Taking on both sides the inner product with  $u_n$ , we find that

$$\lambda_n \langle x, v_n \rangle = \begin{cases} \langle y, u_n \rangle, & 1 \leq n \leq k \\ 0, & n > k. \end{cases}$$

Hence the solutions to (4) are given by

$$x_k = x_0 + \sum_n \frac{1}{\lambda_n} \langle P_k y, u_n \rangle v_n = x_0 + \sum_{n=1}^k \frac{1}{\lambda_n} \langle y, u_n \rangle v_n \in H_1$$

for any  $x_0 \in \text{Ker}(A)$ . Observe that since for increasing  $k$ ,

$$\|Ax_k - Py\|^2 = \|(P - P_k)y\|^2 \xrightarrow{k \rightarrow \infty} 0,$$

the residual of the projected equation can be made arbitrarily small.

Finally, to remove the ambiguity of the sought solution due to the possible noninjectivity of  $A$ , we select  $x_0 = 0$ . This choice minimizes the norm of  $x_k$  since, by orthogonality,

$$\|x_k\|^2 = \|x_0\|^2 + \sum_{n=1}^k \frac{1}{\lambda_n^2} |\langle y, u_n \rangle|^2.$$

## Definition

Let  $H_1$  and  $H_2$  be separable real Hilbert spaces and let  $A: H_1 \rightarrow H_2$  be a compact linear operator with a singular system  $(\lambda_n, v_n, u_n)$ . By the truncated SVD approximation (TSVD) of the problem  $Ax = y$ , we mean the problem of finding  $x \in H_1$  such that

$$Ax = P_k y, \quad x \perp \text{Ker}(A)$$

for some  $k \geq 1$ .

## Theorem

*The solution to the TSVD problem has a unique solution  $x_k$ , called the truncated SVD (TSVD) solution, given by*

$$x_k = \sum_{n=1}^k \frac{1}{\lambda_n} \langle y, u_n \rangle v_n.$$

*The TSVD solution satisfies*

$$\|Ax_k - y\|^2 = \|(I - P)y\|^2 + \|(P - P_k)y\|^2 \xrightarrow{k \rightarrow \infty} \|(I - P)y\|^2.$$

## Truncated SVD for a matrix $A \in \mathbb{R}^{m \times n}$

The truncated SVD solution, i.e., solution of

$$Ax = P_k y \quad \text{and} \quad x \perp \text{Ker}(A), \quad 1 \leq k \leq p := \text{rank}(A),$$

where  $P_k: \mathbb{R}^m \rightarrow \text{span}\{u_1, \dots, u_k\}$  is an orthogonal projection, is given by

$$x_k = \sum_{j=1}^k \frac{1}{\lambda_j} \langle y, u_j \rangle v_j = \sum_{j=1}^k \frac{1}{\lambda_j} v_j (u_j^T y) = V \Lambda_k^\dagger U^T y,$$

where  $A$  has the SVD  $A = U \Lambda V^T$  and we define

$$\Lambda_k^\dagger = \begin{pmatrix} 1/\lambda_1 & 0 & \cdots & 0 & \cdots & 0 \\ 0 & 1/\lambda_2 & & & & \vdots \\ \vdots & & \ddots & & & \\ & & & 1/\lambda_k & & \\ & & & & 0 & \\ \vdots & & & & & \ddots \\ 0 & \cdots & & & & \cdots & 0 \end{pmatrix} \in \mathbb{R}^{n \times m},$$

where  $\lambda_1 \geq \cdots \geq \lambda_p > 0$  are the singular values of  $A$  (i.e., diagonal of  $\Lambda$ ).

## Moore–Penrose pseudoinverse

For the largest possible cut-off  $k = p = \text{rank}(A)$ , the matrix

$$A^\dagger := A_p^\dagger = V\Lambda_p^\dagger U^T =: V\Lambda^\dagger U^T$$

is called the *Moore–Penrose pseudoinverse*. It follows from the above that  $x^\dagger = A^\dagger y$  is the solution of the projected (matrix) equation

$$Ax = Py,$$

where  $P: \mathbb{R}^m \rightarrow \text{Ran}(A)$  is the orthogonal projection.

The solution  $x^\dagger = A^\dagger y$  is called the *minimum norm solution* of the problem  $y = Ax$  since

$$\|A^\dagger y\| = \min\{\|x\| : \|Ax - y\| = \|(I - P)y\|\},$$

where  $P$  is the projection onto the range of  $A$ . The minimum norm solution is the solution that minimizes the residual error and has the minimum norm.



Since the smallest singular value  $\lambda_p$  is extremely small in inverse problems, the use of the pseudoinverse is usually very sensitive to inaccuracies in the data  $y$ .

Spectral regularization using TSVD, i.e., discarding singular values below a certain threshold from the forward model, is a simple and popular technique used to render linear problems less ill-posed while improving the noise robustness of the numerical inversion procedure. More on this next week...