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GENERALIZED HOEFFDING DECOMPOSITION

AND THE (SURPRISING) LINEAR NATURE OF NON-LINEARITY

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Does Hoeffding's functional decomposition hold when the inputs are not mutually independent?

Classical Hoeffding's decomposition: Unique decomposition $G(X) = \sum_{A \in \mathcal{P}_D} G_A(X_A)$ for any square-integrable G(X), where the inputs X are **mutually independent**.

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- Non-perfect <u>functional</u> dependence.
- Non-perfect stochastic dependence.

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However... Achieving this result requires an unusual methodological journey.

<u>In this talk:</u> Mix the fields of probability theory and functional analysis, with a sprinkle of algebraic combinatorics, to generalize Hoeffding's decomposition to dependent inputs.

More context

We're not the first to have worked on this generalization.

(see, e.g., Rabitz and Aliş (1999), Peccati (2004), Hooker (2007), Kuo et al. (2009), and Hart and Gremaud (2018))

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They showed that the generalized decomposition hold, **but under fairly restrictive assumptions on the inputs**.

<u>Our approach:</u> Understand the relationships between these subspaces of \mathbb{L}^2 when the inputs are **not mutually independent**.

Random inputs, black-box model

Let $(\Omega, \mathcal{F}, \mathbb{P})$ be a probability space, and let $(E_1, \mathcal{E}_1), \dots, (E_d, \mathcal{E}_d)$ be standard Borel measurable spaces.

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The random inputs are defined as a measurable mapping (i.e., random element):

$$X: \Omega \rightarrow E$$
,

where $E = \sum_{i=1}^{d} E_i$ is the cartesian product of the d Polish spaces.

(This is just a way to say that $X=(X_1,\ldots,X_d)$ is not necessarily \mathbb{R}^d -valued)

Remark. We are mainly going to treat X as a function: although its law is well-defined, we don't really need to control it directly.

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Let $G: E \to \mathbb{R}$ be a **black-box model**, and denote by G(X) the **random output** (it is a random variable).

Let $D = \{1, ..., d\}$, and denote \mathcal{P}_D the **power-set** of D (i.e., the set of subsets of D).

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For every $A \subset D$, denote by:

- $\sigma_A \subseteq \mathcal{F}$ the σ -algebra generated by X_A ;
- $\sigma_X \subseteq \mathcal{F}$ the σ -algebra generated by X.

And notice that if $B \subseteq A$, then $\sigma_B \subseteq \sigma_A$.

Lemma (Doob-Dynkin). If an \mathbb{R} -valued random variable Y is σ_X -measurable, then there exists some function $f: E \to \mathbb{R}$ such that Y = f(X) a.s.

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Finally, denote by σ_{\emptyset} the \mathbb{P} -trivial σ -algebra, i.e., the σ -algebra that contains every event of \mathcal{F} of probability 0.

Lemma (Kallenberg (2021, Lemma 4.9)). Every σ_{\emptyset} -measurable random variable is a.s. constant. $_{4/40}$

Functional dependence

Assumption 1 (Non-perfect functional dependence). Suppose that:

- $\sigma_{\emptyset} \subset \sigma_i$, i = 1, ..., d (inputs are not constant).
- For $B \subset A$, $\sigma_B \subset \sigma_A$ (inputs add information).
- For every $A, B \in \mathcal{P}_D$, $A \neq B$,

$$\sigma_A \cap \sigma_B = \sigma_{A \cap B}$$
.

This assumption is purely functional: we're just controlling the pre-image of the mappings $(X_A)_{A \in \mathcal{P}_D}$.

Proposition. Suppose that Assumption 1 hold. Then, for any $A, B \in \mathcal{P}_D$ such that $A \cap B \notin \{A, B\}$ (i.e., the sets cannot be subsets of each other), **there is no mapping** T **such that**

$$X_B = T(X_A)$$
 a.e.

In other words, if Assumption 1 hold, then the inputs cannot be functions of each other.

Output space

Recall that $(\Omega, \mathcal{F}, \mathbb{P})$ is our sample space, **and let** \mathcal{G} be a **sub-** σ **-algebra** of \mathcal{F} .

Definition (Lebesgue space). Denote by $\mathbb{L}^2(\mathcal{G})$ the **Lebesgue space** containing every **square-integrable**, \mathbb{R} -valued random variables. It is an (infinite-dimensional) Hilbert space with inner product, $\forall Z_1, Z_2 \in \mathbb{L}^2(\mathcal{G})$:

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 $\mathbb{L}^2(\sigma_X)$ is the space of random outputs: it only contains random variable that can be expressed as functions of X.

For every $A \subset D$, $\mathbb{L}^2(\sigma_A) \subset \mathbb{L}^2(\sigma_X)$ only contains random variables that can be expressed as functions of X_A .

 $\mathbb{L}^2\left(\sigma_\emptyset\right)$ only contains a.s constants.

Generated Lebesgue subspaces

Theorem (Sidák (1957, Theorem 2)). Let $\mathcal{B}_1 \subset \mathcal{B}_2 \subset \mathcal{F}$, then

- $\mathbb{L}^2(\mathcal{B}_1) \subseteq \mathbb{L}^2(\mathcal{B}_2) \subseteq \mathbb{L}^2(\mathcal{F})$; $\mathbb{L}^2(\mathcal{B}_1) \cap \mathbb{L}^2(\mathcal{B}_2) = \mathbb{L}^2(\mathcal{B}_1 \cap \mathcal{B}_2)$.

Recall that, since for $B \subset A \in \mathcal{P}_D$ we have that $\sigma_B \subset \sigma_A$, then:

$$\mathbb{L}^{2}\left(\sigma_{B}
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 is a closed Hilbert subspace of $\mathbb{L}^{2}\left(\sigma_{A}
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and all of them are closed subspaces of $\mathbb{L}^2(\sigma_X)$: They are **nested in very a particular way** (more on that later in the talk).

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Controlling the Lebesgue spaces w.r.t. the σ -algebras allow to express spaces of functions of subsets of inputs (analogously to Chastaing, Gamboa, and Prieur (2012)).

Recall the classical result:

Theorem (Malliavin (1995, Chapter 3)). Let X and Y be two random elements. Then:

$$X \perp\!\!\!\perp Y \iff \forall f(X) \in \mathbb{L}^2\left(\sigma_X\right), \ \forall g(Y) \in \mathbb{L}^2\left(\sigma_Y\right), \ \mathit{Corr}\left(f(X), g(Y)\right) = 0,$$

or, in other words, $\mathbb{L}^2_0\left(\sigma_X\right) \perp \mathbb{L}^2_0\left(\sigma_Y\right)$, where \mathbb{L}^2_0 only contains centered random variables.

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Intuition:

Is it possible to control the **dependence structure** between the inputs by controlling the **angles between the subspaces** $\{\mathbb{L}^2(\sigma_A)\}_{A\in\mathcal{P}_D}$?

Dixmier's angle

Definition (Dixmier's angle (Dixmier 1949)). Let M, N be **closed** subspaces of a Hilbert space H. The cosine of Dixmier's angle between M and N is defined as

$$c_0\left(M,N\right):=\sup\left\{|\langle x,y\rangle|:x\in M,\|x\|\leq 1,\quad y\in N,\|y\|\leq 1\right\}.$$

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$$c_0(M, N) := \sup \{ |\langle x, y \rangle| : x \in M, ||x|| \le 1, \quad y \in N, ||y|| \le 1 \}.$$

Dixmier's angle is closely related to the notion of **maximal correlation** in probability theory (Koyak 1987), as a dependence measure between **random elements**.

Definition (Maximal correlation (Gebelein 1941)). Let Z_1, Z_2 be two **random elements**. The maximal correlation between Z_1 and Z_2 is

$$\rho_0(\mathit{Z}_1,\mathit{Z}_2) := c_0\left(\mathbb{L}_0^2\left(\sigma_{\mathit{Z}_1}\right),\mathbb{L}_0^2\left(\sigma_{\mathit{Z}_2}\right)\right)$$

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Remark .

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Friedrich's angle is used in probability theory as a measure of **partial dependence** between two random elements (Bryc 1984, 1996; Dauxois, Nkiet, and Romain 2004).

Definition (Maximal partial correlation). Let Z_1 and Z_2 be two random elements. The maximal partial correlation is between Z_1 and Z_2 is

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Remark .

$$\rho^*\left(Z_1,Z_2\right) = 0 \iff \mathbb{E}\left[\mathbb{E}\left[.\mid Z_1\right]\mid Z_2\right] = \mathbb{E}\left[\mathbb{E}\left[.\mid Z_2\right]\mid Z_1\right]$$

These two angles are related to the closedness of the sum of the two subspaces:

- $c(M, N) < 1 \iff M + N \text{ is closed in } H$;
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One popular complement of a closed subspace M is its orthogonal complement M^{\perp} .

Feshchenko matrix

Let's go back to our set of subspaces $\left\{ \mathbb{L}^{2}\left(\sigma_{A}\right) \right\}_{A\in\mathcal{P}_{D}}.$

How can we "globally" control all the Friedrichs' angles between them?

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Definition (Maximal coalitional precision matrix). Let Δ be the $(2^d \times 2^d)$, symmetric **set-indexed** matrix, defined element-wise, $\forall A, B \in \mathcal{P}_D$ as

$$\Delta_{AB} = egin{cases} 1 & ext{if } A = B; \ -c\left(\mathbb{L}^2\left(\sigma_A
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These matrices resemble closely the ones used by **Feshchenko (2020)** to study the **closedness of an arbitrary sum of closed subspaces** of a Hilbert space.

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⇒ We're going to call them "Feshchenko matrices".

Stochastic dependence

But why is the Feshchenko matrix interesting?

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Proposition. Suppose that Assumption 1 hold. Then,

$$\Delta = \mathit{I}_{2^d} \quad \Longleftrightarrow \; \mathit{X} \; \textit{is mutually independent}.$$

Remark . Recall that we're working with **abstract-valued random elements** (and not necessarily a random vector).

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Our second assumption:

Assumption 2 (Non-degenerate stochastic dependence). The Feshchenko matrix Δ of the inputs is definite-positive.

Note that this is a restriction of the inner product of $\mathbb{L}^2(\sigma_X)$, and thus **an indirect restriction on** the law of X.

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Definition (Direct-sum decomposition (Axler 2015)). Let W be a vector space and let W_1, \ldots, W_n be **proper subspaces** of W.

W is said to admit a direct-sum decomposition if any $w \in W$ can be written uniquely as

$$w = \sum_{i=1}^n w_i$$
 where $w_i \in W_i$ for $i = 1, \dots, n$.

In this case, we write:

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But which subspaces should be involved in the direct-sum decomposition?

Generalized Hoeffding decomposition

Theorem . Under Assumptions 1 and 2, for every $A \in \mathcal{P}_D$, one has that

$$\mathbb{L}^{2}\left(\sigma_{A}
ight)=igoplus_{B\in\mathcal{P}_{A}}V_{B}.$$

where $V_\emptyset=\mathbb{L}^2\left(\sigma_\emptyset
ight)$, and

$$V_B = \left[\begin{array}{c} + \\ C \in \mathcal{P}_B, C \neq B \end{array} \right]^{\perp_B},$$

where \perp_{B} denotes the orthogonal complement in \mathbb{L}^2 (σ_{B}).

Generalized Hoeffding decomposition

Theorem . Under Assumptions 1 and 2, for every $A \in \mathcal{P}_D$, one has that

$$\mathbb{L}^{2}\left(\sigma_{A}\right)=\bigoplus_{B\in\mathcal{P}_{A}}V_{B}.$$

where $V_\emptyset=\mathbb{L}^2\left(\sigma_\emptyset
ight)$, and

$$V_B = \left[\frac{1}{C \in \mathcal{P}_B, C \neq B} V_C \right]^{\perp_B},$$

where \perp_B denotes the orthogonal complement in \mathbb{L}^2 (σ_B).

Main intuition:

"Inductive generalized centering"

Intuition behind the result: One input

One input:

- 1. Let $i \in D$, and fix $\mathbb{L}^2(\sigma_i)$ as the ambient space.
- 2. We have that $V_{\emptyset} := \mathbb{L}^2(\sigma_{\emptyset})$ is a closed subspace of $\mathbb{L}^2(\sigma_i)$ (thus it is complemented).
- 3. Denote $V_i = [V_{\emptyset}]^{\perp_i}$, the orthogonal complement of V_{\emptyset} in \mathbb{L}^2 (σ_i) .
- 4. One has that $\mathbb{L}^2(\sigma_i) = V_\emptyset \oplus V_i$.
- 5. Since V_{\emptyset} only contains constants, $V_{i}=\mathbb{L}_{0}^{2}\left(\sigma_{i}\right)$.

In other words, we just showed that any $f(X_i) \in \mathbb{L}^2(\sigma_i)$ can be written as

$$f(X_i) = \underbrace{\mathbb{E}\left[f(X_i)\right]}_{\in V_\emptyset} + \underbrace{\mathbb{E}\left[f(X_i) - \mathbb{E}\left[f(X_i)\right]\right]}_{\in V_i}.$$

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And note that we can do this for any $i \in D$.

Intuition behind the result: Two inputs

Two inputs:

- 1. Let $i, j \in D$, and fix $\mathbb{L}^2(\sigma_{ij})$ as the ambient space.
- 2. We have that $\mathbb{L}^2(\sigma_i)$ and $\mathbb{L}^2(\sigma_j)$ are closed subspaces of $\mathbb{L}^2(\sigma_{ij})$.
- 3. Assumptions 1 and 2 imply that $\mathbb{L}^2(\sigma_i) + \mathbb{L}^2(\sigma_j)$ is closed in $\mathbb{L}^2(\sigma_{ij})$ (thus it is complemented).
- 4. Notice (previous step) that $\mathbb{L}^2(\sigma_i) + \mathbb{L}^2(\sigma_j) = V_\emptyset + V_i + V_j$.
- 5. Denote $V_{ij} = \left[V_{\emptyset} + V_i + V_j\right]^{\perp_{ij}}$, the orthogonal complement in $\mathbb{L}^2\left(\sigma_{ij}\right)$.
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In essence, we "centered" a bivariate function from its "univariate and constant parts".

And we can continue the same induction up to d inputs.

Orthocanonical decomposition

As a direct consequence of the previous theorem:

Corollary (Orthocanonical decomposition). Under Assumptions 1 and 2, any $G(X) \in \mathbb{L}^2(\sigma_X)$ can be uniquely decomposed as

$$G(X) = \sum_{A \in \mathcal{P}_D} G_A(X_A),$$

where each $G_A(X_A) \in V_A$.

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The subspaces V_A are comprised of **proper representants**, i.e., either 0 or **functions of exactly** X_A (they do not contain functions of fewer inputs).

Projectors

Recall that for any $G(X) \in \mathbb{L}^2(\sigma_X)$, we have

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Denote the operator

$$Q_A: \mathbb{L}^2\left(\sigma_X
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 Q_A is the (canonical) **oblique projection** onto V_A , parallel to $\bigoplus_{B \in \mathcal{P}_D: B \neq A} V_A$.

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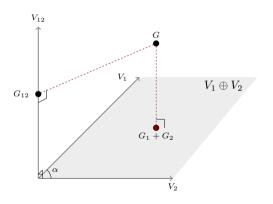
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Illustration : $\mathbb{L}_0^2(\sigma_{12})$

Hence, for any $G(X) \in \mathbb{L}^2(\sigma_X)$, one has that, $\forall A \in \mathcal{P}_D$

$$G_A(X_A) = Q_A(G(X)).$$



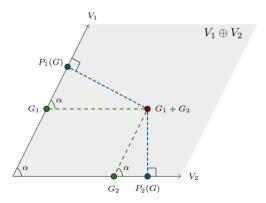
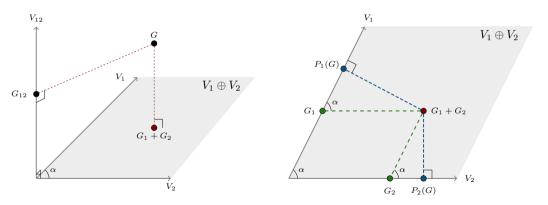


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The oblique projection Q_A usually differ from the oblique projections P_A

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Proposition. Under Assumptions 1 and 2,

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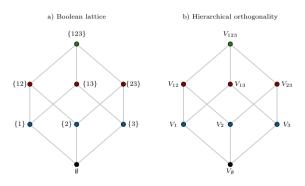
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To illustrate this fact, we need some algebraic combinatorics.

Boolean lattice and hierarchical orthogonality

Our decomposition is **over the power-set** \mathcal{P}_D , which **which is not trivial**.

When endowed with the **binary relation** \subseteq they form an algebraic structure called **a Boolean lattice**.



The subspaces $\{V_A\}_{A\in\mathcal{P}_D}$ are **hierarchically orthogonal** by design: they follow the same algebraic structure, but this time **w.r.t. to** \bot .

Recall that:

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But what about projections onto the subspaces $\left\{\mathbb{L}^2\left(\sigma_A\right)\right\}_{A\in\mathcal{P}_D}$?

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Is it possible to express the projections Q_A using \mathbb{M}_A ?

Generalized Möbius inversion

Yes, because we're working on the power-set $\mathcal{P}_{\mathcal{D}}!$

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Corollary (Möbius inversion on power-sets (Rota 1964)). Let $D = \{1, ..., d\}$. For any two set functions:

$$f: \mathcal{P}_D \to \mathbb{A}, \quad g: \mathcal{P}_D \to \mathbb{A},$$

where \mathbb{A} is an **abelian group**, the following equivalence holds:

$$f(A) = \sum_{B \in \mathcal{P}_A} g(B), \quad \forall A \in \mathcal{P}_D \quad \Longleftrightarrow \quad g(A) = \sum_{B \in \mathcal{P}_A} (-1)^{|A| - |B|} f(B), \quad \forall A \in \mathcal{P}_D.$$

In our case, we have, by definition of the oblique projection onto $\mathbb{L}^2(\sigma_A)$, that

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(This is what we call the "model-centric" approach)

If the inputs are mutually independent, from Hoeffding (1948), we have that:

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Our approach actually generalizes Hoeffding's decomposition!

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Organic variance decomposition: separate pure interaction effects to dependence effects. The dependence structure of X is unwanted, and one wishes to study its effects.

Orthocanonical variance decomposition: the dependence structure of X is inherent in the uncertainty modeling of the studied phenomenon. It amounts to quantify structural and correlative effects.

Organic variance decomposition: Pure interaction

The notion of pure interaction is intrinsically linked with the notion of mutual independence.

Let $\widetilde{X} = (\widetilde{X}_1, \dots, \widetilde{X}_d)^{ op}$ be the random vector such that

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Definition (Pure interaction). For every $A \in \mathcal{P}_D$, define the **pure interaction of** X_A **on** G(X) **as**

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This approach **strongly resembles the "independent Sobol' indices"** proposed by Mara, Tarantola, and Annoni (2015).

(see, also, Lebrun and Dutfoy (2009a, 2009b))

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Canonical variance decomposition

The structural effects represent the variance of each of the $G_A(X_A)$. It amounts to perform a **covariance decomposition** (Hart and Gremaud 2018; Da Veiga et al. 2021).

Definition (Structural effects). For every $A \in \mathcal{P}_D$, define the **structural effects of** X_A **on** G(X) **as**

$$S_A^U = \mathbb{V}(G_A(X_A)).$$

The **correlative effects** represent the part of variance that is due to the correlation between the $G_A(X_A)$.

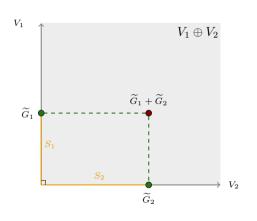
Definition (Correlative effects). For every $A \in \mathcal{P}_D$, define the **correlative effects of** X_A **on** G(X) as

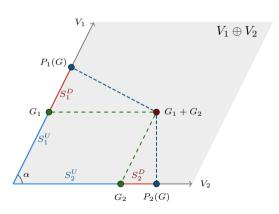
$$S_A^C = Cov\left(G_A(X_A), \sum_{B \in \mathcal{P}_D: B \neq A} G_B(X_B)\right).$$

Variance decomposition: Intuition



Structural and dependence effects





Example: Two Bernoulli inputs

Let
$$E = \{0,1\}^2$$
, and let $X = (X_1, X_2)$, where

$$X_1 \sim \mathcal{B}\left(\mathbf{q}_1\right), \quad \text{and } X_2 \sim \mathcal{B}\left(\mathbf{q}_2\right).$$

The joint law of X can be express using three parameters:

$$p_{00} = 1 - q_1 - q_2 + \rho, \quad p_{01} = q_2 - \rho, \quad p_{10} = q_1 - \rho, \quad p_{11} = \rho$$

where $p_{ij} = \mathbb{P}(\{X_1 = i\} \cap \{X_2 = j\}).$

Any function $G:\{0,1\}^2 \to \mathbb{R}$ can be expressed as the vector $G=(G_{00},G_{01},G_{10},G_{11})^{\top}$.

Each value $G_{ij} = G(i, j)$, can be observed with probability p_{ij} .

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In this case, we can compute everything analytically.

It requires to solving 13 equations with 13 unknowns*.

Feshchenko matrix and the Fréchet bounds

For the **Feshchenko matrix** Δ to be definite positive, one has that:

$$\max\left\{0,q_1q_2-\sqrt{q_1q_2(1-q_1)(1-q_2)}\right\}<\rho<\min\left\{1,q_1q_2-\sqrt{q_1q_2(1-q_1)(1-q_2)}\right\}.$$

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However, the classical Fréchet bounds for ρ for bivariate Bernoulli random variables (Joe 1997, p.210) are equal to

$$\max \left\{ 0, q_1 + q_2 - 1 \right\} \le \rho \le \min \left\{ q1, q2 \right\},\,$$

and are more restrictive than the previous ones.

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ho strictly contained in the Fréchet bounds \implies Assumptions 1 and 2 hold.

Our decomposition hold for virtually any dependence structure between two Bernoullis.

Conclusion

Main take-aways:

- Hoeffding-like decomposition of function with dependent inputs is achievable under fairly reasonable assumptions.
- Mixing probability, functional analysis and combinatorics lead to a linear treatment of multivariate non-linear stochastic problems.
- We can define intuitive model-centric decompositions of quantities of interest.
- We proposed candidates to separate pure interaction and dependence effects.

Perspectives

Main challenge: Estimation.

• We haven't found an off-the-shelf method to estimate the oblique projections...

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A few perspectives:

- Causality and algebraic structures beyond the Boolean lattice.
- Link between Feshchenko matrices and copulas.
- Non \mathbb{R} -valued output.
- Beyond the MSE for surrogate modelling.
- Many methodological questions that seemed unreachable so far, but appear approachable using this framework.

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THANK YOU FOR YOUR ATTENTION!

ANY QUESTIONS?

Annihilating property

Proposition (Annihilating property). For any $A \in \mathcal{P}_D$ and any $B \subset A$

$$P_B\left(Q_A\left(G(X)\right)\right)=P_B\left(G_A(X_A)\right)=0.$$