

Variational analysis for options with stochastic volatility and multiple factors

J. Frédéric Bonnans, Axel Kröner

▶ To cite this version:

J. Frédéric Bonnans, Axel Kröner. Variational analysis for options with stochastic volatility and multiple factors. 2018. hal-01516011v2

HAL Id: hal-01516011 https://inria.hal.science/hal-01516011v2

Preprint submitted on 9 Apr 2018 (v2), last revised 9 Apr 2018 (v3)

HAL is a multi-disciplinary open access archive for the deposit and dissemination of scientific research documents, whether they are published or not. The documents may come from teaching and research institutions in France or abroad, or from public or private research centers. L'archive ouverte pluridisciplinaire **HAL**, est destinée au dépôt et à la diffusion de documents scientifiques de niveau recherche, publiés ou non, émanant des établissements d'enseignement et de recherche français ou étrangers, des laboratoires publics ou privés.

VARIATIONAL ANALYSIS FOR OPTIONS WITH STOCHASTIC VOLATILITY AND MULTIPLE FACTORS*

J. FRÉDÉRIC BONNANS† AND AXEL KRÖNER‡

Abstract. This paper performs a variational analysis for a class of European or American options with stochastic volatility models, including those of Heston and Achdou-Tchou. Taking into account partial correlations and the presence of multiple factors, we obtain the well-posedness of the related partial differential equations, in some weighted Sobolev spaces. This involves a generalization of the commutator analysis introduced by Achdou and Tchou in [3].

Key words. finance, options, partial differential equations, variational formulation, parabolic variational inequalities

AMS subject classifications. 35K20, 35K85, 91G80

1. Introduction. In this paper we consider variational analysis for the partial differential equations associated with the pricing of European or American options. For an introduction to these models, see Fouque et al., [11]. We will set up a general framework of variable volatility models, which is in particular applicable on the following standard models which are well established in mathematical finance. The well-posedness of PDE formulations of variable volatility poblems was studied in [2, 3, 1, 18], and in the recent work [9, 10].

Let the $W_i(t)$ be Brownian motions on a filtered probability space. The variable s denotes a financial asset, and the components of y are factors that influence the volatility:

(i) The Achdou-Tchou model [3], see also Achdou, Franchi, and Tchou [1]:

(1.1)
$$\begin{cases} ds(t) = rs(t)dt + \sigma(y(t))s(t)dW_1(t), \\ dy(t) = \theta(\mu - y(t))dt + \nu dW_2(t), \end{cases}$$

with the interest rate r, the volatility coefficient σ function of the factor y whose dynamics involves a parameter $\nu > 0$, and positive constants θ and μ .

(ii) The Heston model [14]

(1.2)
$$\begin{cases} ds(t) = s(t) \left(r dt + \sqrt{y(t)} dW_1(t) \right), \\ dy(t) = \theta(\mu - y(t)) dt + \nu \sqrt{y(t)} dW_2(t). \end{cases}$$

(iii) The *Double Heston model*, see Christoffersen, Heston and Jacobs [17], and also Gauthier and Possamaï [12]:

(1.3)
$$\begin{cases} ds(t) = s(t) \left(r dt + \sqrt{y_1(t)} dW_1(t) + \sqrt{y_2(t)} dW_2(t) \right), \\ dy_1(t) = \theta_1(\mu_1 - y_1(t)) dt + \nu_1 \sqrt{y_1(t)} dW_3(t), \\ dy_2(t) = \theta_2(\mu_2 - y_2(t)) dt + \nu_2 \sqrt{y_2(t)} dW_4(t). \end{cases}$$

Funding: The first author was suported by the Laboratoire de Finance des Marchés de l'Energie, Paris, France. Both authors were supported by a public grant as part of the Investissement d'avenir project, reference ANR-11-LABX-0056-LMH, LabEx LMH, in a joint call with Gaspard Monge Program for optimization, operations research and their interactions with data sciences.

[†]CMAP, Ecole Polytechnique, CNRS, Université Paris Saclay, and Inria, France (Frederic.Bonnans@inria.fr).

[‡]Department of Mathematics, Humboldt University of Berlin, Berlin, Germany; CMAP, Ecole Polytechnique, CNRS, Université Paris Saclay, and Inria, France (axel.kroener@math.hu-berlin.de)

^{*}Submitted to the editors DATE.

31

32

33

34

53

54

56

57

60

62

63

64

65 66

67 68

70

71

72

In the last two models we have similar interpretations of the coefficients; in the double Heston model, denoting by $\langle \cdot, \cdot \rangle$ the correlation coefficients, we assume that there are correlations only between W_1 and W_3 , and W_2 and W_4 . Consider now the general multiple factor model

35 (1.4)
$$ds = rs(t)dt + \sum_{k=1}^{N} f_k(y_k(t))s^{\beta_k}(t)dW_k(t),$$

$$dy_k = \theta_k(\mu_k - y_k(t))dt + g_k(y_k(t))dW_{N+k}(t), \quad k = 1, \dots, N.$$

Here the y_k are volatility factors, $f_k(y_k)$ represents the volatility coefficient due to y_k , $g_k(y_k)$ is a volatility coefficient in the dynamics of the kth factor with positive constants θ_k and μ_k . Let us denote the correlation between the ith and jth Brownian motions by κ_{ij} : this is a measurable function of (s, y, t) with value in [0, 1] (here $s \in (0, \infty)$ and y_k belongs to either $(0, \infty)$ or \mathbb{R}), see below. We assume that we have nonzero correlations only between the Brownian motions W_k and W_{N+k} , for k=1 to N, i.e.

43 (1.5)
$$\kappa_{ij} = 0 \quad \text{if } i \neq j \text{ and } |j - i| \neq N.$$

Note that, in some of the main results, we will assume for the sake of simplicity that

45 the correlations are constant.

46 We apply the developed analysis to a subclass of stochastic volatility models, obtained

47 by assuming that κ is constant and

48 (1.6)
$$|f_k(y_k)| = |y_k|^{\gamma_k}$$
; $|g_k(y_k)| = \nu_k |y_k|^{1-\gamma_k}$; $\beta_k \in (0,1]$; $\nu_k > 0$; $\gamma_k \in (0,\infty)$.

49 This covers in particular a variant of the Achdou and Tchou model with multiple

factors (VAT), when $\gamma_k = 1$, as well as a generalized multiple factor Heston model

51 (GMH), when $\gamma_k = 1/2$, i.e., for k = 1 to N:

52 (1.7)
$$\begin{array}{ccc} \text{VAT:} & f_k(y_k) = y_k, & g_k(y_k) = \nu_k, \\ \text{GMH:} & f_k(y_k) = \sqrt{y}_k, & g_k(y_k) = \nu_k \sqrt{y}_k. \end{array}$$

For a general class of stochastic volatility models with correlation we refer to Lions and Musiela [16].

The main contribution of this paper is variational analysis for the pricing equation corresponding to the above general class in the sense of the Feynman-Kac theory. This requires in particular to prove continuity and coercivity properties of the corresponding bilinear form in weighted Sobolev spaces H and V, respectively, which have the Gelfand property and allow the application of the Lions and Magenes theory [15] recalled in Appendix A and the regularity theory for parabolic variational inequalities recalled in Appendix B. A special emphasis is given to the continuity analysis of the rate term in the pricing equation. Two approaches are presented, the standard one and an extension of the one based on the commutator of first-order differential operators as in Achdou and Tchou [3], extended to the Heston model setting by Achdou and Pironneau [18]. Our main result is that the commutator analysis gives stronger results for the subclass defined by (1.6), generalizing the particular cases of the VAT and GMH classes, see remarks 6.2 and 6.4. In particular we extend some of the results by [3].

This paper is organized as follows. In section 2 we give the expression of the bilinear form associated with the original PDE, and check the hypotheses of continuity and semi-coercivity of this bilinear form. In section 3 we show how to refine this analysis by taking into account the commutators of the first-order differential operators

associated with the variational formulation. In section 4 we show how to compute the weighting function involved in the bilinear form. In section 5 we develop the results for a general class introduced in the next section. In section 6 we specialize the results to stochastic volatility models. The appendix recalls the main results of the variational theory for parabolic equations, with a discussion on the characterization of the V functional spaces in the case of one dimensional problems.

Notation. We assume that the domain Ω of the PDEs to be considered in the sequel of this paper has the following structure. Let (I, J) be a partition of $\{0, \ldots, N\}$, with $0 \in J$, and

$$\Omega := \prod_{k=0}^{N} \Omega_k; \quad \text{with} \quad \Omega_k := \begin{cases} \mathbb{R} & \text{when } k \in I, \\ (0, \infty) & \text{when } k \in J. \end{cases}$$

Let $L^0(\Omega)$ denote the space of measurable functions over Ω . For a given weighting function $\rho: \Omega \to \mathbb{R}$ of class C^1 , with positive values, we define the weighted space

85 (1.9)
$$L^{2,\rho}(\Omega) := \{ v \in L^0(\Omega); \ \int_{\Omega} v(x)^2 \rho(x) dx < \infty \},$$

86 which is a Hilbert space endowed with the norm

80

87 (1.10)
$$||v||_{\rho} := \left(\int_{\Omega} v(x)^2 \rho(x) dx \right)^{1/2}.$$

88 By $\mathcal{D}(\Omega)$ we denote the space of C^{∞} functions with compact support in Ω . By $H^2_{loc}(\Omega)$ we denote the space of functions over Ω whose product with an element of $\mathcal{D}(\Omega)$ belongs to the Sobolev space $H^2(\Omega)$.

Besides, let Φ be a vector field over Ω (i.e., a mapping $\Omega \to \mathbb{R}^n$). The first-order differential operator associated with Φ is, for $u:\Omega \to \mathbb{R}$ the function over Ω defined by

$$\Phi[u](x) := \sum_{i=0}^{n} \Phi_i(x) \frac{\partial u}{\partial x_i}(x), \quad \text{for all } x \in \Omega.$$

- 2. General setting. Here we give compute the bilinear form associated with the original PDE, in the setting of the general multiple factor model (1.4). Then we will check the hypotheses of continuity and semi-coercivity of this bilinear form.
- 2.1. Variational formulation. We compute the bilinear form of the variational setting, taking into account a general weight function. We will see how to choose the functional spaces for a given ρ , and then how to choose the weight itself.
- 2.1.1. The elliptic operator. In financial models the underlying is solution of stochastic differential equations of the form

103 (2.1)
$$dX(t) = b(t, X(t))dt + \sum_{i=1}^{n_{\sigma}} \sigma_i(t, X(t))dW_i.$$

Here X(t) takes values in Ω , defined in (1.8). That is, X_1 corresponds to the s variable, and X_{k+1} , for k=1 to N, corresponds to y_k . We have that $n_{\sigma}=2N$.

So, b and σ_i , for i = 1 to n_{σ} , are mappings $(0,T) \times \Omega \to \mathbb{R}^n$, and the W_i , for i = 1 to n_{σ} , are standard Brownian processes with correlation $\kappa_{ij} : (0,T) \times \Omega \to \mathbb{R}$

- between W_i and W_j for $i, j \in \{1, ..., n_{\sigma}\}$. The $n_{\sigma} \times n_{\sigma}$ symmetric correlation matrix $\kappa(\cdot, \cdot)$ is nonnegative with unit diagonal:
- 111 (2.2) $\kappa(t,x) \succ 0$; $\kappa_{ii} = 1, i = 1, \dots, n_{\sigma}$, for a.a. $(t,x) \in (0,T) \times \Omega$.
- Here, for symmetric matrices B and C of same size, by " $C \succeq B$ " we mean that C B
- 113 is positive semidefinite. The expression of the second order differential operator A
- corresponding to the dynamics (2.1) is, skipping the time and space arguments, for
- 115 $u: (0,T) \times \Omega \to \mathbb{R}$:

116 (2.3)
$$Au := ru - b \cdot \nabla u - \frac{1}{2} \sum_{i,j=1}^{n_{\sigma}} \kappa_{ij} \sigma_j^{\top} u_{xx} \sigma_i,$$

117 where

118 (2.4)
$$\sigma_j^{\top} u_{xx} \sigma_i := \sum_{k,\ell=1}^{n_{\sigma}} \sigma_{kj} \frac{\partial u^2}{\partial x_k \partial x_{\ell}} \sigma_{\ell i},$$

- r(x,t) represents an interest rate, and u_{xx} is the matrix of second derivatives in space
- of u. The associated backward PDE for a European option is of the form

121 (2.5)
$$\begin{cases} -\dot{u}(t,x) + A(t,x)u(t,x) = f(t,x), & (t,x) \in (0,T) \times \Omega; \\ u(x,T) = u_T(x), & x \in \Omega, \end{cases}$$

- with \dot{u} the notation for the time derivative of u, $u_T(x)$ payoff at final time (horizon)
- 123 T and the r.h.s. f(t,x) represents dividends (often equal to zero).
- In case of an American option we obtain a variational inequality; for details we
- 125 refer to Appendix D.
- 2.1.2. The bilinear form. In the sequel we assume that

127 (2.6)
$$b, \sigma, \kappa \text{ are } C^1 \text{ mappings over } [0, T] \times \Omega.$$

- Multiplying (2.3) by the test function $v \in \mathcal{D}(\Omega)$ and the continuously differentiable
- weight function $\rho \colon \Omega \to \mathbb{R}$ and integrating over the domain we can integrate by parts;
- since $v \in \mathcal{D}(\Omega)$ there will be no contribution from the boundary. We obtain

131 (2.7)
$$-\frac{1}{2} \int_{\Omega} \sigma_{j}^{\top} u_{xx} \sigma_{i} v \kappa_{ij} \rho = \sum_{p=0}^{3} a_{ij}^{p}(u, v),$$

132 with

133 (2.8)
$$a_{ij}^{0}(u,v) := \frac{1}{2} \int_{\Omega} \sum_{k,\ell=1}^{n} \sigma_{kj} \sigma_{\ell i} \frac{\partial u}{\partial x_{k}} \frac{\partial v}{\partial x_{\ell}} \kappa_{ij} \rho = \frac{1}{2} \int_{\Omega} \sigma_{j}[u] \sigma_{i}[v] \kappa_{ij} \rho,$$

134

135 (2.9)
$$a_{ij}^{1}(u,v) := \frac{1}{2} \int_{\Omega} \sum_{k,\ell=1}^{n} \sigma_{kj} \sigma_{\ell i} \frac{\partial u}{\partial x_{k}} \frac{\partial (\kappa_{ij}\rho)}{\partial x_{\ell}} v = \frac{1}{2} \int_{\Omega} \sigma_{j}[u] \sigma_{i}[\kappa_{ij}\rho] \frac{v}{\rho} \rho,$$

137 (2.10)
$$a_{ij}^2(u,v) := \frac{1}{2} \int_{\Omega} \sum_{k,\ell=1}^n \sigma_{kj} \frac{\partial(\sigma_{\ell i})}{\partial x_\ell} \frac{\partial u}{\partial x_k} v \kappa_{ij} \rho = \frac{1}{2} \int_{\Omega} \sigma_j[u] (\operatorname{div} \sigma_i) v \kappa_{ij} \rho,$$

138

139 (2.11)
$$a_{ij}^{3}(u,v) := \frac{1}{2} \int_{\Omega} \sum_{k,\ell=1}^{n} \frac{\partial(\sigma_{kj})}{\partial x_{\ell}} \sigma_{\ell i} \frac{\partial u}{\partial x_{k}} v \kappa_{ij} \rho = \frac{1}{2} \int_{\Omega} \sum_{k=1}^{n} \sigma_{i}[\sigma_{kj}] \frac{\partial u}{\partial x_{k}} v \kappa_{ij} \rho.$$

Also, for the contributions of the first and zero order terms resp. we get

141 (2.12)
$$a^4(u,v) := -\int_{\Omega} b[u]v\rho; \quad a^5(u,v) := \int_{\Omega} ruv\rho.$$

142 Set

143 (2.13)
$$a^p := \sum_{i,j=1}^{n_{\sigma}} a_{ij}^p, \quad p = 0, \dots, 3.$$

144 The bilinear form associated with the above PDE is

145 (2.14)
$$a(u,v) := \sum_{p=0}^{5} a^{p}(u,v).$$

- 146 From the previous discussion we deduce that
- LEMMA 2.1. Let $u \in H^2_{loc}(\Omega)$ and $v \in \mathcal{D}(\Omega)$. Then we have that

148 (2.15)
$$a(u,v) = \int_{\Omega} A(t,x)u(x)v(x)\rho(x)dx.$$

2.1.3. The Gelfand triple. We can view a^0 as the principal term of the bilinear form a(u, v). Let σ denote the $n \times n_{\sigma}$ matrix whose σ_i are the columns. Then

151 (2.16)
$$a^{0}(u,v) = \sum_{i,j=1}^{n_{\sigma}} \int_{\Omega} \sigma_{j}[u] \sigma_{i}[v] \kappa_{ij} \rho = \int_{\Omega} \nabla u^{\top} \sigma \kappa \sigma^{\top} \nabla v \rho.$$

- Since $\kappa \succeq 0$, the above integrand is nonnegative when u = v; therefore, $a^0(u, u) \geq 0$.
- When κ is the identity we have that $a^0(u,u)$ is equal to the seminorm $a^{00}(u,u)$, where

154 (2.17)
$$a^{00}(u,u) := \int_{\Omega} |\sigma^{\top} \nabla u|^2 \rho.$$

- 155 In the presence of correlations it is natural to assume that we have a coercivity of the
- 156 same order. That is, we assume that

157 (2.18) For some
$$\gamma \in (0,1]$$
: $\sigma \kappa \sigma^{\top} \succeq \gamma \sigma \sigma^{\top}$, for all $(t,x) \in (0,T) \times \Omega$.

158 Therefore, we have

159 (2.19)
$$a^{0}(u,u) \ge \gamma a^{00}(u,u).$$

160 Remark 2.2. Condition (2.18) holds in particular if

161 (2.20)
$$\kappa \succeq \gamma I$$
,

but may also hold in other situations, e.g., when $n=1, n_{\sigma}=2, \kappa_{12}=1$, and

 $\sigma_1 = \sigma_2 = 1$. Yet when the σ_i are linearly independent, (2.19) is equivalent to (2.20).

We need to choose a pair (V,H) of Hilbert spaces satisfying the Gelfand conditions for the variational setting of Appendix A, namely V densely and continuously embedded in H, $a(\cdot,\cdot)$ continuous and semi-coercive over V. Additionally, the r.h.s. and final condition of (2.5) should belong to $L^2(0,T;V^*)$ and H resp. (and for the second parabolic estimate, to $L^2(0,T;H)$ and V resp.).

We do as follows: for some measurable function $h:\Omega\to\mathbb{R}_+$ to be specified later we define

171 (2.21)
$$\begin{cases} H := \{ v \in L^0(\Omega); \ hv \in L^{2,\rho}(\Omega) \}, \\ \mathcal{V} := \{ v \in H; \ \sigma_i[v] \in L^{2,\rho}(\Omega), \ i = 1, \dots, n_{\sigma} \}, \\ V := \{ \text{closure of } \mathcal{D}(\Omega) \text{ in } \mathcal{V} \}, \end{cases}$$

endowed with the natural norms.

173 (2.22)
$$||v||_H := ||hv||_{\rho}; \quad ||u||_V^2 := a^{00}(u, u) + ||u||_H^2.$$

- We do not try to characterize the space V since this is problem dependent.
- Obviously, $a^0(u, v)$ is a bilinear continuous form over \mathcal{V} . We next need to choose h so that a(u, v) is a bilinear and semi-coercive continuous form, and $u_T \in H$.
- 2.2. Continuity and semi-coercivity of the bilinear form over \mathcal{V} . We will see that the analysis of a^0 to a^2 is relatively easy. It is less obvious to analyze the term

180 (2.23)
$$a^{34}(u,v) := a^3(u,v) + a^4(u,v).$$

Let $\bar{q}_{ij}(t,x) \in \mathbb{R}^n$ be the vector with kth component equal to

182 (2.24)
$$\bar{q}_{ijk} := \kappa_{ij} \sigma_i [\sigma_{kj}].$$

183 Set

184 (2.25)
$$\hat{q} := \sum_{i,j=1}^{n_{\sigma}} \bar{q}_{ij}, \quad q := \hat{q} - b.$$

185 Then by (2.11)-(2.12), we have that

186 (2.26)
$$a^{34}(u,v) = \int_{\Omega} q[u]v\rho.$$

- We next need to assume that it is possible to choose η_k in $L^0((0,T)\times\Omega)$, for k=1
- 188 to n_{σ} , such that

189 (2.27)
$$q = \sum_{k=1}^{n_{\sigma}} \eta_k \sigma_k.$$

- Often the $n \times n_{\sigma}$ matrix $\sigma(t,x)$ has a.e. rank n. Then the above decomposition is
- 191 possible. However, the choice for η is not necessarily unique. We will see in examples
- 192 how to do it. Consider the following hypotheses:

193 (2.28)
$$h_{\sigma} \leq c_{\sigma} h$$
, where $h_{\sigma} := \sum_{i,j=1}^{n_{\sigma}} |\sigma_i[\kappa_{ij}\rho]/\rho + \kappa_{ij} \operatorname{div} \sigma_i|$, a.e., for some $c_{\sigma} > 0$,

- 194 (2.29) $h_r \le c_r h$, where $h_r := |r|^{1/2}$, a.e., for some $c_r > 0$,
- 195 (2.30) $h_{\eta} \leq c_{\eta} h$, where $h_{\eta} := |\eta|$, a.e., for some $c_{\eta} > 0$.

REMARK 2.3. Let us set for any differentiable vector field $Z: \Omega \to \mathbb{R}^n$

198 (2.31)
$$G_{\rho}(Z) := \operatorname{div} Z + \frac{Z[\rho]}{\rho}.$$

199 Since $\kappa_{ii} = 1$, (2.28) implies that

200 (2.32)
$$|G_{\rho}(\sigma_i)| \le c_{\sigma}h, \ i = 1; \dots, n_{\sigma}.$$

201 Remark 2.4. Since

213

202 (2.33)
$$\sigma_i[\kappa_{ij}\rho] = \sigma_i[\kappa_{ij}]\rho + \sigma_i[\rho]\kappa_{ij},$$

203 and $|\kappa_{ij}| \leq 1$ a.e., a sufficient condition for (2.28) is that there exist a positive con-204 stants c'_{σ} such that

205 (2.34)
$$h'_{\sigma} \leq c'_{\sigma} h; \quad h'_{\sigma} := \sum_{i,j=1}^{n_{\sigma}} |\sigma_{i}[\kappa_{ij}]| + \sum_{i=1}^{n_{\sigma}} (|\operatorname{div} \sigma_{i}| + |\sigma_{i}[\rho]/\rho|).$$

We will see in section 4 how to choose the weight ρ so that $|\sigma_i[\rho]/\rho|$ can be easily estimated as a function of σ .

Lemma 2.5. Let (2.28)-(2.30) hold. Then the bilinear form a(u,v) is both (i) continuous over V, and (ii) semi-coercive, in the sense of (A.5).

210 *Proof.* (i) We have that $a^1 + a^2$ is continuous, since by (2.9)-(2.10), (2.28) and 211 the Cauchy-Schwarz inequality:

$$|a^{1}(u,v) + a^{2}(u,v)| \leq \sum_{\substack{i,j=1\\n_{\sigma}}}^{n_{\sigma}} |a_{ij}^{1}(u,v) + a_{ij}^{2}(u,v)|$$

$$\leq \sum_{j=1}^{n_{\sigma}} \|\sigma_{j}[u]\|_{\rho} \sum_{\substack{i=1\\n_{\sigma}}}^{n_{\sigma}} \|(\sigma_{i}[\kappa_{ij}\rho]/\rho + \kappa_{ij}\operatorname{div}\sigma_{i})v\|_{\rho}$$

$$\leq c_{\sigma}n_{\sigma}\|v\|_{H} \sum_{j=1}^{n_{\sigma}} \|\sigma_{j}[u]\|_{\rho}.$$

214 (ii) Also, a^{34} is continuous, since by (2.27) and (2.30):

$$|a^{34}(u,v)| \le \sum_{k=1}^{n_{\sigma}} \|\sigma_k[u]\|_{\rho} \|\eta_k v\|_{\rho} \le c_{\eta} \|v\|_{H} \sum_{k=1}^{n_{\sigma}} \|\sigma_k[u]\|_{\rho}.$$

216 Set $c := c_{\sigma} n_{\sigma} + c_{\eta}^2$. By (2.35)–(2.36), we have that

$$\begin{cases}
|a^{5}(u,v)| \leq ||r|^{1/2}u||_{2,\rho}||r|^{1/2}v||_{2,\rho} \leq c_{r}^{2}||u||_{H}||v||_{H}, \\
|a^{1}(u,v) + a^{2}(u,v) + a^{34}(u,v)| \leq ca^{00}(u^{1/2}||v||_{H}.
\end{cases}$$

218 Since a^0 is obviously continuous, the continuity of a(u, v) follows.

219 (iii) Semi-coercivity. Using (2.37) and Young's inequality, we get that

$$(2.38) a(u,u) \geq a_0(u,u) - \left| a^1(u,u) + a^2(u,u) + a^{34}(u,u) \right| - \left| a^5(u,u) \right| \\ \geq \gamma a^{00}(u) - ca^{00}(u)^{1/2} ||u||_H - c_r ||u||_H^2 \\ \geq \frac{1}{2} \gamma a^{00}(u) - \left(\frac{1}{2} \frac{c^2}{\gamma} + c_r \right) ||u||_H^2,$$

which means that a is semi-coercive.

The above consideration allow to derive well-posedness results for parabolic equations and parabolic variational inequalities.

- THEOREM 2.6. (i) Let (V, H) be given by (2.21), with h satisfying (2.28)-(2.30), $(f, u_T) \in L^2(0, T; V^*) \times H$. Then equation (2.5) has a unique solution u in $L^2(0, T; V)$ with $\dot{u} \in L^2(0, T; V^*)$, and the mapping $(f, u_T) \mapsto u$ is nondecreasing. (ii) If in addition the semi-symmetry condition (A.8) holds, then u in $L^{\infty}(0, T; V)$ and $\dot{u} \in L^2(0, T; H)$.
- 229 *Proof.* This is a direct consequence of Propositions A.1, A.2 and C.1. \square
- 230 We next consider the case of parabolic variational inequalities associated with the set

231 (2.39)
$$K := \{ \psi \in V : \psi(x) \ge \Psi(x) \text{ a.e. in } \Omega \},$$

- where $\Psi \in V$. The strong and weak formulations of the parabolic variational inequality are defined in (B.2) and (B.5) resp. The abstract notion of monotonicity is discussed in appendix B. We denote by \mathcal{K} the closure of K in V.
- THEOREM 2.7. (i) Let the assumptions of theorem 2.6 hold, with $u_T \in \mathcal{K}$. Then the weak formulation (B.5) has a unique solution u in $L^2(0,T;K) \cap C(0,T;H)$, and the mapping $(f,u_T) \mapsto u$ is nondecreasing.
- 238 (ii) Let in addition the semi-symmetry condition (A.8) be satisfied. Then u is the 239 unique solution of the strong formulation (B.2), belongs to $L^{\infty}(0,T;V)$, and \dot{u} belongs to $L^{2}(0,T;H)$.

- 241 *Proof.* This follows from Propositions B.1 and C.2.
- 3. Variational analysis using the commutator analysis. In the following a commutator for first order differential operators is introduced, and calculus rules are derived.
- 3.1. Commutators. Let $u: \Omega \to \mathbb{R}$ be of class C^2 . Let Φ and Ψ be two vector fields over Ω , both of class C^1 . Recalling (1.11), we may define the *commutator* of the first-order differential operators associated with Φ and Ψ as

248 (3.1)
$$[\Phi, \Psi][u] := \Phi[\Psi[u]] - \Psi[\Phi[u]].$$

249 Note that

$$\Phi[\Psi[u]] = \sum_{i=1}^{n} \Phi_i \frac{\partial (\Psi u)}{\partial x_i} = \sum_{i=1}^{n} \Phi_i \left(\sum_{k=1}^{n} \frac{\partial \Psi_k}{\partial x_i} \frac{\partial u}{\partial x_k} + \Psi_k \frac{\partial^2 u}{\partial x_k \partial x_i} \right).$$

251 So, the expression of the commutator is

$$[\Phi, \Psi][u] = \sum_{i=1}^{n} \left(\Phi_{i} \sum_{k=1}^{n} \frac{\partial \Psi_{k}}{\partial x_{i}} \frac{\partial u}{\partial x_{k}} - \Psi_{i} \sum_{k=1}^{n} \frac{\partial \Phi_{k}}{\partial x_{i}} \frac{\partial u}{\partial x_{k}} \right)$$

$$= \sum_{k=1}^{n} \left(\sum_{i=1}^{n} \Phi_{i} \frac{\partial \Psi_{k}}{\partial x_{i}} - \Psi_{i} \frac{\partial \Phi_{k}}{\partial x_{i}} \right) \frac{\partial u}{\partial x_{k}}.$$

It is another first-order differential operator associated with a vector field (which happens to be the Lie bracket of Φ and Ψ , see e.g.[4]).

3.2. Adjoint. Remembering that H was defined in (2.21), given two vector fields Φ and Ψ over Ω , we define the spaces

257 (3.4)
$$V(\Phi, \Psi) := \{ v \in H; \ \Phi[v], \Psi[v] \in H \},$$

258 (3.5)
$$V(\Phi, \Psi) := \{\text{closure of } \mathcal{D}(\Omega) \text{ in } \mathcal{V}(\Phi, \Psi)\}.$$

- We define the adjoint Φ^{\top} of Φ (view as an operator over say $C^{\infty}(\Omega,\mathbb{R})$, the latter
- being endowed with the scalar product of $L^{2,\rho}(\Omega)$, by

262 (3.6)
$$\langle \Phi^{\top}[u], v \rangle_{\rho} = \langle u, \Phi[v] \rangle_{\rho} \text{ for all } u, v \in \mathcal{D}(\Omega),$$

where $\langle \cdot, \cdot \rangle_{\rho}$ denotes the scalar product in $L^{2,\rho}(\Omega)$. Thus, there holds the identity

264 (3.7)
$$\int_{\Omega} \Phi^{\top}[u](x)v(x)\rho(x)dx = \int_{\Omega} u(x)\Phi[v](x)\rho(x)dx \text{ for all } u,v \in \mathcal{D}(\Omega).$$

265 Furthermore,

$$\int_{\Omega} u \sum_{i=1}^{n} \Phi_{i} \frac{\partial v}{\partial x_{i}} \rho dx = -\sum_{i=1}^{n} \int_{\Omega} v \frac{\partial}{\partial x_{i}} (u \rho \Phi_{i}) dx$$

$$= -\sum_{i=1}^{n} \int_{\Omega} v \left(\frac{\partial}{\partial x_{i}} (u \Phi_{i}) + \frac{u}{\rho} \Phi_{i} \frac{\partial \rho}{\partial x_{i}} \right) \rho dx.$$

267 Hence,

268 (3.9)
$$\Phi^{\top}[u] = -\sum_{i=1}^{n} \frac{\partial}{\partial x_{i}} (u\Phi_{i}) - u\Phi_{i} \frac{\partial \rho}{\partial x_{i}} / \rho = -u \operatorname{div} \Phi - \Phi[u] - u\Phi[\rho] / \rho.$$

Remembering the definition of $G_{\rho}(\Phi)$ in (2.31), we obtain that

270 (3.10)
$$\Phi[u] + \Phi^{\top}[u] + G_o(\Phi)u = 0.$$

3.3. Continuity of the bilinear form associated with the commutator.

272 Setting, for v and w in $V(\Phi, \Psi)$:

$$\Delta(u,v) := \int_{\Omega} [\Phi, \Psi][u](x)v(x)\rho(x)\mathrm{d}x,$$

274 we have

(3.12)

$$\Delta(u,v) = \int_{\Omega} (\Phi[\Psi[u]]v - \Psi[\Phi[u]]v)\rho dx = \int_{\Omega} \Psi[u]\Phi^{\top}[v] - \Phi[u]\Psi^{\top}[v])\rho dx$$

$$= \int_{\Omega} (\Phi[u]\Psi[v] - \Psi[u]\Phi[v])\rho dx + \int_{\Omega} (\Phi[u]G_{\rho}(\Psi)v - \Psi[u]G_{\rho}(\Phi)v)\rho dx.$$

LEMMA 3.1. For $\Delta(\cdot, \cdot)$ to be a continuous bilinear form on $V(\Phi, \Psi)$, it suffices that, for some $c_{\Delta} > 0$:

278 (3.13)
$$|G_{\rho}(\Phi)| + |G_{\rho}(\Psi)| \le c_{\Delta}h$$
 a.e.,

279 and we have then:

280 (3.14)
$$|\Delta(u,v)| \le \|\Psi[u]\|_{\rho} \left(\|\Phi[v]\|_{\rho} + c_{\Delta} \|v\|_{H} \right) + \|\Phi[u]\|_{\rho} \left(\|\Psi[v]\|_{\rho} + c_{\Delta} \|v\|_{H} \right).$$

281 *Proof.* Apply the Cauchy Schwarz inequality to (3.12), and use (3.13) combined 282 with the definition of the space H.

We apply the previous results with $\Phi := \sigma_i$, $\Psi := \sigma_j$. Set for v, w in V:

284 (3.15)
$$\Delta_{ij}(u,v) := \int_{\Omega} [\sigma_i, \sigma_j][u](x)v(x)\rho(x)dx, \quad i, j = 1, \dots, n_{\sigma}.$$

- We recall that V was defined in (2.21).
- COROLLARY 3.2. Let (2.28) hold. Then the $\Delta_{ij}(u,v)$, $i,j=1,\ldots,n_{\sigma}$, are continuous bilinear forms over V.

- 288 Proof. Use remark 2.3 and conclude with lemma 3.1.
- 3.4. Redefining the space H. In section 2.2 we have obtained the continuity and semi-coercivity of a by decomposing q, defined in (2.26), as a linear combination (2.27) of the σ_i . We now take advantage of the previous computation of commutators and assume that, more generally, instead of (2.27), we can decompose q in the form

293 (3.16)
$$q = \sum_{k=1}^{n_{\sigma}} \eta_k'' \sigma_k + \sum_{1 \le i < j \le n_{\sigma}} \eta_{ij}' [\sigma_i, \sigma_j] \quad \text{a.e.}$$

- We assume that η' and η'' are measurable functions over $[0, T] \times \Omega$, that η' is weakly differentiable, and that for some $c'_{\eta} > 0$:
- 296 (3.17) $h'_{\eta} \leq c'_{\eta} h$, where $h'_{\eta} := |\eta''| + \sum_{i,j=1}^{N} |\sigma_{i}[\eta'_{ij}]|$ a.e., $\eta' \in L^{\infty}(\Omega)$.
- Lemma 3.3. Let (2.28), (2.29), and (3.17) hold. Then the bilinear form a(u,v) defined in (2.14) is both (i) continuous and (ii) semi-coercive over V.
- 299 *Proof.* (i) We only have to analyze the contribution of a^{34} (defined in (2.23)), 300 since the other contributions to $a(\cdot,\cdot)$ do not change. For the terms in the first sum 301 in (3.16) we have, as was done in (2.36):

302 (3.18)
$$\left| \int_{\Omega} \sigma_{k}[u] \eta_{k}'' v \rho \right| \leq \|\sigma_{k}[u]\|_{\rho} \|\sigma_{k}[u] \eta_{k}'' v\|_{\rho} \leq \|\sigma_{k}[u]\|_{\rho} \|v\|_{H}.$$

303 (ii) Setting $w := \eta'_{ij}v$ and taking here $(\Phi, \Psi) = (\sigma_i, \sigma_j)$, we get that

304 (3.19)
$$\int_{\Omega} \eta'_{ij}[\sigma_i, \sigma_j)[u] v \rho = \Delta(u, w),$$

where $\Delta(\cdot, \cdot)$ was defined in (3.11). Combining with lemma 3.1, we obtain

$$|\Delta_{ij}(u,v)| \leq \|\sigma_{j}[u]\|_{\rho} \left(\|\sigma_{i}[w]\|_{\rho} + c_{\sigma}\|\eta'_{ij}\|_{\infty} \|v\|_{H}\right) + \|\sigma_{i}[u]\|_{\rho} \left(\|\sigma_{j}[w]\|_{\rho} + c_{\sigma}\|\eta'_{ij}\|_{\infty} \|v\|_{H}\right).$$

307 Since

308 (3.21)
$$\sigma_i[\eta'_{ij}v] = \eta'_{ij}\sigma_i[v] + \sigma_i[\eta'_{ij}]v,$$

309 by (3.17):

310 (3.22)
$$\|\sigma_i[w]\|_{\rho} \le \|\eta'_{ij}\|_{\infty} \|\sigma_i[v]\|_{\rho} + \|\sigma_i[\eta'_{ij}]v\|_{\rho} \le \|\eta'_{ij}\|_{\infty} \|\sigma_i[v]\|_{\rho} + c_{\eta}\|v\|_{H}.$$

- 311 Combining these inequalities, point (i) follows.
- 312 (ii) Use u = v in (3.21) and (3.12). We find after cancellation in (3.12) that

$$\Delta_{ij}(u,\eta'_{ij}u) = \int_{\Omega} u(\sigma_i[u]\sigma_j[\eta'_{ij}] - \sigma_j[u]\sigma_i(\eta'_{ij}))\rho + \int_{\Omega} (\sigma_i[u]G_{\rho}(\sigma_j) - \sigma_j[u]G_{\rho}(\sigma_i))\eta'_{ij}u\rho.$$

314 By (3.17), an upper bound for the absolute value of the first integral is

315 (3.24)
$$\left(\|\sigma_i[u]\|_{\rho} + \|\sigma_j[u]\|_{\rho} \right) \|hu\|_{\rho} \le 2 \|u\|_{\mathcal{V}} \|u\|_{H}.$$

- 316 With (2.28), we get an upper bound for the absolute value of the second integral in
- 317 the same way, so, for any $\varepsilon > 0$:

$$|\Delta_{ij}(u, \eta'_{ij}u)| \le 4 \|u\|_{\mathcal{V}} \|u\|_{H}.$$

We finally have that for some c > 0

$$a(u,u) \geq a_{0}(u,u) - c \|u\|_{\mathcal{V}} \|u\|_{H},$$

$$\geq a_{0}(u,u) - \frac{1}{2} \|u\|_{\mathcal{V}}^{2} - \frac{1}{2}c^{2} \|u\|_{H}^{2},$$

$$= \frac{1}{2} \|u\|_{\mathcal{V}}^{2} - \frac{1}{2}(c^{2} + 1) \|u\|_{H}^{2}.$$

- 321 The conclusion follows.
- Remark 3.4. The statements analogous to theorems 2.6 and 2.7 hold, assuming now that h satisfies (2.28), (2.29), and (3.17) (instead of (2.28)-(2.30)).
- 324 **4. The weight** ρ **.** Classes of weighting functions characterized by their growth 325 are introduced. A major result is the independence of the growth order of the function h on the choice of the weighting function ρ in the class under consideration.
- 4.1. Classes of functions with given growth. In financial models we usually have nonnegative variables and the related functions have polynomial growth. Yet, after a logarithmic transformation, we get real variables whose related functions have exponential growth. This motivates the following definitions.
- We remind that (I, J) is a partition of $\{0, ..., N\}$, with $0 \in J$ and that Ω was defined in (1.8).
- Definition 4.1. Let γ' and γ'' belong to \mathbb{R}^{N+1}_+ , with index from 0 to N. Let $\mathcal{G}(\gamma', \gamma'')$ be the class of functions $\varphi: \Omega \to \mathbb{R}$ such that for some c > 0:

335 (4.1)
$$|\varphi(x)| \le c \left(\prod_{k \in I} (e^{\gamma'_k x_k} + e^{-\gamma''_k x_k}) \right) \left(\prod_{k \in J} (x_k^{\gamma'_k} + x_k^{-\gamma''_k}) \right).$$

- 336 We define \mathcal{G} as the union of $\mathcal{G}(\gamma', \gamma'')$ for all nonnegative (γ', γ'') . We call γ'_k and γ''_k 337 the growth order of φ , w.r.t. x_k , at $-\infty$ and $+\infty$ (resp. at zero and $+\infty$).
- Observe that the class \mathcal{G} is stable by the operations of sum and product, and that if f, g belong to that class, so does h = fg, h having growth orders equal to the sum of the growth orders of f and g. For $a \in \mathbb{R}$, we define

341 (4.2)
$$a^+ := \max(0, a); \quad a^- := \max(0, -a); \quad N(a) := (a^2 + 1)^{1/2}.$$

342 as well as

$$\beta := \rho_I \rho_J,$$

344 where

345 (4.4)
$$\rho_I(x) := \prod_{k \in I} e^{-\alpha'_k N(x_k^+) - \alpha''_k N(x_k^-)},$$

346 (4.5)
$$\rho_J(x) := \prod_{k \in J} \frac{x_k^{\alpha'_k}}{1 + x_k^{\alpha'_k + \alpha''_k}},$$

348 for some nonnegative constants α'_k , α''_k , to be specified later.

LEMMA 4.2. Let $\varphi \in \mathcal{G}(\gamma', \gamma'')$. Then $\varphi \in L^{1,\rho}(\Omega)$ whenever ρ is as above, with α satisfying, for some positive ε' and ε'' , for all k = 0 to N:

$$\begin{cases} \alpha'_k = \varepsilon' + \gamma'_k, & \alpha''_k = \varepsilon'' + \gamma''_k, & k \in I, \\ \alpha'_k = (\varepsilon' + \gamma''_k - 1)_+, & \alpha''_k = 1 + \varepsilon'' + \gamma'_k, & k \in J. \end{cases}$$

352 In addition we can choose for k = 0 (if element of J):

353 (4.7)
$$\begin{cases} \alpha'_0 := (\varepsilon' + \gamma''_0 - 1)_+; \ \alpha''_0 := 0 & \text{if } \varphi(s, y) = 0 \text{ when } s \text{ is far from } 0, \\ \alpha'_0 := 0, \ \alpha''_0 := 1 + \varepsilon'' + \gamma'_0, & \text{if } \varphi(s, y) = 0 \text{ when } s \text{ is close to } 0. \end{cases}$$

Proof. It is enough to prove (4.6), the proof of (4.7) is similar. We know that φ satisfy (4.1) for some c > 0 and γ . We need to check the finiteness of

356 (4.8)
$$\int_{\Omega} \left(\prod_{k \in I} (e^{\gamma'_k y_k} + e^{-\gamma''_k y_k}) \right) \left(\prod_{k \in J} (y_k^{\gamma'_k} + y_k^{-\gamma''_k}) \right) \rho(s, y) d(s, y).$$

357 But the above integral is equal to the product $p_I p_J$ with

358 (4.9)
$$p_I := \prod_{k \in I} \int_{\mathbb{R}} (e^{\gamma'_k x_k} + e^{-\gamma''_k x_k}) e^{-\alpha'_k N(x_k^+) - \alpha''_k N(x_k^-)} dx_k,$$

359 (4.10)
$$p_J := \prod_{k \in J} \int_{\mathbb{R}_+} \frac{x_k^{\alpha'_k + \gamma'_k} + x_k^{\alpha'_k - \gamma''_k}}{1 + x_k^{\alpha'_k + \alpha''_k}} dx_k.$$

Using (4.6) we deduce that p_I is finite since for instance

$$\int_{\mathbb{R}_{+}} (e^{\gamma'_{k}x_{k}} + e^{-\gamma''_{k}x_{k}})e^{-\alpha'_{k}N(x_{k}^{+}) - \alpha''_{k}N(x_{k}^{-})} dx_{k}
\leq 2 \int_{\mathbb{R}_{+}} e^{\gamma'_{k}x_{k}}e^{-(1+\gamma'_{k})x_{k}} dx_{k} = 2 \int_{\mathbb{R}_{+}} e^{-x_{k}} dx_{k} = 2,$$

363 and p_J is finite since

364 (4.12)
$$p_J = \prod_{k \in J} \int_{\mathbb{R}_+} \frac{x_k^{\varepsilon' + \gamma_k' + \gamma_k''} + x_k^{\varepsilon' - 1}}{1 + x_L^{\varepsilon' + \varepsilon'' + \gamma_k' + \gamma_k''}} \mathrm{d}x_k < \infty.$$

 \Box The conclusion follows. \Box

366 **4.2.** On the growth order of h. Set for all k

$$\alpha_k := \alpha_k' + \alpha_k''.$$

- Remember that we take ρ in the form (4.3)-(4.4).
- 369 Lemma 4.3. We have that:
- 370 (i) We have that

371 (4.14)
$$\left\| \frac{\rho_{x_k}}{\rho} \right\|_{\infty} \leq \alpha_k, \quad k \in I; \quad \left\| \frac{x}{\rho} \rho_{x_k} \right\|_{\infty} \leq \alpha_k, \quad k \in J.$$

- (ii) Let h satisfying either (2.28)-(2.30) or (2.28)-(2.29), and (3.17). Then the growth order of h does not depend on the choice of the weighting function ρ .
- *Proof.* (i) For $k \in I$ this is an easy consequence of the fact that $N(\cdot)$ is non expansive. For $k \in J$, we have that

376 (4.15)
$$\frac{x}{\rho} \rho_{x_k} = \frac{x}{\rho} \frac{\alpha_k' x^{\alpha_k' - 1} (1 + x^{\alpha_k}) - x^{\alpha_k'} \alpha_k x^{\alpha_k - 1}}{(1 + x^{\alpha_k})^2} = \frac{\alpha_k' - \alpha_k'' x^{\alpha_k}}{1 + x^{\alpha_k}}.$$

- 377 We easily conclude, discussing the sign of the numerator.
- 378 (ii) The dependence of h w.r.t. ρ is only through the last term in (2.28), namely,
- 379 $\sum_{i} |\sigma_{i}[\rho]/\rho$. By (i) we have that

380 (4.16)
$$\left| \frac{\sigma_i^k[\rho]}{\rho} \right| \le \left\| \frac{\rho_{x_k}}{\rho} \right\|_{\infty} |\sigma_i^k| \le \alpha_k |\sigma_i^k|, \quad k \in I,$$

381

382 (4.17)
$$\left| \frac{\sigma_i^k[\rho]}{\rho} \right| \le \left\| \frac{x_k \rho_{x_k}}{\rho} \right\|_{\infty} \left| \frac{\sigma_i^k}{x_k} \right| \le \alpha_k \left| \frac{\sigma_i^k}{x_k} \right|, \quad k \in J.$$

- In both cases, the choice of α has no influence on the growth order of h.
- 384 **4.3. European option.** In the case of a European option with payoff $u_T(x)$, we need to check that $u_T \in H$, that is, ρ must satisfy

386 (4.18)
$$\int_{\Omega} |u_T(x)|^2 h(x)^2 \rho(x) dx < \infty.$$

- In the framework of the semi-symmetry hypothesis (A.8), we need to check that $u_T \in V$, which gives the additional condition
- 389 (4.19) $\sum_{i=1}^{n_{\sigma}} \int_{\Omega} |\sigma_i[u_T](x)|^2 \rho(x) \mathrm{d}x < \infty.$
- 390 In practice the payoff depends only on s and this allows to simplify the analysis.
- 5. Applications using the commutator analysis. The commutator analysis is applied to the general multiple factor model and estimates for the function h characterizing the space H (defined in (2.21)) are derived. The estimates are compared to the case when the commutator analysis is not applied. The resulting improvement will be established in the next section.

- 5.1. Commutator and continuity analysis. We analyze the general multiple factor model (1.4), which belongs to the class of models (2.1) with $\Omega \subset \mathbb{R}^{1+N}$, $n_{\sigma} = 2N$, and for i = 1 to N:
- 399 (5.1) $\sigma_i[v] = f_i(y_i)s^{\beta_i}v_s; \quad \sigma_{N+i}[v] = g_i(y_i)v_i,$
- with f_i and g_i of class C^1 over Ω . We need to compute the commutators of the first-
- order differential operators associated with the σ_i . The correlations will be denoted
- 402 by
- 403 (5.2) $\hat{\kappa}_k := \kappa_{k,N+k}, \quad k = 1, \dots, N.$
- REMARK 5.1. We use many times the following rule. For $\Omega \subset \mathbb{R}^n$, where n =
- 405 1+N, $u \in H^1(\Omega)$, $a,b \in L^0$, and vector fields $Z[u] := au_{x_1}$ and $Z'[u] := bu_{x_2}$, we
- 406 have $Z[Z'[u]] = a(bu_{x_2})_{x_1} = ab_{x_1}u_{x_2} + abu_{x_1x_2}$, so that
- $[Z, Z'][u] = ab_{x_1}u_{x_2} ba_{x_2}u_{x_1}.$
- 408 We obtain that
- 409 (5.4) $[\sigma_i, \sigma_\ell][u] = (\beta_\ell \beta_i) f_i(y_i) f_\ell(y_\ell) s^{\beta_i + \beta_\ell 1} u_s, \quad 1 \le i < \ell \le N,$
- 410
- 411 (5.5) $[\sigma_i, \sigma_{N+i}][u] = -s^{\beta_i} f_i'(y_i) g_i(y_i) u_s, \quad i = 1, \dots, N,$
- 412 and
- 413 (5.6) $[\sigma_i, \sigma_{N+\ell}][u] = [\sigma_{N+i}, \sigma_{N+\ell}][u] = 0, \quad i \neq \ell.$
- 414 Also,
- div $\sigma_i + \frac{\sigma_i[\rho]}{\rho} = f_i(y_i)s^{\beta_i 1}(\beta_i + s\frac{\rho_s}{\rho}),$ div $\sigma_{N+i} + \frac{\sigma_{N+i}[\rho]}{\rho} = g'_i(y_i) + g_i(y_i)\frac{\rho_i}{\rho}.$
- 5.1.1. Computation of q. Remember the definitions of \bar{q} , \hat{q} and q in (2.24) and (2.25), where δ_{ij} denote the Kronecker operator. We obtain that, for $1 \leq i, j, k \leq N$:
- $\begin{cases}
 \bar{q}_{ij0} &= \delta_{ij}\beta_j f_i^2(y_i) s^{2\beta_i 1}; \quad \bar{q}_{iik} = 0; \\
 \bar{q}_{i,N+j} &= 0; \\
 \bar{q}_{N+i,j,0} &= \delta_{ij}\hat{\kappa}_i f'(y_i) g_i(y_i) s^{\beta_i}; \quad \bar{q}_{N+i,j,k} = 0; \\
 \bar{q}_{N+i,N+j,k} &= \delta_{ijk} g_i(y_i) g_i'(y_i).
 \end{cases}$
- That means, we have for $\hat{q} = \sum_{i,j=1}^{2N} \bar{q}_{ij}$ and $q = \hat{q} b$ that
- $\hat{q}_{0} = \sum_{i=1}^{N} \left(\beta_{i} f_{i}^{2}(y_{i}) s^{2\beta_{i}-1} + \hat{\kappa}_{i} f'(y_{i}) g_{i}(y_{i}) s^{\beta_{i}} \right); \quad q_{0} = \hat{q}_{0} rs,$ $q_{k} = g_{k}(y_{k}) g'_{k}(y_{k}); \quad q_{k} = \hat{q}_{k} - \theta_{k}(\mu_{k} - y_{k}),$ k = 1, ..., N.

5.2. Computation of η' and η'' . The coefficients η' , η'' are solution of (3.16).

422 We can write $\eta = \hat{\eta} + \tilde{\eta}$, where

423 (5.10)
$$\hat{q} = \sum_{i=1}^{n_{\sigma}} \hat{\eta}_{i}^{"}\sigma_{i} + \sum_{1 \leq i,j \leq n_{\sigma}} \hat{\eta}_{ij}^{'}[\sigma_{i},\sigma_{j}], \quad \hat{\eta}_{ij}^{'} = 0 \text{ if } i = j.$$

$$-b = \sum_{i=1}^{n_{\sigma}} \tilde{\eta}_{i}^{"}\sigma_{i} + \sum_{1 \leq i,j \leq n_{\sigma}} \tilde{\eta}_{ij}^{'}[\sigma_{i},\sigma_{j}], \quad \tilde{\eta}_{ij}^{'} = 0 \text{ if } i = j.$$

424 For k = 1 to N, this reduces to

425 (5.11)
$$\begin{cases} \hat{\eta}_{N+k}'' g_k(y_k) = g_k'(y_k) g_k(y_k); \\ \tilde{\eta}_{N+k}'' g_k(y_k) = -\theta_k(\mu_k - y_k). \end{cases}$$

426 So, we have that

$$\begin{cases} \hat{\eta}_{N+k}'' = g_k'(y_k); \\ \tilde{\eta}_{N+k}'' = \frac{-\theta_k(\mu_k - y_k)}{g_k(y_k)}. \end{cases}$$

428 For the 0th component, (5.10) can be expressed as

$$\begin{cases}
\sum_{k=1}^{N} \left(-\hat{\eta}'_{k,N+k} f'_{k}(y_{k}) g_{k}(y_{k}) s^{\beta_{k}} - \hat{\kappa}_{k} f'_{k}(y_{k}) g_{k}(y_{k}) s^{\beta_{k}} \right) \\
+ \sum_{k=1}^{N} \left(\hat{\eta}''_{k} f_{k}(y_{k}) s^{\beta_{k}} - \beta_{k} f^{2}_{k}(y_{k}) s^{2\beta_{k}-1} \right) \\
+ \sum_{k=1}^{N} \left(-\tilde{\eta}'_{k,N+k} f'_{k}(y_{k}) g_{k}(y_{k}) s^{\beta_{k}} + \tilde{\eta}''_{k} f_{k}(y_{k}) s^{\beta_{k}} \right) - rs = 0.
\end{cases}$$

430 We choose to set each term in parenthesis in the first two lines above to zero. It

431 follows that

$$\hat{\eta}'_{k,N+k} = -\hat{\kappa}_k \in L^{\infty}(\Omega), \quad \hat{\eta}''_k = \beta_k f_k(y_k) s^{\beta_k - 1}.$$

434 If N > 1 we (arbitrarily) choose then to set the last line to zero with

435 (5.15)
$$\tilde{\eta}_{k}^{"} = \tilde{\eta}_{k}^{'} = 0, \quad k = 2, \dots, N.$$

436 It remains that

$$\tilde{\eta}_{138}^{"} \quad (5.16) \qquad \qquad \tilde{\eta}_{1}^{"} f_{1}(y_{1}) s^{\beta_{1}} - \tilde{\eta}_{1,N+1}^{"} f_{1}^{"}(y_{1}) g_{1}(y_{1}) s^{\beta_{1}} = rs.$$

Here, we can choose to take either $\tilde{\eta}_1'' = 0$ or $\tilde{\eta}_{1,N+1}' = 0$. We obtain then two

440 possibilities:

441 (5.17)
$$\begin{cases} (i) & \tilde{\eta}_1'' = 0 \text{ and } \tilde{\eta}_{1,N+1}' = \frac{-rs^{1-\beta_1}}{f_1'(y_1)g_1(y_1)}, \\ (ii) & \tilde{\eta}_1'' = \frac{rs^{1-\beta_1}}{f_1(y_1)} \text{ and } \tilde{\eta}_{1,N+1}' = 0. \end{cases}$$

5.2.1. Estimate of the h function. We decide to choose case (i) in (5.17). The function h needs to satisfy (2.28), (2.29), and (3.17) (instead of (2.30)). Instead of (2.28), we will rather check the stronger condition (2.34). We compute

445 (5.18)
$$h'_{\sigma} := \sum_{k=1}^{N} |f_k(y_k)| s^{\beta_k} \left(|(\hat{\kappa}_k)_s| + |\frac{\rho_s}{\rho}| \right) + |g_k(y_k)| \left(|(\hat{\kappa}_k)_k| + |\frac{\rho_k}{\rho}| \right)$$

$$+\sum_{k=1}^{N} \left(\beta_{k} |f_{k}(y_{k}) s^{\beta_{k}-1}| + |g'_{k}(y_{k})|\right),\,$$

447 (5.19)
$$h_r := |r|^{\frac{1}{2}},$$

448 (5.20)
$$h'_{\eta} := \hat{h}'_{\eta} + \tilde{h}'_{\eta},$$

450 where we have

(5.21)

451
$$\hat{h}'_{\eta} := \sum_{k=1}^{N} \left(\beta_k |f_k(y_k)| s^{\beta_k - 1} + |g'_k(y_k)| + \left| f_k(y_k)| s^{\beta_k} \frac{\partial \hat{\kappa}_k}{\partial s} \right| + \left| g_k(y_k) \frac{\partial \hat{\kappa}_k}{\partial y_k} \right| \right),$$

(5.22)

$$\tilde{h}'_{\eta} := \sum_{k=1}^{N} \left| \frac{\theta_k(\mu_k - y_k)}{g_k(y_k)} \right| + \left| r \frac{f_1(y_1)}{f'_1(y_1)g_1(y_1)} \right| + \left| r g_1(y_1) s^{1-\beta_1} \frac{\partial}{\partial y_1} \left[\frac{1}{f'_1(y_1)g_1(y_1)} \right] \right|.$$

REMARK 5.2. Had we chosen (ii) instead of (i) in (5.17), this would only change the expression of \tilde{h}'_n that would then be

456 (5.23)
$$\tilde{h}'_{\eta} = \sum_{k=1}^{N} \left| \frac{\theta_k(\mu_k - y_k)}{g_k(y_k)} \right| + \left| \frac{rs^{1-\beta_1}}{f_1(y_1)} \right|.$$

5.2.2. Estimate of the h function without the commutator analysis. The only change in the estimate of h will be the contribution of h'_{η} and h''_{η} . We have to satisfy (2.28)-(2.30). In addition, ignoring the commutator analysis, we would solve (5.13) with $\hat{\eta}' = 0$, meaning that we choose

461 (5.24)
$$\hat{\eta}_k'' := \beta_k f_k(y_k) s^{\beta_k - 1} + \hat{\kappa}_k \frac{f_k'(y_k) g_k(y_k)}{f_k(y_k)}, \ k = 1, \dots, N,$$

and take $\tilde{\eta}_1''$ out of (5.16). Then condition (3.17), with here $\hat{\eta}' = 0$, would give

463 (5.25)
$$h \ge c_{\eta} h_{\eta}$$
, where $h_{\eta} := h_{\hat{\eta}} + h_{\tilde{\eta}}$,

464 with

465 (5.26)
$$h_{\hat{\eta}} := \sum_{k=1}^{N} \left(\beta_k |f_k(y_k)| s^{\beta_k - 1} + |\hat{\kappa}_k| \left| \frac{f'_k(y_k) g_k(y_k)}{f_k(y_k)} \right| + |g'_k(y_k)| \right),$$

466 (5.27)
$$h_{\tilde{\eta}} := \sum_{k=1}^{N} \left| \frac{\theta_k(\mu_k - y_k)}{g_k(y_k)} \right| + \left| \frac{rs^{1-\beta_1}}{f_1(y_1)} \right|.$$

468 We will see in applications that this is in general worse.

- 6. Application to stochastic volatility models. The results of Section 5 are specified for a subclass of the multiple factor model, in particular for the VAT and GMH models. We show that the commutator analysis allows to take smaller values for the function h (and consequently to include a larger class of payoff functions).
- 6.1. A useful subclass. Here we assume that

474 (6.1)
$$|f_k(y_k)| = |y_k|^{\gamma_k}$$
; $|g_k(y_k)| = \nu_k |y_k|^{1-\gamma_k}$; $\beta_k \in (0,1]$; $\nu_k > 0$; $\gamma_k \in (0,\infty)$.

Furthermore, we assume κ to be constant and

476 (6.2)
$$|f'_k(y_k)g_k(y_k)| = \text{const} \text{ for all } y_k, \ k = 1, \dots, N.$$

477 Set

$$c_{s} := \|s\rho_{s}/\rho\|_{\infty};$$

$$c'_{k} = \begin{cases} \|\rho_{k}/\rho\|_{\infty} & \text{if } \Omega_{k} = \mathbb{R}, \\ 0 & \text{otherwise.} \end{cases}$$

$$c''_{k} = \begin{cases} 0 & \text{if } \Omega_{k} = \mathbb{R}, \\ \|y_{k}\rho_{k}/\rho\|_{\infty} & \text{otherwise.} \end{cases}$$

We get, assuming that $\gamma_1 \neq 0$:

480 (6.4)
$$h'_{\sigma} := \sum_{k=1}^{N} \left(c_{s} |y_{k}|^{\gamma_{k}} s^{\beta_{k}-1} + \nu_{k} c'_{k} |y_{k}|^{1-\gamma_{k}} + \nu_{k} c''_{k} |y_{k}|^{-\gamma_{k}} + \beta_{k} |y_{k}|^{\gamma_{k}} s^{\beta_{k}-1} + (1-\gamma_{k}) \nu_{k} |y_{k}|^{-\gamma_{k}} \right),$$

481

482 (6.5)
$$\hat{h}'_{\eta} := \sum_{k=1}^{N} \left(\beta_k |y_k|^{\gamma_k} s^{\beta_k - 1} + (1 - \gamma_k) \nu_k |y_k|^{-\gamma_k} \right),$$

483 (6.6)
$$\tilde{h}'_{\eta} := \sum_{k=1}^{N} \left(\frac{\theta_{k} |\mu_{k} - y_{k}|}{\nu_{k} |y_{k}|^{1 - \gamma_{k}}} + \frac{r |y_{1}|^{\gamma_{1}}}{\gamma_{1} \nu_{1}} \right).$$

Therefore when all $y_k \in \mathbb{R}$, we can choose h' as

486 (6.7)
$$h' := 1 + \sum_{k=1}^{N} (|y_k|^{\gamma_k} (1 + s^{\beta_k - 1}) + (1 - \gamma_k) |y_k|^{-\gamma_k} + |y_k|^{\gamma_k - 1}) + \sum_{k \in I} |y_k|^{1 - \gamma_k} + \sum_{k \in J} |y_k|^{-\gamma_k}.$$

487 Without the commutator analysis we would get

488 (6.8)
$$\hat{h}_{\eta} := \sum_{k=1}^{N} (\beta_k |y_k|^{\gamma_k} s^{\beta_k - 1} + \nu_k |\hat{\kappa}_k| |y_k|^{-\gamma_k} + (1 - \gamma_k) \nu_k |y_k|^{-\gamma_k}),$$

489 (6.9)
$$\tilde{h}_{\eta} := \sum_{k=1}^{N} \left(\theta_{k} \frac{|\mu_{k} - y_{k}|}{\nu_{k} |y_{k}|^{1 - \gamma_{k}}} + r s^{1 - \beta_{1}} |y_{1}|^{-\gamma_{1}} \right).$$

491 Therefore we can choose

492 (6.10)
$$h := h''; \quad h'' := h' + rs^{1-\beta_1}/|y_1|^{\gamma_1} + \sum_k \nu_k |\hat{\kappa}_k||y_k|^{-\gamma_k}.$$

- So, we always have that $h' \leq h''$, meaning that it is advantageous to use the commu-
- 494 tator analysis, due to the term $rs^{1-\beta_1}/|y_1|^{\gamma_1}$ above in particular. The last term in
- 495 the above r.h.s. has as contribution only when $\gamma_k \neq 1$ (since otherwise h' includes a
- 496 term of the same order).
- 497 **6.2.** Application to the VAT model. For the variant of the Achdou and
- Tchou model with multiple factors (VAT), i.e. when $\gamma_k = 1$, for k = 1 to N, we can
- 499 take h equal to

500 (6.11)
$$h'_{TA} := 1 + \sum_{k=1}^{N} |y_k| (1 + s^{\beta_k - 1}),$$

when the commutator analysis is used, and when it is not, take h equal to

502 (6.12)
$$h_{TA} := h_{TA} + rs^{1-\beta_1} |y_1|^{-1} + \sum_{k=1}^{N} \nu_k |\hat{\kappa}_k| |y_k|^{-1}.$$

- Remember that $u_T(s) = (s K)_+$ for a call option, and $u_T(s) = (K s)_+$ for a put option, both with strike K > 0.
- LEMMA 6.1. For the VAT model, using the commutator analysis, in case of a call (resp. put) option with strike K > 0, we can take $\rho = \rho_{call}$, (resp. $\rho = \rho_{put}$), with

$$\rho_{call}(s,y) := (1 + s^{3+\varepsilon''})^{-1} \Pi_{k=1}^{N} e^{-\varepsilon N(y_k)},
\rho_{put}(s,y) := \frac{s^{\alpha_P}}{1 + s^{\alpha_P}} \Pi_{k=1}^{N} e^{-\varepsilon N(y_k)},$$

508 where
$$\alpha_P := \left(\varepsilon' + 2 \sum_{k=1}^{N} (1 - \beta_k) - 1 \right)_{\perp}$$
.

- 509 Proof. (i) In the case of a call option, we have that
- 510 (6.14) $1 \ge c_0 s^{\beta_k 1}$ for $c_0 > 0$ small enough over the domain of integration,
- 511 so that we can as well take

512 (6.15)
$$h(s,y) = 1 + \sum_{k=1}^{N} |y_k| \le \prod_{k=1}^{N} (1 + |y_k|).$$

513 So, we need that $\varphi(s,y) \in L^{1,\rho}(\Omega)$, with

514 (6.16)
$$\varphi(s,y) = h^2(s,y)u_T^2(s) = (s-K)_+^2 \prod_{k=1}^N (1+|y_k|)^2.$$

By lemma 4.2, where here $J = \{0\}$ and $I = \{1, ..., N\}$, we may take resp.

516 (6.17)
$$\gamma_0' = 2, \ \gamma_0'' = 0, \ \gamma_k' > 0, \ \gamma_k'' > 0, \ k = 1, \dots, N,$$

and so we may choose for $\varepsilon' > 0$ and $\varepsilon'' > 0$:

518 (6.18)
$$\alpha'_0 = 0, \ \alpha''_0 = 3 + \varepsilon'', \ \alpha'_k = \varepsilon', \ \alpha''_k = \varepsilon'', \ k = 1, \dots, N,$$

- 519 so that setting $\varepsilon := \varepsilon' + \varepsilon''$, we can take $\rho = \rho_{call}$.
- 520 (ii) For a put option with strike K>0, $1\leq c_0s^{\beta_k-1}$ for big enough $c_0>0$, over the
- 521 domain of integration, so that we can as well take

522
$$(6.19)$$
 $h(s,y) = 1 + \sum_{k=1}^{N} |y_k| s^{\beta_k - 1} \le \prod_{k=1}^{N} (1 + |y_k| s^{\beta_k - 1})^2 \le \prod_{k=1}^{N} s^{2\beta_k - 2} (1 + |y_k|)^2$

- 523 and
- 524 (6.20) $\varphi(s,y) = h^2(s,y)u_T^2(s) \le (K-s)_+^2 \prod_{k=1}^N s^{2\beta_k-2} (1+|y_k|)^2.$
- 525 By lemma 4.2, in the case of a put option and since $(K-s)_+^2$ is bounded, we can take
- 526 $\gamma'_k, \gamma''_k, \alpha'_k, \alpha''_k$ as before, for k = 1 to N, and

527 (6.21)
$$\gamma_0' = 0, \ \gamma_0'' = 2 \sum_{k=1}^{N} (1 - \beta_k), \ \alpha_0' = \left(\varepsilon' + 2 \sum_{k=1}^{N} (1 - \beta_k) - 1\right)_+, \ \alpha_0'' = 0$$

- 528 the result follows.
- REMARK 6.2. If we do not use the commutator analysis, then we have a greater
- 530 "h" function; we can check that our previous choice of ρ does not apply any more (so
- 531 we should consider a smaller weight function, but we do not need to make it explicit).
- And indeed, we have then a singularity when say y_1 is close to zero so that the previous
- 533 choice of ρ makes the p integral undefined.
- 6.3. Application to the GMH model. For the generalized multiple factor Heston model (GMH), i.e. when $\gamma_k = 1/2$, k = 1 to N, we can take h equal to

536 (6.22)
$$h'_{H} := 1 + \sum_{k=1}^{N} \left(|y_{k}|^{\frac{1}{2}} (1 + s^{\beta_{k} - 1}) + |y_{k}|^{-\frac{1}{2}} \right),$$

537 when the commutator analysis is used, and when it is not, take h equal to

538 (6.23)
$$h_H := h_H + rs^{1-\beta_1}|y_1|^{-\frac{1}{2}}.$$

- LEMMA 6.3. (i) For the GMH model, using the commutator analysis, in case of
- 540 a call option with strike K, meaning that $u_T(s) = (s K)_+$, we can take $\rho = \rho_{call}$,
- 541 with

542 (6.24)
$$\rho_{call}(s,y) := (1 + s^{\varepsilon''+3})^{-1} \prod_{k=1}^{N} y_k^{\varepsilon'} (1 + y_k^{\varepsilon+2})^{-1}.$$

543 (ii) For a put option with strike K > 0, we can take $\rho = \rho_{put}$, with

544 (6.25)
$$\rho_{put}(s,y) := \prod_{k=1}^{N} y_k^{\varepsilon'} (1 + y_k^{\varepsilon+2})^{-1}.$$

545 *Proof.* (i) For the call option, using (6.14) we see that we can as well take

546 (6.26)
$$h(s,y) \le 1 + \sum_{k=1}^{N} \left(y_k^{1/2} + y_k^{-1/2} \right) \le (s - K)_+^2 \prod_{k=1}^{N} (1 + y_k^{1/2} + y_k^{-1/2}).$$

547 So, we need that $\varphi(s,y) \in L^{1,\rho}(\Omega)$, with

548 (6.27)
$$\varphi(s,y) = h^2(s,y)u_T^2(s) = (s-K)_+^2 \prod_{k=1}^N (1+y_k^{1/2}+y_k^{-1/2}).$$

By lemma 4.2, where here $J = \{0, ..., N\}$, we may take resp.

550 (6.28)
$$\gamma_0' = 2, \ \gamma_0'' = 0, \ \gamma_k' = 1, \ \gamma_k'' = 1, \ k = 1, \dots, N,$$

and so we may choose for $\varepsilon' > 0$ and $\varepsilon'' > 0$:

552 (6.29)
$$\alpha'_0 = 0, \ \alpha''_0 = 3 + \varepsilon'', \ \alpha'_k = \varepsilon', \ \alpha''_k = \varepsilon'' + 2, \quad k = 1, \dots, N,$$

- so that setting $\varepsilon := \varepsilon' + \varepsilon''$, we can take $\rho = \rho_{call}$. (ii) For a put option with strike $K > 0, 1 \le c_0 s^{\beta_k 1}$ for big enough $c_0 > 0$, over the 554
- domain of integration, so that we can as well take

556 (6.30)
$$h(s,y) = 1 + \sum_{k=1}^{N} |y_k| s^{\beta_k - 1} \le \prod_{k=1}^{N} (1 + |y_k| s^{\beta_k - 1})^2 \le \prod_{k=1}^{N} s^{2\beta_k - 2} (1 + |y_k|)^2$$

and 557

568

558 (6.31)
$$\varphi(s,y) = h^2(s,y)u_T^2(s) \le (K-s)_+^2 \prod_{k=1}^N s^{2\beta_k - 2} (1 + |y_k|)^2.$$

By lemma 4.2, in the case of a put option and since $(K-s)_+^2$ is bounded, we can take

 $\gamma'_k, \gamma''_k, \alpha'_k, \alpha''_k$ as before, for k = 1 to N, and

561 (6.32)
$$\gamma_0' = 0, \ \gamma_0'' = 0, \ \alpha_0' = 0, \ \alpha_0'' = 0.$$

the result follows. 562

REMARK 6.4. If we do not use the commutator analysis, then, again, we have a 563 greater "h" function; we can check that our previous choice of ρ does not apply any 564 more (so we should consider a smaller weight function, but we do not need to make it 565 explicit). And indeed, by the behaviour of the integral for large s the previous choice 566 567 of ρ makes the p integral undefined.

Appendix A. Regularity results by Lions and Magenes [15, Ch. 1].

Let H be a Hilbert space identified with its dual and scalar product denoted by 569 (\cdot,\cdot) . Let V be a Hilbert space, densely and continuously embedded in H, with duality product denoted by $\langle \cdot, \cdot \rangle_V$. Set

572 (A.1)
$$W(0,T) := \{ u \in L^2(0,T;V); \ \dot{u} \in L^2(0,T;V^*) \}.$$

It is known [15, Ch. 1] that 573

574 (A.2)
$$W(0,T) \subset C(0,T;H)$$
 with continuous inclusion,

and that for any u, v in W(0,T), and $0 \le t < t' \le T$, the following integration by parts formula holds:

577 (A.3)
$$\int_{t}^{t'} (\langle \dot{u}(s), v(s) \rangle_{V} + \langle \dot{v}(s), u(s) \rangle_{V}) \, \mathrm{d}s = (u(t'), v(t'))_{H} - (u(t), v(t))_{H}.$$

Equivalently, 578

579 (A.4)
$$2\int_{t}^{t'} \langle \dot{u}(s), u(s) \rangle_{V} ds = \|u(t')\|_{H}^{2} - \|u(t)\|_{H}^{2}, \text{ for all } u \in W(0, T).$$

- Let $A(t) \in L^{\infty}(0, T; L(V, V^*))$ satisfy the hypotheses of uniform continuity and semicoercivity, i.e., for some $\alpha > 0$, $\lambda \ge 0$, and c > 0:
- Given $(f, u_T) \in L^2(0, T; V^*) \times H$, we consider the following (backward) parabolic equation: find u in W(0, T) such

(A.6)
$$\begin{cases} -\dot{u}(t) + A(t)u(t) = f & \text{in } L^2(0, T; V^*), \\ u(T) = u_T & \text{in } H, \end{cases}$$

- and recall classical results from [15, Ch. 1].
- PROPOSITION A.1 (first parabolic estimate). The parabolic equation (A.6) has a unique solution $u \in W(0,T)$, and for some c > 0 not depending on (f, u_T) :
- 589 (A.7) $||u||_{L^2(0,T;V)} + ||u||_{L^{\infty}(0,T;H)} \le c(||u_T||_H + ||f||_{L^2(0,T;V^*)}).$
- We next derive a stronger result with the hypothesis of *semi-symmetry* below: (A.8)
- $\begin{cases} A(t) = A_0(t) + A_1(t), \ A_0(t) \ \text{and} \ A_1(t) \ \text{continuous linear mappings} \ V \to V^*, \\ A_0(t) \ \text{symmetric and continuously differentiable} \ V \to V^* \ \text{w.r.t.} \ t, \\ A_1(t) \ \text{is measurable with range in} \ H, \ \text{and for positive numbers} \ \alpha_0, \ c_{A,1} \text{:} \\ (i) \quad \langle A_0(t)u,u\rangle_V \ge \alpha_0 \|u\|_V^2, \quad \text{for all} \ u \in V, \ \text{and a.a.} \ t \in [0,T], \\ (ii) \quad \|A_1(t)u\|_H \le c_{A,1} \|u\|_V, \quad \text{for all} \ u \in V, \ \text{and a.a.} \ t \in [0,T], \\ f \in L^2(0,T;H) \ \text{and} \ u_T \in V. \end{cases}$
- PROPOSITION A.2 (second parabolic estimate). Let (A.8) hold. Then the solufor $u \in W(0,T)$ of (A.6) belongs to $L^{\infty}(0,T;V)$, \dot{u} belongs to $L^{2}(0,T;H)$, and for some c>0 not depending on (f,u_{T}) :
- 595 (A.9) $||u||_{L^{\infty}(0,T;V)} + ||\dot{u}||_{L^{2}(0,T;H)} \le c(||u_{T}||_{V} + ||f||_{L^{2}(0,T;H)}).$
- 596 Appendix B. Parabolic variational inequalities.
- Let $K \subset V$ be a non-empty, closed and convex set, K be the closure of K in H, and $u_T \in K$. Let
- 599 (B.1) $\begin{cases} L^2(0,T;K) := \{u \in L^2(0,T;V); \ u(t) \in K \text{ a.e.}\}, \\ W(0,T;K) := W(0,T) \cap L^2(0,T;K). \end{cases}$
- We consider parabolic variational inequalities as follows: find $u \in W(0,T;K)$ such that
- 602 (B.2) $\begin{cases} \langle -\dot{u}(t) + A(t)u(t) f(t), v u(t) \rangle_V \ge 0 & \text{for all } v \in K, \quad \text{a.a. } t, \\ u(T) = u_T & \text{in } H. \end{cases}$
- Take $v \in W(0,T;K)$. Adding to the previous inequality the integration by parts formula
- 605 (B.3) $-\int_0^T \langle \dot{v}(s) \dot{u}(s), v(s) u(s) \rangle_V ds = \frac{1}{2} \|u(0) v(0)\|_H^2 \frac{1}{2} \|u(T) v(T)\|_H^2$

and since $u(T) = u_T$ we find that

$$\begin{cases} \int_{0}^{T} \langle -\dot{v}(t) + A(t)u(t) - f(t), v - u(t) \rangle_{V} \ge \frac{1}{2} \|u(0) - v(0)\|_{H}^{2} - \frac{1}{2} \|u(T) - v(T)\|_{H}^{2} \\ \text{for all } v \in W(0, T; K), \quad u(T) = u_{T}. \end{cases}$$

- 608 It can be proved that the two formulation (B.2) and (B.4) are equivalent (they have the
- 609 same set of solutions), and that they have at most one solution. The weak formulation
- 610 is as follows: find $u \in L^2(0,T;K) \cap C(0,T;H)$ such that

(B.5)
$$\begin{cases} \int_0^T \langle -\dot{v}(t) + A(t)u(t) - f(t), v - u(t) \rangle_V \ge -\frac{1}{2} \|u(T) - v(T)\|_H^2 \\ \text{for all } v \in L^2(0, T; K), \quad u(T) = u_T. \end{cases}$$

- 612 Clearly a solution of the strong formulation (B.2) is solution of the weak one.
- PROPOSITION B.1 (Brézis [6]). The following holds:
- 614 (i) Let $u_T \in \mathcal{K}$ and $f \in L^2(0,T;V^*)$. Then the weak formulation (B.5) has a unique solution u and, for some c > 0, given $v_0 \in K$:
- 616 (B.6) $||u||_{L^{\infty}(0,T;H)} + ||u||_{L^{2}(0,T;V)} \le c(||u_{T}||_{H} + ||f||_{L^{2}(0,T;V^{*})} + ||v_{0}||_{V}).$
- 617 (ii) Let in addition the semi-symmetry hypothesis (A.8) hold, and let u_T belong 618 to K. Then $u \in L^{\infty}(0,T;V)$, $\dot{u} \in L^2(0,T;H)$, and u is the unique solution 619 of the original formulation (B.2). Furthermore, for some c > 0:

620 (B.7)
$$||u||_{L^{\infty}(0,T;V)} + ||\dot{u}||_{L^{2}(0,T;H)} \le c(||u_{T}||_{V} + ||f||_{L^{2}(0,T;H)}).$$

Appendix C. Monotonicity. Assume that H is an Hilbert lattice, i.e., is endowed with an order relation \succeq compatible with the vector space structure:

623 (C.1)
$$x_1 \succeq x_2$$
 implies that $\gamma x_1 + x \succeq \gamma x_2 + x$, for all $\gamma \geq 0$ and $x \in H$,

- 624 such that the maxima and minima denoted by $\max(x_1, x_2)$ and $\min(x_1, x_2)$ are well
- defined, the operator max, min be continuous, with $\min(x_1, x_2) = -\max(-x_1, -x_2)$.
- Setting $x_+ := \max(x, 0)$ and $x_- := -\min(x, 0)$ we have that $x = x_+ x_-$. Assuming
- 627 that the maximum of two elements of V belong to V we see that we have an induced
- lattice structure on V. The induced dual order over V^* is as follows: for v_1^* and v_2^* in
- 629 V^* , we say that $v_1^* \ge v_2^*$ if $\langle v_1^* v_2^*, v \rangle_V \ge 0$ whenever $v \ge 0$.
- Assume that we have the following extension of the integration by parts formula (B.3): for all u, v in W(0,T) and $0 \le t < t' \le T$,

632 (C.2)
$$2 \int_{t}^{t'} \langle \dot{u}(s), u_{+}(s) \rangle_{V} ds = \|u_{+}(t')\|_{H}^{2} - \|u_{+}(t)\|_{H}^{2}.$$

633 and that

634 (C.3)
$$\langle A(t)u, u_{+}\rangle_{V} = \langle A(t)u_{+}, u_{+}\rangle_{V}.$$

PROPOSITION C.1. Let u_i be solution of the parabolic equation (A.6) for $(f, u_T) = (f^i, u_T^i)$, i = 1, 2. If $f^1 \ge f2$ and $u_T^1 \ge u_T^2$, then $u_1 \ge u_2$.

This type of result may be extended to the case of variational inequalities. If K 637 and K' are two subsets of V, we say that K dominates K' if for any $u \in K$ and 638 $u' \in K'$, $\max(u, u') \in K$ and $\min(u, u') \in K'$. 639

Proposition C.2. Let u_i be solution of the weak formulation (B.5) of the parabolic 640 variational inequality for $(f, u_T, K) = (f^i, u_T^i, K^i)$, i = 1, 2. If $f_1 \ge f_2$, $u_T^1 \ge u_T^2$, and 641 K^1 dominates K^2 , then $u_1 > u_2$. 642

The monotonicity w.r.t. the convex K is due to Haugazeau [13] (in an elliptic 643 setting, but the result is easily extended to the parabolic one). See also Brézis [7].

Appendix D. Link with American options. 645 An American option is the right to get a payoff $\Psi(t,x)$ at any time t < T and u_T at time T. We can motivate 646 as follows the derivation of the associated variational inequalities. If the option can 647 be exercized only at times $t_k = hk$, with h = T/M and k = 0 to M (Bermudean 648 option), then the same PDE as for the European option holds over $(t_k, t_{k+1}), k=0$ 649 650 to M-1. Denoting by \tilde{u}_k the solution of this PDE, we have that $u(t_k) = \max(\Psi, \tilde{u}_k)$. Assuming that A does not depend on time and that there is a flux f(t,x) of dividents, 651 we compute the approximation u_k of $u(t_k)$ as follows. Discretizing the PDE with the 652 implicit Euler scheme we obtain the continuation value \hat{u}_k solution of 653

654 (D.1)
$$\frac{\hat{u}_k - u_{k+1}}{h} + A\hat{u}_k = f(t_k, \cdot), \ k = 0, \dots, M - 1; \quad u_M = \max(\Psi, 0),$$

so that $u_k = u_{k+1} - hA\hat{u}_k + hf(t_k, \cdot)$, we find that 655

656 (D.2)
$$u_k = \max(\hat{u}_k, \Psi) = \max(u_{k+1} - hA\hat{u}_k + hf(t_k, \cdot), \Psi),$$

which is equivalent to

644

658 (D.3)
$$\min(u_k - \Psi, \frac{u_k - u_{k+1}}{h} + A\hat{u}_k - f(t_k, \cdot)) = 0.$$

This suggest for the continuous time model and general operators A and r.h.s. f the 659 following formulation 660

661 (D.4)
$$\min(u(t,x) - \Psi(x), -\dot{u}(t,x) + A(t,x)u(t,x) - f(t,x)) = 0$$
, $(t,x) \in (0,T) \times \Omega$.

The above equation has a rigorous mathematical sense in the context of viscosity 662 663 solution, see Barles [5]. However we rather need the variational formulation which can be derived as follows. Let v(x) satisfy $v(x) \geq \Psi(x)$ a.e., be smooth enough. Then 664

$$\int_{\Omega} \left(-\dot{u}(t,x) + A(t,x)u(t,x) - f(t,x) \right) (v(x) - u(t,x)) dx =
\int_{\{u(t,x) = \Psi(x)\}} \left(-\dot{u}(t,x) + A(t,x)u(t,x) - f(t,x) \right) (v(x) - u(t,x)) dx
+ \int_{\{u(t,x) > \Psi(x)\}} \left(-\dot{u}(t,x) + A(t,x)u(t,x) - f(t,x) \right) (v(x) - u(t,x)) dx.$$

The first integrand is nonnegative, being a product of nonnegative terms, and the 666 second integrand is equal to 0 since by (D.3), $-\dot{u}(t,x) + A(t,x)u(t,x) - f(t,x) = 0$ 667 a.e. when $u(t,x) > \Psi(x)$. So we have that, for all $v \geq \Psi$ smooth enough: 668

669 (D.6)
$$\int_{\Omega} \left(-\dot{u}(t,x) + A(t,x)u(t,x) - f(t,x) \right) \left(v(x) - u(t,x) \right) dx \ge 0.$$

- We see that this is of the same nature as a parabolic variational inequality, where K is the set of functions greater or equal to Ψ (in an appropriate Sobolev space).
- Appendix E. Some one dimensional problems. It is not always easy to characterize the space V. Let us give a detailed analysis in a simple case.
- **E.1. The Black-Scholes setting.** For the Black-Scholes model with zero interest rate (the extension to a constant nonzero interest rate is easy) and unit volatility coefficient, we have that $Au = -\frac{1}{2}x^2u''(x)$, with $x \in (0, \infty)$. In the case of a put option: $u_T(x) = (K x)_+$ we may take $H := L^2(\mathbb{R}_+)$. For $v \in \mathcal{D}(0, \infty)$ and u sufficiently smooth we have that $-\frac{1}{2}\int_0^\infty x^2u''(x)\mathrm{d}x = a(u,v)$ with

679 (E.1)
$$a(u,v) := \frac{1}{2} \int_0^\infty x^2 u'(x) v'(x) dx + \int_0^\infty x u'(x) v(x) dx.$$

This bilinear form a is continuous and semi coercive over the set

681 (E.2)
$$V := \{ u \in H; \ xu'(x) \in H \}.$$

- It is easily checked that $\bar{u}(x) := x^{-1/3}/(1+x)$ belongs to V. So, some elements of V are unbounded near zero.
- We now claim that $\mathcal{D}(0,\infty)$ is a dense subset of V. First, it follows from a standard truncation argument and the dominated convergence theorem that $V_{\infty} := V \cap L^{\infty}(0,\infty)$ is a dense subset of V. Note that elements of V are continuous over $(0,\infty)$. Given $\varepsilon > 0$ and $u \in V_{\infty}$, define
- 688 (E.3) $u_{\varepsilon}(x) := \begin{cases} 0 & \text{if } x \in (0, \varepsilon), \\ u(2\varepsilon)(x/\varepsilon 1) & \text{if } x \in [\varepsilon, 2\varepsilon], \\ u(2\varepsilon) & \text{if } x > 2\varepsilon. \end{cases}$

Obviously
$$u_{\varepsilon} \in V_{\infty}$$
. By the dominated convergence theorem, $u_{\varepsilon} \to u$ in H . Set for $w \in V$

691 (E.4)
$$\Phi_{\varepsilon}(w) := \int_0^{2\varepsilon} x^2 w'(x)^2 \mathrm{d}x.$$

692 Since Φ_{ε} is quadratic and $v_{\varepsilon} \to u$ in H, we have that:

693 (E.5)
$$\frac{1}{2} \int_0^\infty x^2 (u_{\varepsilon}' - u')^2 dx = \frac{1}{2} \Phi_{\varepsilon}(u_{\varepsilon} - u) \le \Phi_{\varepsilon}(u_{\varepsilon}) + \Phi_{\varepsilon}(u).$$

Since $u \in V$, $\Phi_{\varepsilon}(u) \to 0$ and

695 (E.6)
$$\Phi_{\varepsilon}(u_{\varepsilon}) \leq \|u\|_{\infty}^{2} \int_{0}^{2\varepsilon} \varepsilon^{-2} x^{2} dx = O(\|u\|_{\infty}^{2} \varepsilon).$$

- So, the l.h.s. of (E.5) has limit 0 when $\varepsilon \downarrow 0$. We have proved that the set V^0 of functions in V^{∞} equal to zero near zero, is a dense subset of V. Now define for N > 0
- 698 (E.7) $\varphi_N(x) = \begin{cases} 1 & \text{if } x \in (0, N), \\ 1 \log(x/N) & \text{if } x \in [N, eN], \\ 0 & \text{if } x > eN. \end{cases}$
- Given $u \in V_0$, set $u_N := u\varphi_N$. Then $u_N \in H$ and, by a dominated convergence argument, $u_N \to u$ in H. The weak derivative of u_N is $u'_N = u'\varphi_N + u\varphi'_N$. By a

dominated convergence argument, $xu'\varphi_N \to xu'$ in $L^2(\mathbb{R}_+)$. It remains to prove that $xu\varphi'_N \to 0$ in $L^2(\mathbb{R}_+)$. But φ'_N is equal to 1/x over its support, so that

703 (E.8)
$$||xu\varphi_N'||_{L^2(\mathbb{R}_+)}^2 = \int_N^{eN} u^2(x) dx \le \int_N^\infty u^2(x) dx \to 0$$

- 704 when $N \uparrow \infty$. The claim is proved.
- 705 **E.2. The CIR setting.** In the Cox-Ingersoll-Ross model [8] the stochastic pro-706 cess satisfies

707 (E.9)
$$ds(t) = \theta(\mu - s(t))dt + \sigma\sqrt{s} dW(t), \quad t \ge 0$$

- 708 We assume the coefficients θ , μ and σ to be constant and positive. The associated
- 709 PDE is given by

710 (E.10)
$$\begin{cases} Au := -\theta(\mu - x)u' - \frac{1}{2}x\hat{\sigma}^2 u'' = 0 & (x, t) \in \mathbb{R}^+ \times (0, T), \\ u(x, T) = u_T(x) & x \in \mathbb{R}^+. \end{cases}$$

- Again for the sake of simplicity we will take $\rho(x) = 1$, which is well-adapted in the
- case of a payoff with compact support in $(0, \infty)$. For $v \in \mathcal{D}(0, \infty)$ and u sufficiently
- 713 smooth we have that $\int_0^\infty Au(x)v(x) dx = a(u,v)$ with

714
$$a(u,v) := \theta \int_0^\infty (\mu - x)u'(x)v(x)dx + \frac{1}{2}\hat{\sigma}^2 \int_0^\infty xu'(x)v'(x)dx + \frac{1}{2}\hat{\sigma}^2 \int_0^\infty u'(x)v(x)dx.$$

715 So one should take V of the form

716 (E.12)
$$\mathcal{V} := \{ u \in H; \ \sqrt{x}u'(x) \in L^2(\mathbb{R}_+) \}.$$

- 717 We next determine H by requiring that the bilinear form is continuous; by the Cauchy-
- 718 Schwarz inequality

719
$$\left| \int_0^\infty u'(x)v(x) dx \right| \le \|x^{1/2}u'\|_2 \|x^{-1/2}v\|_2; \quad \left| \int_0^\infty xu'(x)v(x) dx \right| \le \|x^{1/2}u'\|_2 \|x^{1/2}v\|_2.$$

- 720 We easily deduce that the bilinear form a is continuous and semi coercive over \mathcal{V} ,
- 721 when choosing

722 (E.14)
$$H := \{ v \in L^2(\mathbb{R}_+); \ (x^{1/2} + x^{-1/2})v \in L^2(\mathbb{R}_+) \},$$

Note that then the integrals below are well defined and finite for any $v \in \mathcal{V}$:

724 (E.15)
$$\int_0^\infty (x^{1/2}v')(x^{-1/2}v) = \int_0^\infty vv' = \frac{1}{2} \int_0^\infty (v^2)'.$$

- 725 So $w := v^2$ is the primitive of an integrable function and therefore has a limit at zero.
- Since v is continuous over $(0, \infty)$ it follows that v has a limit at zero.
- However if this limit is nonzero we get a contradiction with the condition that
- 728 $x^{-1/2}v \in L^2(\mathbb{R}_+)$. So, every element of \mathcal{V} has zero value at zero.
- We now claim that $\mathcal{D}(0,\infty)$ is a dense subset of \mathcal{V} . First, $\mathcal{V}_{\infty} := \mathcal{V} \cap L^{\infty}(0,\infty)$ is
- 730 a dense subset of \mathcal{V} . Note that elements of \mathcal{V} are continuous over $(0,\infty)$. Given $\varepsilon>0$

and $u \in \mathcal{V}_{\infty}$, define $u_{\varepsilon}(x)$ as in (E.3). Then $u_{\varepsilon} \in \mathcal{V}_{\infty}$. By the dominated convergence

732 theorem, $u_{\varepsilon} \to u$ in H. Set for $w \in \mathcal{V}$

733 (E.16)
$$\Phi_{\varepsilon}(w) := \int_{0}^{2\varepsilon} xw'(x)^{2} dx.$$

734 Since Φ_{ε} is quadratic and $u_{\varepsilon} \to u$ in H, we have that:

735 (E.17)
$$\frac{1}{2} \int_0^\infty x^2 (u_{\varepsilon}' - u')^2 dx = \frac{1}{2} \Phi_{\varepsilon}(u_{\varepsilon} - u) \le \Phi_{\varepsilon}(u_{\varepsilon}) + \Phi_{\varepsilon}(u).$$

736 Since $u \in \mathcal{V}$, $\Phi_{\varepsilon}(u) \to 0$ and

737 (E.18)
$$\Phi_{\varepsilon}(u_{\varepsilon}) \leq \varepsilon^{-2} u(2\varepsilon)^2 \int_0^{2\varepsilon} x dx = 2u(2\varepsilon)^2 \to 0.$$

- 738 So, the l.h.s. of (E.17) has limit 0 when $\varepsilon \downarrow 0$. We have proved that the set \mathcal{V}^0 of
- functions in \mathcal{V}^{∞} equal to zero near zero, is a dense subset of \mathcal{V} . Define φ_N as in (E.7)
- Given $u \in \mathcal{V}_0$, set $u_N := u\varphi_N$. As before, $u_N \to u$ in H, is $u'_N = u'\varphi_N + u\varphi'_N$,
- 741 $xu'\varphi_N \to xu$ in $L^2(\mathbb{R}_+)$, and it remains to prove that $xu\varphi'_N \to 0$ in $L^2(\mathbb{R}_+)$. But φ'_N
- 742 is equal to 1/x over its support, so that when $N \uparrow \infty$:

743 (E.19)
$$||x^{1/2}u\varphi_N'||_{L^2(\mathbb{R}_+)}^2 = \int_N^{eN} x^{-1}u^2(x)\mathrm{d}x \le \int_N^\infty u^2(x)\mathrm{d}x \to 0.$$

744 The claim is proved.

Appendix F. Dependence of the solution w.r.t. the weights. Let $\rho = \rho_1 \rho_2$, we index the spaces H, V, W(0,T) as $V_1 := V_{\rho_1}, V := V_{\rho}$, etc.. Then $V_1 \subset V$ with dense inclusion, since $\mathcal{D}(\Omega)$ is a dense subspace of V. So, $V^* \subset V_1$ with dense inclusion. In addition, let $\hat{v} \in \mathcal{D}(\Omega)$ and set $v := \hat{v} \rho_2$: since our weights belong to $C^{\infty}(\Omega)$ we have that $v \in \mathcal{D}(\Omega)$. It is not difficult to check that, for all $u \in V_1$:

750 (F.1)
$$a_{\rho_1}(u,v) = a_{\rho}(u,\hat{v})$$

Let $f \in L^2(0,T;V^*)$ and $u \in W_1(0,T)$ be solution of the variational formulation with the weight ρ_1 . Taking v as a test function we see that for a.a. t:

753 (F.2)
$$-\langle \dot{u}(t), v \rangle_1 + a_{\rho_1}(u(t), v) = \langle J^* f(t), v \rangle_1 = \langle f(t), Jv \rangle;$$

we denote by $\langle \cdot, \cdot \rangle_1$ (resp. $\langle \cdot, \cdot \rangle$) the duality product in V_1 (resp. in V). If J denote the injection of V_1 in V, we have for $f \in V^*$ and $v_1 \in V_1$ that

756 (F.3)
$$\langle f, Jv_1 \rangle = \langle J^*f, v_1 \rangle_1$$

and if $f \in H^*$ and $J^*f \in H_1^*$ (spaces identified with their dual) we deduce that

758 (F.4)
$$\int_{\Omega} f v_1 \rho = \int_{\Omega} J^* f v_1 \rho_1$$

759 implying that $J^*f = \rho_2 f$.

760 REFERENCES

- [1] Y. Achdou, B. Franchi, and N. Tchou, A partial differential equation connected to option
 pricing with stochastic volatility: regularity results and discretization, Math. Comp. 74
 (2005), no. 251, 1291–1322.
 - [2] Y. Achdou and O. Pironneau, Computational methods for option pricing, Society for Industrial and Applied Mathematics (SIAM), Philadelphia, PA, 2005.
 - [3] Y. Achdou and N. Tchou, Variational analysis for the Black and Scholes equation with stochastic volatility, M2AN Math. Model. Numer. Anal. 36 (2002), no. 3, 373–395.
 - [4] T. Aubin, A course in differential geometry, Graduate Studies in Mathematics, vol. 27, American Mathematical Society, Providence, RI, 2001.
 - [5] G. Barles, Convergence of numerical schemes for degenerate parabolic equations arising in finance theory, Numerical methods in finance, Publ. Newton Inst., vol. 13, Cambridge Univ. Press, Cambridge, 1997, pp. 1–21.
 - [6] H. Brézis, Inéquations variationnelles paraboliques, Séminaire Jean Leray.

764

765

766

767

768

769

770

771

772

773

774

 $775 \\ 776$

777

 $778 \\ 779$

780

 $781 \\ 782$

783

784

785

786

787 788

789

790

791

792

793 794

795

796

797

798

799

800

801

802

- [7] ______, Problèmes unilatéraux, Journal de Mathématiques pures et appliquées **51** (1972), 1–168.
- [8] J.C. Cox, J. E. Ingersoll, Jr., and S. A. Ross, A theory of the term structure of interest rates, Econometrica 53 (1985), no. 2, 385–407. MR 785475
- [9] P.M. N. Feehan and C. A. Pop, Schauder a priori estimates and regularity of solutions to boundary-degenerate elliptic linear second-order partial differential equations, J. Differential Equations 256 (2014), no. 3, 895–956. MR 3128929
- [10] P.M.N. Feehan and C. A. Pop, Degenerate-elliptic operators in mathematical finance and higher-order regularity for solutions to variational equations., Adv. Differ. Equ. 20 (2015), no. 3-4, 361–432 (English).
- [11] J.-P. Fouque, G. Papanicolaou, and K.R. Sircar, Derivatives in financial markets with stochastic volatility, Cambridge University Press, Cambridge, 2000.
- [12] P. Gauthier and D. Possamaï, Efficient simulation of the double Heston model, Tech. report, SSRN working paper series, 2009, Revised in Jan. 2010.
- [13] Y. Haugazeau, Sur des inéquations variationnelles, C. R. Acad. Sci. Paris Sér. A-B 265 (1967), A95–A98.
- [14] S.L. Heston, A closed-form solution for options with stochastic volatility with applications to bond and currency options, The Review of Financial Studies 6 (1993), no. 2, 327–343.
- [15] J.-L. Lions and E. Magenes, Non-homogeneous boundary value problems and applications. Vol. I, Springer-Verlag, New York, 1972, Translated from the French by P. Kenneth, Die Grundlehren der mathematischen Wissenschaften, Band 181.
- [16] P.-L. Lions and M. Musiela, Correlations and bounds for stochastic volatility models, Ann. Inst. H. Poincaré Anal. Non Linéaire 24 (2007), no. 1, 1–16.
- [17] K. Jacobs P. Christoffersen, S. Heston, The shape and term structure of the index option smirk: Why multifactor stochastic volatility models work so well, Management Science 55 (2009), no. 12, 1914–1932.
- [18] O. Pironneau and Y. Achdou, Partial differential equations for option pricing, Handbook of numerical analysis Vol XV. Special Volume: Mathematical modeling and numerical methods in finance, Amsterdam: Elsevier/North-Holland, 2009, pp. 369–495.