## City, University of London Institutional Repository

Citation: Černý, A. \& Kallsen, J. (2007). On the structure of general mean-variance hedging strategies. Annals of Probability, 35(4), pp. 1479-1531. doi: 10.1214/009117906000000872

This is the published version of the paper.

This version of the publication may differ from the final published version.

Permanent repository link: http://openaccess.city.ac.uk/17876/
Link to published version: http://dx.doi.org/10.1214/009117906000000872

Copyright and reuse: City Research Online aims to make research outputs of City, University of London available to a wider audience. Copyright and Moral Rights remain with the author(s) and/or copyright holders. URLs from City Research Online may be freely distributed and linked to.

[^0]
# ON THE STRUCTURE OF GENERAL MEAN-VARIANCE HEDGING STRATEGIES 

By Aleš ČERNÝ and Jan Kallsen<br>City University London and Technische Universität München

We provide a new characterization of mean-variance hedging strategies in a general semimartingale market. The key point is the introduction of a new probability measure $P^{\star}$ which turns the dynamic asset allocation problem into a myopic one. The minimal martingale measure relative to $P^{\star}$ coincides with the variance-optimal martingale measure relative to the original probability measure $P$.

## Contents

1. Introduction
1.1. Overview
1.2. Semimartingale characteristics and notation
2. Admissible strategies and quadratic hedging
2.1. Admissible strategies
2.2. Mean-variance hedging
3. On the pure investment problem
3.1. Opportunity process
3.2. Adjustment process
3.3. Variance-optimal signed martingale measure
3.4. Opportunity-neutral measure
3.5. Characterization of $L$ and $\tilde{a}$
3.6. When does $P^{\star}=P$ hold?
3.7. Determination of the opportunity process
4. On the pure investment problem
4.1. Mean value process and pure hedge coefficient
4.2. Main results
4.3. Connections to the literature

Appendix
A.1. Locally square-integrable semimartingales
A.2. $\sigma$-Martingales

References

[^1]
## 1. Introduction.

1.1. Overview. In incomplete market models perfect replication of contingent claims is typically impossible. A classical way out is to minimize the mean squared hedging error

$$
E\left(\left(v+\vartheta \cdot S_{T}-H\right)^{2}\right)
$$

over all reasonable hedging strategies $\vartheta$ and possibly all initial endowments $v$. Here, the random variable $H$ denotes the discounted payoff of the claim, the semimartingale $S$ stands for the discounted price process of the underlying, the dot refers to stochastic integration, and $T$ is the time horizon. Mathematically speaking, one seeks to compute the orthogonal projection of $H$ on some space of stochastic integrals.

This problem has been extensively studied both as far as general theory as well as concrete results in specific setups are concerned. In order to render equal justice (or rather injustice) to most contributions, we refer the reader to [38] and [45] for excellent overviews of the literature. More recent publications in this context include $[2-4,7-11,17,23-26,32-36,46]$.

The purpose of this piece of research is to provide a deeper understanding of the structure of the mean-variance hedging problem in a general semimartingale context. More specifically, we aim at concrete formulas for the objects of interestto the extent that this is possible without restricting to more specific situations.

If $S$ is a square-integrable martingale, the answer to the above hedging problem is provided by the Galtchouk-Kunita-Watanabe decomposition of the claim (cf. [19]). In particular, the optimal hedge $\vartheta$ is of the form

$$
\begin{equation*}
\vartheta_{t}=\frac{d\langle V, S\rangle_{t}}{d\langle S, S\rangle_{t}} \tag{1.1}
\end{equation*}
$$

where $V_{t}=E\left(H \mid \mathscr{F}_{t}\right)$ denotes the martingale generated by the contingent claim $H$.
If $S$ fails to be a martingale, the hedging problem becomes much more involved. Relatively explicit results have been obtained by Schweizer [42] under the condition of deterministic mean-variance tradeoff, which can be intepreted as a certain homogeneity property of the asset price process $S$. In this case the optimal hedge is the sum of two terms. The first satisfies an equation resembling (1.1). The second can be interpreted in terms of a pure investment problem under quadratic utility.

In the current paper we reduce the general case to the expressions of [42]. This is done by a specific nonmartingale change of measure. If the formulas of [42] are evaluated relative to the new opportunity-neutral measure $P^{\star}$ rather than $P$, they yield the optimal hedge relative to the original probability measure $P$. We discuss the links to the literature more thoroughly in Section 4.3.

The paper is structured as follows. Section 2 explains the setup of the meanvariance problem at hand. In particular, we define a notion of admissibility which ensures the existence of an optimal hedge. The measure change alluded to above
and related objects are introduced in Section 3. Subsequently, we turn to the hedging problem itself. Finally, the appendix contains and summarizes auxiliary statements on semimartingales. In particular, we prove a sufficient condition for square integrability of exponential semimartingales which is needed in Section 4.
1.2. Semimartingale characteristics and notation. Unexplained notation is typically used as in [28]. Superscripts refer generally to coordinates of a vector or vector-valued process rather than powers. The few exceptions should be obvious from the context. If $X$ is a semimartingale, $L(X)$ denotes the set of $X$-integrable predictable processes in the sense of [28], III.6.17.

In the subsequent sections, optimal hedging strategies are expressed in terms of semimartingale characteristics.

Definition 1.1. Let $X$ be an $\mathbb{R}^{d}$-valued semimartingale with characteristics $(B, C, v)$ relative to some truncation function $h: \mathbb{R}^{d} \rightarrow \mathbb{R}^{d}$. By [28], II.2.9 there exists some predictable process $A \in \mathscr{A}_{\text {loc }}^{+}$, some predictable $\mathbb{R}^{d \times d}$-valued process $c$ whose values are nonnegative, symmetric matrices, and some transition kernel $F$ from $\left(\Omega \times \mathbb{R}_{+}, \mathscr{P}\right)$ into $\left(\mathbb{R}^{d}, \mathscr{B}^{d}\right)$ such that

$$
B_{t}=b \cdot A_{t}, \quad C_{t}=c \cdot A_{t}, \quad \nu([0, t] \times G)=F(G) \cdot A_{t}
$$

$$
\text { for } t \in[0, T], G \in \mathscr{B}^{d} .
$$

We call $(b, c, F, A)$ differential characteristics of $X$.
One should observe that the differential characteristics are not unique: for example, $\left(2 b, 2 c, 2 F, \frac{1}{2} A\right)$ yields another version. Especially for $A_{t}=t$, one can interpret $b_{t}$ or rather $b_{t}+\int(x-h(x)) F_{t}(d x)$ as a drift rate, $c_{t}$ as a diffusion coefficient, and $F_{t}$ as a local jump measure. The differential characteristics are typically derived from other "local" representations of the process, for example, in terms of a stochastic differential equation.

From now on, we choose the same fixed process $A$ for all the (finitely many) semimartingales in this paper. The results do not depend on its particular choice. In concrete models, $A$ is often taken to be $A_{t}=t$ (e.g., for Lévy processes, diffusions, Itô processes, etc.) and $A_{t}=[t]:=\max \{n \in \mathbb{N}: n \leq t\}$ (discrete-time processes). Since almost all semimartingales of interest in this paper are actually special semimartingales, we use from now on the (otherwise forbidden) "truncation" function

$$
h(x):=x,
$$

which simplifies a number of expressions considerably.
By $\langle X, Y\rangle$ we denote the $P$-compensator of $[X, Y]$ provided that $X, Y$ are semimartingales such that [ $X, Y$ ] is $P$-special (cf. [27], page 37). If $X$ and $Y$ are vectorvalued, then $[X, Y]$ and $\langle X, Y\rangle$ are to be understood as matrix-valued processes with components $\left[X^{i}, Y^{j}\right.$ ] and $\left\langle X^{i}, Y^{j}\right\rangle$, respectively. Moreover, if both $Y$ and a
predictable process $\vartheta$ are $\mathbb{R}^{d}$-valued, then the notation $\vartheta \bullet[X, Y]$ (and accordingly $\vartheta \cdot\langle X, Y\rangle)$ refers to the vector-valued process whose components $\vartheta \cdot\left[X^{i}, Y\right]$ are the vector-stochastic integral of $\left(\vartheta^{j}\right)_{j=1, \ldots, d}$ relative to $\left(\left[X^{i}, Y^{j}\right]\right)_{j=1, \ldots, d}$. If $P^{\star}$ denotes another probability measure, we write $\langle X, Y\rangle^{P^{\star}}$ for the $P^{\star}$-compensator of $[X, Y]$.

In the whole paper, we write $M^{X}$ for the local martingale part and $A^{X}$ for the predictable part of finite variation in the canonical decomposition

$$
X=X_{0}+M^{X}+A^{X}
$$

of a special semimartingale $X$. If $P^{\star}$ denotes another probability measure, we write accordingly

$$
X=X_{0}+M^{X \star}+A^{X \star}
$$

for the $P^{\star}$-canonical decomposition of $X$.
If $(b, c, F, A)$ denote differential characteristics of an $\mathbb{R}^{d}$-valued special semimartingale $X$, we use the notation $\tilde{c}, \hat{c}$ for modified second characteristics in the following sense (provided that the integrals exist):

$$
\begin{align*}
& \tilde{c}:=c+\int x x^{\top} F(d x),  \tag{1.2}\\
& \hat{c}:=c+\int x x^{\top} F(d x)-b b^{\top} \Delta A . \tag{1.3}
\end{align*}
$$

Observe that $x^{\top} \hat{c} x \leq x^{\top} \tilde{c} x$ for any $x \in \mathbb{R}^{d}$. The notion of modified second characteristics is motivated by the following:

Proposition 1.2. Let $X$ be an $\mathbb{R}^{d}$-valued special semimartingale with differential characteristics ( $b, c, F, A$ ) and modified second characteristics as in (1.2) and (1.3). If the corresponding integrals exist, then

$$
\begin{aligned}
\langle X, X\rangle & =\tilde{c} \cdot A, \\
\left\langle M^{X}, M^{X}\right\rangle & =\hat{c} \cdot A .
\end{aligned}
$$

Proof. The first equation follows from [28], I.4.52, the second from [28], II.2.17 (adjusted for the truncation function).

From now on we use the notation $\left(b^{X}, c^{X}, F^{X}, A\right)$ to denote differential characteristics of a special semimartingale $X$. Accordingly, $\tilde{c}^{X}, \hat{c}^{X}$ stands for the modified second characteristics of $X$. If they refer to some probability measure $P^{\star}$ rather than $P$, we write instead $\left(b^{X \star}, c^{X \star}, F^{X \star}, A\right)$ and $\tilde{c}^{X \star}, \hat{c}^{X \star}$, respectively. We denote the joint characteristics of two special semimartingales $X, Y$ [i.e., the characteristics of $(X, Y)]$ as

$$
\left(b^{X, Y}, c^{X, Y}, F^{X, Y}, A\right)=\left(\binom{b^{X}}{b^{Y}},\left(\begin{array}{cc}
c^{X} & c^{X Y} \\
c^{Y X} & c^{Y}
\end{array}\right), F^{X, Y}, A\right)
$$

and

$$
\tilde{c}^{X, Y}=\left(\begin{array}{cc}
\tilde{c}^{X} & \tilde{c}^{X Y} \\
\tilde{c}^{Y X} & \tilde{c}^{Y}
\end{array}\right), \quad \hat{c}^{X, Y}=\left(\begin{array}{cc}
\hat{c}^{X} & \hat{c}^{X Y} \\
\hat{c}^{Y X} & \hat{c}^{Y}
\end{array}\right) .
$$

In the whole paper, we write $c^{-1}$ for the Moore-Penrose pseudoinverse of a matrix or matrix-valued process $c$, which is a particular matrix satisfying $c c^{-1} c=c$ (cf. [1]). From the construction it follows that the mapping $c \mapsto c^{-1}$ is measurable. Moreover, $c^{-1}$ is nonnegative and symmetric if this holds for $c$.

Finally, we write $X \sim Y$ (resp. $X \sim^{\star} Y$ ) if two semimartingales differ only by some $P-\sigma$-martingale (or some $P^{\star}-\sigma$-martingale, resp.). Some facts on $\sigma$-martingales are summarized in Appendix A.2.
2. Admissible strategies and quadratic hedging. We work on a filtered probability space $\left(\Omega, \mathscr{F},\left(\mathscr{F}_{t}\right)_{t \in[0, T]}, P\right)$, where $T \in \mathbb{R}_{+}$denotes a fixed terminal time. The $\mathbb{R}^{d}$-valued process $S=\left(S_{t}^{1}, \ldots, S_{t}^{d}\right)_{t \in[0, T]}$ represents the discounted prices of $d$ securities. We assume that

$$
\begin{equation*}
\sup \left\{E\left(\left(S_{\tau}^{i}\right)^{2}\right): \tau \text { stopping time }, i=1, \ldots, d\right\}<\infty \tag{2.1}
\end{equation*}
$$

that is, $S$ is a $L^{2}(P)$-semimartingale in the sense of [15].
Moreover, we make the following standing:
ASSUMPTION 2.1. There exists some equivalent $\sigma$-martingale measure with square-integrable density, that is, some probability measure $Q \sim P$ with $E\left(\left(\frac{d Q}{d P}\right)^{2}\right)<\infty$ and such that $S$ is a $Q-\sigma$-martingale.

This can be interpreted as a natural no-free-lunch condition in the present quadratic context. More specifically, Théorème 2 in [47] and standard arguments show that Assumption 2.1 is equivalent to the absence of $L^{2}$-free lunches in the sense that

$$
\overline{K_{2}^{s}(0)-L_{+}^{2}} \cap L_{+}^{2}=\{0\}
$$

where $K_{2}^{s}(0)$ denotes the set of payoffs of simple trading defined below, $L_{+}^{2}$ contains the nonnegative square-integrable random variables, and the closure is to be taken in $L^{2}(P)$.
2.1. Admissible strategies. The choice of the set of admissible trading strategies in continuous time is a delicate point. If it is too large, arbitrage opportunities occur even in the Black-Scholes model, if it is too small, optimal strategies as, for example, the replicating portfolio of a European call in the Black-Scholes model fail to exist. Inspired by Delbaen and Schachermayer [15], we consider the closure (in a proper $L^{2}$-sense) of the set of simple strategies.

More specifically, an $\mathbb{R}^{d}$-valued process $\vartheta$ is called simple if it is a linear combination of processes of the form $Y 1_{\rrbracket} \tau_{1}, \tau_{2} \rrbracket$, where $\tau_{1} \leq \tau_{2}$ denote stopping times
and $Y$ a bounded $\mathscr{F}_{\tau_{1}}$-measurable random variable. We call a payoff attainable by simple trading with initial endowment $v \in L^{2}\left(\Omega, \mathscr{F}_{0}, P\right)$ if it belongs to the set

$$
K_{2}^{s}(v):=\left\{v+\vartheta \cdot S_{T}: \vartheta \text { simple }\right\}
$$

If the initial endowment $v$ is not fixed beforehand, we consider instead the set

$$
K_{2}^{s}\left(\mathscr{F}_{0}\right):=\left\{v+\vartheta \cdot S_{T}: v \in L^{2}\left(\Omega, \mathscr{F}_{0}, P\right), \vartheta \text { simple }\right\} .
$$

Since the hedging problems in this paper concern the approximation of arbitrary payoffs $H$ in $L^{2}(P)$ by attainable outcomes, it makes perfect sense from an economical point of view to call the elements of the $L^{2}(P)$-closures $K_{2}(v):=\overline{K_{2}^{s}(v)}$, respectively, $K_{2}\left(\mathscr{F}_{0}\right):=\overline{K_{2}^{S}\left(\mathscr{F}_{0}\right)}$ attainable as well. These outcomes can be written as a stochastic integral $v+\vartheta \cdot S_{T}$ with some strategy $\vartheta \in L(S)$ that can be approximated in the following sense by simple strategies (cf. Lemmas 2.4 and 2.6 below).

DEFINITION 2.2. We call $\vartheta \in L(S)$ admissible strategy if there exists some sequence $\left(\vartheta^{(n)}\right)_{n \in \mathbb{N}}$ of simple strategies such that

$$
\begin{aligned}
& \vartheta^{(n)} \cdot S_{t} \rightarrow \vartheta \cdot S_{t} \quad \text { in probability for any } t \in[0, T] \quad \text { and } \\
& \vartheta^{(n)} \cdot S_{T} \rightarrow \vartheta \cdot S_{T} \quad \text { in } L^{2}(P) .
\end{aligned}
$$

Similarly, we call $(v, \vartheta) \in L^{0}\left(\Omega, \mathscr{F}_{0}, P\right) \times L(S)$ admissible endowment/strategy pair if there exist some sequences $\left(v^{(n)}\right)_{n \in \mathbb{N}}$ in $L^{2}\left(\Omega, \mathscr{F}_{0}, P\right)$ and $\left(\vartheta^{(n)}\right)_{n \in \mathbb{N}}$ of simple strategies such that

$$
\begin{aligned}
v^{(n)}+\vartheta^{(n)} \cdot S_{t} & \rightarrow v+\vartheta \cdot S_{t} \\
v^{(n)}+\vartheta^{(n)} \cdot S_{T} & \rightarrow v+\vartheta \cdot S_{T}
\end{aligned} \quad \text { in probability for any } t \in[0, T] \quad \text { and } L^{2}(P) . ~ \$ ~ l
$$

We set

$$
\begin{aligned}
\bar{\Theta} & :=\{\vartheta \in L(S): \vartheta \text { admissible }\}, \\
\overline{L^{2}\left(\mathscr{F}_{0}\right) \times \Theta}: & :\left\{(v, \vartheta) \in L^{0}\left(\Omega, \mathscr{F}_{0}, P\right) \times L(S):(v, \vartheta) \text { admissible }\right\} .
\end{aligned}
$$

One easily verifies that $\overline{L^{2}\left(\mathscr{F}_{0}\right) \times \Theta}=\mathbb{R} \times \bar{\Theta}$ if the initial $\sigma$-field $\mathscr{F}_{0}$ is trivial. Admissible strategies are linked via duality to martingale measures of the following kind:

DEFINITION 2.3. We call a signed measure $Q \ll P$ with $Q(\Omega)=1 a b-$ solutely continuous signed $\sigma$-martingale measure ( $S \sigma M M$ ) if $S Z^{Q}$ is a $P-\sigma$-martingale for the density process

$$
Z_{t}^{Q}:=E\left(\left.\frac{d Q}{d P} \right\rvert\, \mathscr{F}_{t}\right)
$$

of $Q$.

A probability measure $Q \sim P$ is a $\mathrm{S} \sigma \mathrm{MM}$ if and only if $S$ is a $Q-\sigma$-martingale (cf. Lemma A.8).

Lemma 2.4. For $H \in L^{2}(P)$ and $v \in L^{2}\left(\Omega, \mathscr{F}_{0}, P\right)$ the following statements are equivalent:

1. $H \in K_{2}(v)$.
2. $E_{Q}(H-v)=0$ for any $S \sigma M M Q$ with $\frac{d Q}{d P} \in L^{2}(P)$.
3. $H=v+\vartheta \cdot S_{T}$ with some $\vartheta \in \bar{\Theta}$.
4. $H=v+\vartheta \cdot S_{T}$ with some $\vartheta \in L(S)$ such that $(\vartheta \cdot S) Z^{Q}$ is a martingale for any $S \sigma M M Q$ with density process $Z^{Q}$ and $\frac{d Q}{d P} \in L^{2}(P)$.
In particular, we have $K_{2}(v)=\left\{v+\vartheta \cdot S_{T}: \vartheta \in \bar{\Theta}\right\}$.
Proof. It suffices to consider the case $v=0$.
$1 \Rightarrow 3$, 4: Step 1: We start by showing that statement 4 holds for $H \in K_{2}^{s}(0)$, that is, for $H=\vartheta \cdot S_{T}$ with some simple $\vartheta$. Integration by parts yields

$$
\begin{align*}
(\vartheta \cdot S) Z^{Q} & =(\vartheta \cdot S)_{-} \cdot Z^{Q}+\vartheta \cdot\left(Z_{-}^{Q} \cdot S+\left[Z^{Q}, S\right]\right) \\
& =\left(\vartheta \cdot S_{-}-\vartheta^{\top} S_{-}\right) \cdot Z^{Q}+\vartheta \cdot\left(S Z^{Q}\right), \tag{2.2}
\end{align*}
$$

which implies that $(\vartheta \cdot S) Z^{Q}$ is a $\sigma$-martingale. Since $\sup _{t \in[0, T]}\left|Z_{t}^{Q}\right| \in L^{2}(P)$ by Doob's inequality and $\vartheta \cdot S$ is a $L^{2}$-semimartingale in the sense of (2.1), we have that $(\vartheta \cdot S) Z^{Q}$ is of class ( D$)$ and hence a martingale (cf. Lemma A.7).

Step 2: Let $H^{n}=\vartheta^{(n)} \cdot S_{T}$ be an approximating sequence for $H \in K_{2}(0)$. From [15], Theorem 1.2, it follows that $H$ has a representation $H=\vartheta \cdot S_{T}$ for some $\vartheta \in L(S)$. In the proof of this theorem it is actually shown that $\vartheta$ can be chosen such that $\vartheta^{(n)} \cdot S_{t}$ converges in probability to $\vartheta \cdot S_{t}$ for any $t \in[0, T]$.

Since $H^{n} Z_{T}^{Q} \rightarrow H Z_{T}^{Q}$ in $L^{1}(P)$, we have that

$$
E\left(\left(\vartheta^{(n)} \cdot S_{T}\right) Z_{T}^{Q} \mid \mathscr{F}_{t}\right) \rightarrow E\left(\left(\vartheta \cdot S_{T}\right) Z_{T}^{Q} \mid \mathscr{F}_{t}\right)
$$

in $L^{1}(P)$ and hence in probability. Step 1 yields $E\left(\left(\vartheta^{(n)} \cdot S_{T}\right) Z_{T}^{Q} \mid \mathscr{F}_{t}\right)=\left(\vartheta^{(n)} \cdot\right.$ $\left.S_{t}\right) Z_{t}^{Q}$. Together, it follows that $E\left(\left(\vartheta \cdot S_{T}\right) Z_{T}^{Q} \mid \mathscr{F}_{t}\right)=\left(\vartheta \cdot S_{t}\right) Z_{t}^{Q}$.
$3 \Rightarrow 1$ : This is obvious.
$4 \Rightarrow 2$ : This is obvious as well.
$2 \Rightarrow 1$ : It suffices to show that $K_{2}(0)^{\perp} \subset\left(V^{\perp}\right)^{\perp}$ for

$$
V:=\left\{\frac{d Q}{d P}: Q \text { S } \sigma \mathrm{MM} \text { with } \frac{d Q}{d P} \in L^{2}(P)\right\},
$$

where the orthogonal complements refer to $L^{2}(P)$. Let $Y \in K_{2}(0)^{\perp}$ and set $Z_{t}:=$ $E\left(Y \mid \mathscr{F}_{t}\right)$. For $s \leq t$ and $F \in \mathscr{F}_{s}$ we have

$$
\begin{aligned}
& E\left(1_{F}\left(S_{t} Z_{t}-S_{s} Z_{s}\right)\right) \\
& \quad=E\left(1_{F}\left(S_{t}-S_{s}\right) Y\right)-E\left(1_{F} S_{t}\left(Z_{T}-Z_{t}\right)\right)+E\left(1_{F} S_{s}\left(Z_{T}-Z_{s}\right)\right) \\
& \quad=0
\end{aligned}
$$

because $Z$ is a martingale and $1_{F \times(s, t]} \cdot S_{T} \in K_{2}(0)$. If $E(Y) \neq 0$, then $Y$ is a multiple of a $\mathrm{S} \sigma \mathrm{MM}$ and hence in $\left(V^{\perp}\right)^{\perp}$. If $E(Y)=0$, then $Y+\frac{d Q}{d P} \in V \subset\left(V^{\perp}\right)^{\perp}$ for the $\mathrm{S} \sigma$ MM $Q$ from Assumption 2.1, which implies that $Y \in\left(V^{\perp}\right)^{\perp}$ as well.

This leads to the following characterization of admissible strategies:
COROLLARY 2.5. We have equivalence between:

1. $\vartheta$ is an admissible strategy.
2. $\vartheta \in L(S), \vartheta \cdot S_{T} \in L^{2}(P)$, and $(\vartheta \cdot S) Z^{Q}$ is a martingale for any $S \sigma M M Q$ with density process $Z^{Q}$ and $\frac{d Q}{d P} \in L^{2}(P)$.

Proof. $\quad 1 \Rightarrow 2$ : This follows from the argument in step 2 of the proof of Lemma 2.4.
$2 \Rightarrow 1$ : We have $\vartheta \cdot S_{T} \in K_{2}(0)$ by Lemma 2.4. Let $Q$ be a $\sigma$-martingale measure as in Assumption 2.1. By the proof of Lemma $2.4(1 \Rightarrow 3,4)$ there exists some $\tilde{\vartheta} \in \bar{\Theta}$ such that $\tilde{\vartheta} \cdot S$ is a $Q$-martingale with $\tilde{\vartheta} \cdot S_{T}=\vartheta \cdot S_{T}$. Since $\vartheta \cdot S$ is a $Q$-martingale as well, we have $\tilde{\vartheta} \cdot S=\vartheta \cdot S$ and hence $\vartheta \in \bar{\Theta}$.

In the case without fixed initial endowment we have:
Lemma 2.6. There exists

1. $K_{2}\left(\mathscr{F}_{0}\right)=\left\{v+\vartheta \cdot S_{T}:(v, \vartheta) \in \overline{L^{2}\left(\mathscr{F}_{0}\right) \times \Theta}\right\}$.
2. If $(v, \vartheta) \in \overline{L^{2}\left(\mathscr{F}_{0}\right) \times \Theta}$, then $(v+\vartheta \cdot S) Z^{Q}$ is a martingale for any $S \sigma M M Q$ with density process $Z^{Q}$ and $\frac{d Q}{d P} \in L^{2}(P)$.

Proof. This follows by rather obvious extension of the proof of Lemma 2.4 $(1 \Rightarrow 3,4)$ and the underlying arguments in [15].

REMARK 2.7. An inspection of the proof reveals that statement 2 in Corollary 2.5 and Lemma 2.6 holds for any square-integrable martingale $Z^{Q}$ such that $S Z^{Q}$ is a $\sigma$-martingale, that is, the property $E\left(Z_{T}^{Q}\right)=1$ is not needed.

If necessary the whole setup can be relaxed to slightly more general price processes:

REMARK 2.8. Instead of (2.1), Delbaen and Schachermayer [15] assume only that $S$ is a local $L^{2}(P)$-semimartingale, that is, that there is a localizing sequence of stopping times $\left(U_{n}\right)_{n \in \mathbb{N}}$ such that:

$$
\sup \left\{E\left(\left(S_{\tau}^{i}\right)^{2}\right): \tau \leq U_{n} \text { stopping time, } i=1, \ldots, d\right\}<\infty
$$

for any $n \in \mathbb{N}$. Equivalently, $S^{1}, \ldots, S^{d}$ are locally square-integrable semimartingales (cf. Definition A. 1 and Lemma A. 2 in the Appendix). In this case Delbaen and Schachermayer [15] call a linear combination of strategies $Y 1_{\rrbracket \tau_{1}, \tau_{2} \rrbracket}$ simple if the corresponding stopping times $\tau_{1} \leq \tau_{2}$ are dominated by some $U_{n}$. One easily verifies that all results in this paper extend to this slightly more general setup.

The corresponding admissible sets $\bar{\Theta}$ and $\overline{L^{2}\left(\mathscr{F}_{0}\right) \times \Theta}$ from Definition 2.2 do not depend on the chosen sequence $\left(U_{n}\right)_{n \in \mathbb{N}}$ : For $\bar{\Theta}$ this follows from the characterization in Corollary 2.5. Moreover, $K_{2}\left(\mathscr{F}_{0}\right)=\overline{L^{2}\left(\Omega, \mathscr{F}_{0}, P\right)+K_{2}(0)}$ does not depend on $\left(U_{n}\right)_{n \in \mathbb{N}}$ by Lemma 2.4. Using Lemma 2.6 and arguing similarly as in the proof of Corollary $2.5(2 \Rightarrow 1)$, we have that the same is true for $\overline{L^{2}\left(\mathscr{F}_{0}\right) \times \Theta}$.

Many results in the subsequent sections could also be expressed in terms of the generally different set of strategies considered in [42] and other papers on meanvariance hedging, namely

$$
\Theta:=\left\{\vartheta \in L(S): \vartheta \cdot S \in \mathscr{S}^{2}\right\}
$$

where $\mathscr{S}^{2}$ denotes the set of square-integrable semimartingales (cf. Definition A.1). In contrast to $\left\{v+\vartheta \cdot S_{T}: \vartheta \in \bar{\Theta}\right\}$, the set $\left\{v+\vartheta \cdot S_{T}: \vartheta \in \Theta\right\}$ is not necessarily closed. This issue is discussed in detail by Monat and Stricker [37], Delbaen et al. [13] and Choulli, Krawczyk and Stricker [12]. By considering $L^{2}$-closures in the above sense, one avoids the problem that optimal hedging strategies may fail to exist. In the context of continuous processes, our notion of admissible strategies coincides with the one of Gourieroux, Laurent and Pham [20] and Laurent and Pham [30]. Recently, the question of how to choose a reasonable set of strategies in a quadratic context has been discussed by Xia and Yan [48]. Their notion of admissibility differs from ours but their set of terminal payoffs coincides with $K_{2}(0)$.

The relationship between $\bar{\Theta}$ and $\Theta$ is clarified by the following result. The first assertion is inspired by a similar statement in Grandits and Rheinläender [21], Lemma 2.1 for continuous processes. Loosely speaking, it says that $\bar{\Theta}$ is a kind of $L^{2}$-closure of $\Theta$.

## Corollary 2.9. We have

1. $\Theta \subset \bar{\Theta}$ and $\overline{\left\{\vartheta \cdot S_{T}: \vartheta \in \Theta\right\}}=K_{2}(0)=\left\{\vartheta \cdot S_{T}: \vartheta \in \bar{\Theta}\right\}$.
2. $L^{2}\left(\Omega, \mathscr{F}_{0}, P\right) \times \Theta \subset \overline{L^{2}\left(\mathscr{F}_{0}\right) \times \Theta}$ and

$$
\begin{aligned}
& \quad \overline{\left\{v+\vartheta \cdot S_{T}: v \in L^{2}\left(\Omega, \mathscr{F}_{0}, P\right), \vartheta \in \Theta\right\}} \\
& \quad=K_{2}\left(\mathscr{F}_{0}\right)=\left\{v+\vartheta \cdot S_{T}:(v, \vartheta) \in \overline{L^{2}\left(\mathscr{F}_{0}\right) \times \Theta}\right\}
\end{aligned}
$$

In both cases the closure $\overline{\{\cdots\}}$ refers to the $L^{2}(P)$-norm.

Proof. 1. For $\vartheta \in \Theta$ we have $E\left(\sup _{t \in[0, T]}\left|\vartheta \cdot S_{t}\right|^{2}\right)<\infty$ by Protter [39], Theorem IV.5. $\vartheta \in \bar{\Theta}$ now follows easily from Corollary $2.5(2 \Rightarrow 1)$ together with (2.2) and Lemma A.7. The second equality is shown in Lemma 2.4. In order to verify the first equality, it suffices to prove that any simple strategy is in $\Theta$. This may not be true in the first place. But if the sequence $\left(U_{n}\right)_{n \in \mathbb{N}}$ in Remark 2.8 is chosen such that $\left(S^{i}\right)^{U_{n}} \in \mathscr{S}^{2}$ for $n \in \mathbb{N}, i=1, \ldots, d$, then $\vartheta \in \Theta$ for any simple $\vartheta$. Since $\bar{\Theta}$ does not depend on the chosen sequence $\left(U_{n}\right)_{n \in \mathbb{N}}$, the claim follows.
2. By statement 1 we have

$$
L^{2}\left(\Omega, \mathscr{F}_{0}, P\right) \times \Theta \subset L^{2}\left(\Omega, \mathscr{F}_{0}, P\right) \times \bar{\Theta} \subset \overline{L^{2}\left(\mathscr{F}_{0}\right) \times \Theta} .
$$

The equalities follow similarly as in statement 1 , this time using Lemma 2.6.
2.2. Mean-variance hedging. The goal of this paper is to hedge a fixed contingent claim with discounted payoff $H \in L^{2}(\Omega, \mathscr{F}, P)$. We consider two closely related optimization problems.

DEFINITION 2.10. 1. We call an admissible endowment/strategy pair $\left(v_{0}, \varphi\right)$ optimal if $(v, \vartheta)=\left(v_{0}, \varphi\right)$ minimizes the expected squared hedging error

$$
\begin{equation*}
E\left(\left(v+\vartheta \cdot S_{T}-H\right)^{2}\right) \tag{2.3}
\end{equation*}
$$

over all admissible endowment/strategy pairs $(v, \vartheta)$.
2. If the initial endowment $v=v_{0} \in L^{2}\left(\Omega, \mathscr{F}_{0}, P\right)$ is given beforehand, a minimizer $\vartheta=\varphi$ of (2.3) over all $\vartheta \in \bar{\Theta}$ is called optimal hedging strategy for given initial endowment $v_{0}$.

Due to the chosen notion of admissibility, optimal hedges always exist:
Lemma 2.11. There exist optimal hedges in the sense of Definition 2.10(1) and (2). In both cases, the value process $v_{0}+\varphi \cdot S$ of the optimal hedge is unique up to a $P$-null set.

Proof. The existence follows from Lemmas 2.4, 2.6 and the closedness of $K_{2}\left(\mathscr{F}_{0}\right)$ and $K_{2}\left(v_{0}\right)$, respectively.

Denote by $v_{0}+\varphi \cdot S$ and $\tilde{v}_{0}+\tilde{\varphi} \cdot S$ value processes of two optimal hedges [which implies that $v_{0}=\tilde{v}_{0}$ in the situation of Definition 2.10(2)]. A simple convexity argument yields $v_{0}+\varphi \cdot S_{T}=\tilde{v}_{0}+\tilde{\varphi} \cdot S_{T}$. It remains to be shown that this implies $v_{0}+\varphi \cdot S=\tilde{v}_{0}+\tilde{\varphi} \cdot S$ up to a $P$-null set. Otherwise, there exists some $n \in \mathbb{N}$ such that $P(\tau<T)>0$ for the stopping time

$$
\tau:=\inf \left\{t \in[0, T]: v_{0}+\varphi \cdot S_{t} \geq \tilde{v}_{0}+\tilde{\varphi} \cdot S_{t}+\frac{1}{n}\right\} \wedge T
$$

(or possibly with exchanged roles of $\varphi, \tilde{\varphi}$ ). From Corollary 2.5 and Lemma 2.6 if follows that $M:=v_{0}-\tilde{v}_{0}+(\varphi-\tilde{\varphi}) \cdot S$ is a martingale with respect to the $\sigma$-martingale measure $Q$ from Assumption 2.1. Consequently, $E_{Q}\left(M_{\tau}\right)=$ $E_{Q}\left(M_{T}\right)=0$, which is impossible if $P(\tau<T)>0$.
3. On the pure investment problem. In many papers the mean-variance hedging problem is partially reduced to pure portfolio optimization with quadratic utility. This is done here as well.
3.1. Opportunity process. In the spirit of Markowitz, we call an admissible strategy $\lambda^{(\tau)}$ efficient on a stochastic interval $\rrbracket \tau, T \rrbracket$ if it minimizes

$$
\begin{equation*}
E\left(\left(1-\vartheta \cdot S_{T}\right)^{2}\right) \tag{3.1}
\end{equation*}
$$

over all $\vartheta \in \bar{\Theta}$ vanishing on $\llbracket 0, \tau \rrbracket$. Indeed, by standard arguments there exists no strategy with at most the same variance yielding a higher expected return. Alternatively, one may view $\lambda^{(\tau)}$ as optimal hedging strategy on $\rrbracket \tau, T \rrbracket$ for the constant option $H=1$. A crucial role will be played by the related opportunity process

$$
L_{t}=E\left(\left(1-\lambda^{(t)} \cdot S_{T}\right)^{2} \mid \mathscr{F}_{t}\right),
$$

whose existence and properties are yet to be derived.
LEMMA 3.1. 1. For any stopping time $\tau$ there exists an efficient strategy $\lambda^{(\tau)}$ on $\rrbracket \tau, T \rrbracket$. Its value process $1-\lambda^{(\tau)} \cdot S$ is uniquely determined.
2. $1-\lambda^{(\varrho)} \cdot S_{\tau}=\left(1-\lambda^{(\varrho)} \cdot S_{\sigma}\right)\left(1-\lambda^{(\sigma)} \cdot S_{\tau}\right)$ for all stopping times $\varrho \leq \sigma \leq \tau$.
3. If $1-\lambda^{(\sigma)} \cdot S_{\tau}=0$, then $1-\lambda^{(\sigma)} \cdot S_{T}=0$ for all stopping times $\sigma \leq \tau$.
4. $E\left(\left(1-\lambda^{(\tau)} \cdot S_{T}\right)^{2} \mid \mathscr{F}_{\sigma}\right) \leq E\left(\left(1-\vartheta \cdot S_{T}\right)^{2} \mid \mathscr{F}_{\sigma}\right)$ for all stopping times $\sigma \leq \tau$ and any strategy $\vartheta \in \bar{\Theta}$ with $\vartheta 1_{\llbracket 0, \tau \rrbracket}=0$.
5. $E\left(1-\lambda^{(\tau)} \cdot S_{T} \mid \mathscr{F}_{\sigma}\right)=E\left(\left(1-\lambda^{(\tau)} \cdot S_{T}\right)^{2} \mid \mathscr{F}_{\sigma}\right) \in(0,1]$ almost surely for all stopping times $\sigma \leq \tau$.

Proof. 1. If $G$ denotes the orthogonal projection of 1 on

$$
\overline{\left\{\vartheta \cdot S_{T}: \vartheta \in \bar{\Theta} \text { and } \vartheta 1_{\llbracket 0, \tau \rrbracket}=0\right\}} \subset K_{2}(0)
$$

then there is a sequence $\left(\vartheta^{(n)}\right)_{n \in \mathbb{N}}$ of strategies in $\bar{\Theta}$ that vanish on $\llbracket 0, \tau \rrbracket$ and satisfy $\vartheta^{(n)} \cdot S_{T} \rightarrow G$ in $L^{2}(P)$. By Lemma 2.4 we have $G=\vartheta \cdot S_{T}$ for some $\vartheta \in \bar{\Theta}$. Moreover, $\vartheta^{(n)} \cdot S_{T} \rightarrow \vartheta \cdot S_{T}$ in $L^{1}(Q)$ for the $\sigma$-martingale measure $Q$ from Assumption 2.1. This implies $0=\vartheta^{(n)} \cdot S_{t} \rightarrow \vartheta \cdot S_{t}$ in $L^{1}(Q)$ because both $\vartheta^{(n)} \cdot S$ and $\vartheta \cdot S_{t}$ are $Q$-martingales by Corollary 2.5 . Hence we have $\vartheta 1_{\llbracket 0, \tau \rrbracket}=0$ without loss of generality. Uniqueness follows as in the proof of Lemma 2.11.
2. We start by showing that

$$
\begin{align*}
& E\left(\left(1-\lambda^{(\varrho)} \cdot S_{T}\right)^{2} \mid \mathscr{F}_{\sigma}\right) \\
& \quad \leq E\left(\left(1-\left(\lambda^{(\varrho)} 1_{\llbracket 0, \sigma \rrbracket}+\left(1-\lambda^{(\varrho)} \cdot S_{\sigma}\right) \vartheta\right) \cdot S_{T}\right)^{2} \mid \mathscr{F}_{\sigma}\right) \tag{3.2}
\end{align*}
$$

holds almost surely for any $\vartheta \in \bar{\Theta}$ with $\vartheta 1_{\llbracket 0, \sigma \rrbracket}=0$. Otherwise, there exists some $\vartheta \in \bar{\Theta}$ with $\vartheta 1_{\llbracket 0, \sigma \rrbracket}=0$ such that the reverse inequality holds on some set $F \in \mathscr{F}_{\sigma}$ with $P(F)>0$. Define the strategy

$$
\psi:= \begin{cases}\lambda^{(\varrho)} 1_{\llbracket 0, \sigma \rrbracket}+\left(1-\lambda^{(\varrho)} \cdot S_{\sigma}\right) \vartheta, & \text { on } F, \\ \lambda^{(\varrho)}, & \text { on } F^{C} .\end{cases}
$$

We have

$$
\begin{align*}
E((1- & \left.\left.\psi \cdot S_{T}\right)^{2}\right) \\
= & E\left(E\left(\left(1-\lambda^{(\varrho)} \cdot S_{T}\right)^{2} \mid \mathscr{F}_{\sigma}\right) 1_{F}{ }^{c}\right)  \tag{3.3}\\
& +E\left(E\left(\left(1-\left(\lambda^{(\varrho)} 1_{\llbracket 0, \sigma \rrbracket}+\left(1-\lambda^{(\varrho)} \cdot S_{\sigma}\right) \vartheta\right) \cdot S_{T}\right)^{2} \mid \mathscr{F}_{\sigma}\right) 1_{F}\right) \\
< & E\left(\left(1-\lambda^{(\varrho)} \cdot S_{T}\right)^{2}\right) .
\end{align*}
$$

This contradicts the optimality of $\lambda^{(\varrho)}$ if $\psi \in \bar{\Theta}$.
In order to show $\psi \in \bar{\Theta}$, let $Z$ be the density process of some $\mathrm{S} \sigma \mathrm{MM}$ with square-integrable density. Integration by parts yields that $(\psi \cdot S) Z$ is a $\sigma$-martingale [cf. (2.2)]. Since $P\left(\left|\lambda^{(\varrho)} \cdot S_{\sigma}\right| \leq n\right) \uparrow 1$ and $P\left(\left|\vartheta \cdot S_{\sigma}\right| \leq n\right) \uparrow 1$ for $n \uparrow \infty$, we may assume w.l.o.g. that $\left|\lambda^{(\varrho)} \cdot S_{\sigma}\right|$ and $\left|\vartheta \cdot S_{\sigma}\right|$ are bounded on $F$, say by $n \in \mathbb{N}$. On $\rrbracket \sigma, T \rrbracket$ we have

$$
\begin{aligned}
\mid(\psi \cdot & \left.S-\lambda^{(\varrho)} \cdot S\right) Z \mid \\
& \leq\left(\left|\lambda^{(\varrho)} \cdot S-\lambda^{(\varrho)} \cdot S_{\sigma}\right|+\left|1-\lambda^{(\varrho)} \cdot S_{\sigma}\right|\left|\vartheta \cdot S-\vartheta \cdot S_{\sigma}\right|\right)|Z| 1_{F} \\
& \leq\left(\left|\left(\lambda^{(\varrho)} \cdot S\right) Z\right|+n|Z|+(n+1)(|(\vartheta \cdot S) Z|+n)\right) 1_{F} .
\end{aligned}
$$

The processes in the last line are of class (D) by Corollary 2.5. This in turn implies that $(\psi \cdot S) Z$ is of class (D) as well and hence a martingale. Another application of Corollary 2.5 yields $\psi \in \bar{\Theta}$. Thus (3.3) yields a true contradiction, which means that (3.2) holds.

Note that (3.2) implies

$$
\begin{equation*}
E\left(\left.\left(1-\frac{\lambda^{(\varrho)} 1 \rrbracket \sigma, T \rrbracket}{1-\lambda^{(\varrho)} \cdot S_{\sigma}} \cdot S_{T}\right)^{2} \right\rvert\, \mathscr{F}_{\sigma}\right) \leq E\left(\left(1-\vartheta \cdot S_{T}\right)^{2} \mid \mathscr{F}_{\sigma}\right) \tag{3.4}
\end{equation*}
$$

almost surely on $\left\{1-\lambda^{(\varrho)} \cdot S_{\sigma} \neq 0\right\}$ for any $\vartheta \in \bar{\Theta}$ with $\vartheta 1_{\llbracket 0, \sigma \rrbracket}=0$. Moreover, we have on the set $\left\{1-\lambda^{(\varrho)} \cdot S_{\sigma}=0\right\}$ that

$$
E\left(\left(1-\lambda^{(\varrho)} \cdot S_{T}\right)^{2} \mid \mathscr{F}_{\sigma}\right) \leq E\left(\left(1-\left(\lambda^{(\varrho)} 1_{\llbracket 0, \sigma \rrbracket}\right) \cdot S_{T}\right)^{2} \mid \mathscr{F}_{\sigma}\right)=0
$$

and hence $1-\lambda^{(\varrho)} \cdot S_{T}=0$.
Similarly as (3.2), one shows that

$$
\begin{equation*}
E\left(\left(1-\lambda^{(\sigma)} \cdot S_{T}\right)^{2} \mid \mathscr{F}_{\sigma}\right) \leq E\left(\left.\left(1-\left(\alpha \frac{\lambda^{(\varrho)} 1 \rrbracket \sigma, T \rrbracket}{1-\lambda^{(\varrho)} \cdot S_{\sigma}}+\vartheta\right) \cdot S_{T}\right)^{2} \right\rvert\, \mathscr{F}_{\sigma}\right) \tag{3.5}
\end{equation*}
$$

holds almost surely on $\left\{1-\lambda^{(\varrho)} \cdot S_{\sigma} \neq 0\right\}$ for any $\alpha \in \mathbb{R}_{+}$and any $\vartheta \in \bar{\Theta}$ with $\vartheta 1_{\llbracket 0, \sigma \rrbracket}=0$. Using a convexity argument, (3.4) and (3.5) yield that

$$
1-\lambda^{(\sigma)} \cdot S_{T}=1-\frac{\lambda^{(\varrho)} 1 \rrbracket \sigma, T \rrbracket}{1-\lambda^{(\varrho)} \cdot S_{\sigma}} \cdot S_{T}
$$

on $\left\{1-\lambda^{(\varrho)} \cdot S_{\sigma} \neq 0\right\}$ and hence

$$
\lambda^{(\sigma)} \cdot S_{T}\left(1-\lambda^{(\varrho)} \cdot S_{\sigma}\right)=\left(\lambda^{(\varrho)} 1_{\rrbracket \sigma, T \rrbracket}\right) \cdot S_{T} .
$$

By taking conditional expectation relative to the $\sigma$-martingale measure $Q$ from Assumption 2.1, it follows that

$$
\lambda^{(\sigma)} \cdot S_{\tau}\left(1-\lambda^{(\varrho)} \cdot S_{\sigma}\right)=\left(\lambda^{(\varrho)} 1_{\rrbracket \sigma, T \rrbracket}\right) \cdot S_{\tau}
$$

for any $\tau \geq \sigma$ (cf. Corollary 2.5), which yields the claim.
3. This is shown in the proof of statement 2.
4. This follows from (3.2) for $\varrho=\sigma$.
5. If $E\left(\left(1-\lambda^{(\tau)} \cdot S_{T}\right)^{2} \mid \mathscr{F}_{\sigma}\right)=0$ on some set $F \in \mathscr{F}_{\sigma}$ with $P(F)>0$, then

$$
\lambda^{(\tau)} \cdot S_{T}-\lambda^{(\tau)} \cdot S_{\sigma}=1
$$

which contradicts the fact that $\lambda^{(\tau)} \cdot S$ is a $Q$-martingale for the $\sigma$-martingale measure $Q$ from Assumption 2.1 (cf. Corollary 2.5). Hence, $E\left(\left(1-\lambda^{(\tau)} \cdot S_{T}\right)^{2} \mid\right.$ $\left.\mathscr{F}_{\sigma}\right)>0$ almost surely. Moreover,

$$
\begin{aligned}
E\left(\left(1-(1+\varepsilon) \lambda^{(\tau)} \cdot S_{T}\right)^{2} \mid \mathscr{F}_{\sigma}\right)= & E\left(\left(1-\lambda^{(\tau)} \cdot S_{T}\right)^{2} \mid \mathscr{F}_{\sigma}\right) \\
& -2 \varepsilon E\left(\lambda^{(\tau)} \cdot S_{T}\left(1-\lambda^{(\tau)} \cdot S_{T}\right) \mid \mathscr{F}_{\sigma}\right) \\
& +\varepsilon^{2} E\left(\left(\lambda^{(\tau)} \cdot S_{T}\right)^{2} \mid \mathscr{F}_{\sigma}\right)
\end{aligned}
$$

for any $\varepsilon \in \mathbb{R}$. By statement 4 this implies $E\left(\lambda^{(\tau)} \cdot S_{T}\left(1-\lambda^{(\tau)} \cdot S_{T}\right) \mid \mathscr{F}_{\sigma}\right)=0$. Together, the assertion follows.

Lemma 3.2. 1. There exists a unique semimartingale $L$ with $L_{T}=1$ such that the process $M^{(\tau)}-\left(M^{(\tau)}\right)^{\tau}$ is a martingale for any stopping time $\tau$, where

$$
\begin{equation*}
M^{(\tau)}:=\left(1-\lambda^{(\tau)} \cdot S\right) L \tag{3.6}
\end{equation*}
$$

2. The process $1_{\rrbracket \tau, T \rrbracket} \bullet\left(S M^{(\tau)}\right)$ is a martingale for any stopping time $\tau$. (In the slightly more general setup of Remark 2.8, the upper bound $T$ is to be replaced by $U_{n}$ for arbitrary $n$.)
3. The process $\left(\left(v+\vartheta \cdot S_{s}\right) M_{s}^{(t)}\right)_{s \in[t, T]}$ is a martingale for any $(v, \vartheta) \in$ $\overline{L^{2}\left(\mathscr{F}_{0}\right) \times \Theta}$ and any $t \in[0, T]$.

Proof. 1. Our reasoning relies heavily on the proofs of Lemma 3.4 and Theorem 1.3 in [16]. For any stopping time $\sigma$ we introduce the process

$$
{ }^{\sigma} M_{t}:=\frac{E\left(1-\lambda^{(\sigma)} \cdot S_{T} \mid \mathscr{F}_{t}\right)}{E\left(1-\lambda^{(\sigma)} \cdot S_{T} \mid \mathscr{F}_{\sigma \wedge t}\right)}
$$

Define stopping times $\left(\tau_{n}\right)_{n \in \mathbb{N}}$ recursively by $\tau_{0}:=0$ and

$$
\tau_{n+1}:=\inf \left\{t \geq \tau_{n}:\left|\frac{1-\lambda^{\left(\tau_{n}\right)} \cdot S_{t}}{E\left(1-\lambda^{\left(\tau_{n}\right)} \cdot S_{T} \mid \mathscr{F}_{\tau_{n}}\right)}\right| \leq \frac{1}{2}\right\} \wedge T .
$$

Then

$$
\left.\left.\right|^{\tau_{n}} M_{\tau_{n+1}}\left|=\frac{\left|1-\lambda^{\left(\tau_{n}\right)} \cdot S_{\tau_{n+1}}\right|}{E\left(1-\lambda^{\left(\tau_{n}\right)} \cdot S_{T} \mid \mathscr{F}_{\tau_{n}}\right)}\right| E\left(1-\lambda^{\left(\tau_{n+1}\right)} \cdot S_{T} \mid \mathscr{F}_{\tau_{n+1}}\right) \right\rvert\, \leq \frac{1}{2}
$$

on $\left\{\tau_{n+1}<T\right\}$ by Lemma 3.1. Using Lemma 3.1(2) one easily verifies that

$$
{ }^{\tau_{n}} M_{t}={ }^{\tau_{n}} M_{\tau_{n+1}}{ }^{\tau_{n+1}} M_{t}
$$

for $t \geq \tau_{n+1}$. Consequently, $\lim _{m \rightarrow \infty}{ }^{\tau_{n}} M_{\tau_{m}}=0$ on $D:=\left\{\tau_{n}<T\right.$ for all $\left.n \in \mathbb{N}\right\}$. Letting

$$
\widetilde{M}_{t}^{\left(\tau_{n}\right)}:=E\left(1-\lambda^{\left(\tau_{n}\right)} \cdot S_{T} \mid \mathscr{F}_{t}\right)
$$

we have

$$
\begin{aligned}
1 & =\lim _{m \rightarrow \infty} \frac{E\left(\widetilde{M}_{\tau_{m}}^{\left(\tau_{n}\right)} \mid \mathscr{F}_{\tau_{n}}\right)}{\widetilde{M}_{\tau_{n}}^{\left(\tau_{n}\right)}} \\
& =E\left(\left.\frac{\lim _{m \rightarrow \infty} \widetilde{M}_{\tau_{m}}^{\left(\tau_{n}\right)}}{\widetilde{M}_{\tau_{n}}^{\left(\tau_{n}\right)}} \right\rvert\, \mathscr{F}_{\tau_{n}}\right) \\
& =E\left(\lim _{m \rightarrow \infty}^{\tau_{n}} M_{\tau_{m}} \mid \mathscr{F}_{\tau_{n}}\right) \\
& =E\left({ }^{\tau_{n}} M_{T} 1_{D^{c}} \mid \mathscr{F}_{\tau_{n}}\right) \\
& \leq \sqrt{E\left(\left(\tau_{n} M_{T}\right)^{2} \mid \mathscr{F}_{\tau_{n}}\right)} \sqrt{E\left(1_{D^{c}} \mid \mathscr{F}_{\tau_{n}}\right)} .
\end{aligned}
$$

Since the last term converges to 0 on $D$, it follows that

$$
\begin{equation*}
\lim _{n \rightarrow \infty} E\left(\left({ }^{\tau_{n}} M_{T}\right)^{2} \mid \mathscr{F}_{\tau_{n}}\right)=\infty \quad \text { on } D \tag{3.7}
\end{equation*}
$$

Denote by $Z$ the density process of the measure $Q$ from Assumption 2.1. By Corollary 2.5 we have

$$
\begin{equation*}
E\left(\left.\left(1-\lambda^{\left(\tau_{n}\right)} \cdot S_{T}\right) \frac{Z_{T}}{Z_{\tau_{n}}} \right\rvert\, \mathscr{F}_{\tau_{n}}\right)=1 \tag{3.8}
\end{equation*}
$$

Observe that

$$
\begin{aligned}
E\left(\left.\left(\frac{Z_{T}}{Z_{\tau_{n}}}\right)^{2} \right\rvert\, \mathscr{F}_{\tau_{n}}\right)= & E\left(\left({ }^{\tau_{n}} M_{T}\right)^{2} \mid \mathscr{F}_{\tau_{n}}\right)+2 E\left(\left.{ }^{\tau_{n}} M_{T}\left(\frac{Z_{T}}{Z_{\tau_{n}}}-{ }^{\tau_{n}} M_{T}\right) \right\rvert\, \mathscr{F}_{\tau_{n}}\right) \\
& +E\left(\left.\left(\frac{Z_{T}}{Z_{\tau_{n}}}-\tau_{n} M_{T}\right)^{2} \right\rvert\, \mathscr{F}_{\tau_{n}}\right)
\end{aligned}
$$

for any $n \geq 1$. Due to (3.8) and Lemma 3.1(5) the second term on the right-hand side vanishes. It follows that

$$
\begin{equation*}
E\left(\left(\left(^{\tau_{n}} M_{T}\right)^{2} \mid \mathscr{F}_{\tau_{n}}\right) \leq E\left(\left.\left(\frac{Z_{T}}{Z_{\tau_{n}}}\right)^{2} \right\rvert\, \mathscr{F}_{\tau_{n}}\right)\right. \tag{3.9}
\end{equation*}
$$

Together we have $P(D)=0$ : Indeed, otherwise (3.7) yields

$$
P\left(\sup _{t \in[0, T]} \frac{E\left(Z_{T}^{2} \mid \mathscr{F}_{t}\right)}{Z_{t}^{2}}<E\left(\left({ }^{\tau_{n}} M_{T}\right)^{2} \mid \mathscr{F}_{\tau_{n}}\right)\right)>0
$$

for large $n$. Consequently,

$$
\left\{\frac{E\left(Z_{T}^{2} \mid \mathscr{F}_{\tau_{n}}\right)}{Z_{\tau_{n}}^{2}}<E\left(\left({ }^{\tau_{n}} M_{T}\right)^{2} \mid \mathscr{F}_{\tau_{n}}\right)\right\} \in \mathscr{F}_{\tau_{n}}
$$

has positive probability as well in contradiction to (3.9).
Now define the semimartingale $L$ by

$$
L_{t}:=\frac{E\left(1-\lambda^{\left(\tau_{n}\right)} \cdot S_{T} \mid \mathscr{F}_{t}\right)}{1-\lambda^{\left(\tau_{n}\right)} \cdot S_{t}} \quad \text { for } \tau_{n} \leq t<\tau_{n+1}
$$

The claimed martingale property follows from Lemma 3.1(2).
Uniqueness of $L$ follows from

$$
E\left(1-\lambda^{(t)} \cdot S_{T} \mid \mathscr{F}_{t}\right)-L_{t}=E\left(M_{T}^{(t)} \mid \mathscr{F}_{t}\right)-M_{t}^{(t)}=0 .
$$

2. It suffices to verify that $E\left(1_{\rrbracket \tau, T \rrbracket} \cdot\left(S M^{(\tau)}\right)_{\sigma}\right)=0$ for any stopping time $\sigma$. By substituting $\sigma \vee \tau$ for $\sigma$, we may assume $\sigma \geq \tau$ w.l.o.g. Since $1_{\rrbracket \tau, T \rrbracket} \cdot M^{(\tau)}$ is a square-integrable martingale, we have $E\left(S_{\sigma}\left(M_{T}^{(\tau)}-M_{\sigma}^{(\tau)}\right)\right)=0$ and similarly $E\left(S_{\tau}\left(M_{T}^{(\tau)}-M_{\tau}^{(\tau)}\right)\right)=0$. Consequently,

$$
E\left(S_{\sigma} M_{\sigma}^{(\tau)}-S_{\tau} M_{\tau}^{(\tau)}\right)=E\left(\left(S_{\sigma}-S_{\tau}\right) M_{T}^{(\tau)}\right)=E\left(\left(\psi \cdot S_{T}\right) M_{T}^{(\tau)}\right)
$$

for $\psi:=1_{\rrbracket \tau, \sigma \rrbracket}$. The optimality of $\lambda^{(\tau)}$ implies that

$$
\begin{aligned}
0 & \leq E\left(\left(1-\left(\lambda^{(\tau)}+\varepsilon \psi\right) \cdot S_{T}\right)^{2}\right)-E\left(\left(1-\lambda^{(\tau)} \cdot S_{T}\right)^{2}\right) \\
& =2 \varepsilon E\left(\left(\psi \cdot S_{T}\right) M_{T}^{(\tau)}\right)+\varepsilon^{2} E\left(\left(\psi \cdot S_{T}\right)^{2}\right)
\end{aligned}
$$

for any $\varepsilon \in \mathbb{R}$ and hence $E\left(\left(\psi \cdot S_{T}\right) M_{T}^{(\tau)}\right)=0$.
3. By statement 2 we have that $1_{\rrbracket t, T \rrbracket} \cdot\left(S M^{(t)}\right)$ and hence $\left(S_{S} M_{s}^{(t)}\right)_{s \in[t, T]}$ is a martingale. Consequently, the signed measure with density process $\left(M^{(t)} /\right.$ $\left.E\left(L_{t}\right)\right)_{s \in[t, T]}$ is a $\mathrm{S} \sigma \mathrm{MM}$ in the sense of Definition 2.3 if the time set $[0, T]$ is replaced with $[t, T]$. By Lemma 2.6 (also adapted to $[t, T]$ instead of $[0, T]$ as time set), the assertion follows.

DEFINITION 3.3. We call the process $L$ from Lemma 3.2 opportunity process.
The terminology is inspired by the fact that $L$ is linked to optimal investment opportunities. Indeed, the following corollary states that $L$ represents both first and second moments of efficient strategies in the sense of (3.1).

Corollary 3.4. For any $t \in[0, T]$ we have

$$
\begin{align*}
L_{t} & =E\left(1-\lambda^{(t)} \cdot S_{T} \mid \mathscr{F}_{t}\right) \\
& =E\left(\left(1-\lambda^{(t)} \cdot S_{T}\right)^{2} \mid \mathscr{F}_{t}\right)  \tag{3.10}\\
& =\inf \left\{E\left(\left(1-\vartheta \cdot S_{T}\right)^{2} \mid \mathscr{F}_{t}\right): \vartheta \in \bar{\Theta} \text { with } \vartheta 1_{\llbracket 0, t \rrbracket}=0\right\} .
\end{align*}
$$

In particular, $L$ is a submartingale.
Proof. This follows from Lemmas 3.2 and 3.1.
These equations can be interpreted in terms of dynamic Sharpe ratios (cf. also [31], (5.16)):

DEFINITION 3.5. For $t \in[0, T]$ we call

$$
\begin{equation*}
\varrho_{t}:=\sup \left\{\frac{E\left(\vartheta \cdot S_{T} \mid \mathscr{F}_{t}\right)}{\sqrt{\operatorname{Var}\left(\vartheta \cdot S_{T} \mid \mathscr{F}_{t}\right)}}: \vartheta \in \bar{\Theta} \text { with } \vartheta 1_{\llbracket 0, t \rrbracket}=0\right\} \tag{3.11}
\end{equation*}
$$

maximal Sharpe ratio on $(t, T]$, where we set $\operatorname{Var}\left(X \mid \mathscr{F}_{t}\right):=E\left(X^{2} \mid \mathscr{F}_{t}\right)-$ $\left(E\left(X \mid \mathscr{F}_{t}\right)\right)^{2}$ and $\frac{0}{0}:=0$.

PROPOSITION 3.6. The relation between opportunity process $L$ and maximal Sharpe ratio $\varrho$ is given by

$$
\varrho=\sqrt{\frac{1}{L}-1}
$$

and

$$
L=\frac{1}{1+\varrho^{2}}
$$

respectively.
Proof. On the set

$$
D:=\left\{\omega \in \Omega: E\left(\vartheta \cdot S_{T} \mid \mathscr{F}_{t}\right)(\omega)=0 \text { for all } \vartheta \in \bar{\Theta} \text { with } \vartheta 1_{\llbracket 0, t \rrbracket}=0\right\}
$$

we have $\varrho_{t}=0$. Moreover, the infimum in (3.10) is attained in $\lambda^{(t)}=0$, which implies that $L_{t}=1$ on $D$.

For $\omega \in D^{C}$ there exists some $\vartheta \in \bar{\Theta}$ with $\vartheta 1_{\llbracket 0, t \rrbracket=0}$ and $E\left(\vartheta \cdot S_{T} \mid \mathscr{F}_{t}\right)(\omega)>0$. For sufficiently small $\varepsilon>0$ we have that $E\left(\left(1-\varepsilon \vartheta \cdot S_{T}\right)^{2} \mid \mathscr{F}_{t}\right)(\omega)<1$, which implies that $L_{t}<1$ on $D^{C}$ (cf. Corollary 3.4). By scaling invariance it suffices to consider $\vartheta$ with $E\left(\vartheta \cdot S_{T} \mid \mathscr{F}_{t}\right)=1-L_{t}$ in the supremum of (3.11). For these $\vartheta$ we have

$$
\frac{E\left(\vartheta \cdot S_{T} \mid \mathscr{F}_{t}\right)}{\sqrt{\operatorname{Var}\left(\vartheta \cdot S_{T} \mid \mathscr{F}_{t}\right)}}=\frac{1-L_{t}}{\sqrt{E\left(\left(1-\vartheta \cdot S_{T}\right)^{2} \mid \mathscr{F}_{t}\right)-L_{t}^{2}}}
$$

which implies that the supremum is attained in $\vartheta=\lambda^{(t)}$. The assertion follows now from Corollary 3.4.
3.2. Adjustment process. The optimal number of shares $\lambda^{(\tau)}$ in (3.1) depends on $\tau$. However, the optimal number of shares per unit of wealth does not. It is denoted by $\tilde{a}$ in the following lemma.

Lemma 3.7. We use the notation from Lemma 3.1. There exists some $\tilde{a} \in$ $L(S)$ such that

$$
1-\lambda^{(\tau)} \cdot S=\mathscr{E}\left(\left(-\tilde{a} 1_{\rrbracket \tau, T \rrbracket}\right) \cdot S\right)=1-\left(\tilde{a} 1_{\rrbracket \tau, T \rrbracket} \mathscr{E}\left(\left(-\tilde{a} 1_{\rrbracket \tau, T \rrbracket}\right) \cdot S\right)_{-}\right) \cdot S
$$

for any stopping time $\tau$. Consequently, we may assume

$$
\begin{equation*}
\lambda^{(\tau)}=\tilde{a} 1_{\rrbracket \tau, T \rrbracket} \mathscr{E}\left(\left(-\tilde{a} 1_{\rrbracket \tau, T \rrbracket}\right) \cdot S\right)_{-} . \tag{3.12}
\end{equation*}
$$

Proof. Let

$$
\tilde{a}:=\sum_{n=0}^{\infty} \frac{\lambda^{\left(\tau_{n}\right)}}{1-\lambda^{\left(\tau_{n}\right)} \cdot S_{-}} 1_{\rrbracket \tau_{n}, \tau_{n+1} \rrbracket},
$$

where $\left(\tau_{n}\right)_{n \in \mathbb{N}}$ denotes the sequence of stopping times from the proof of Lemma 3.2. On $\llbracket 0, \tau_{n+1} \rrbracket$ we have

$$
1-\lambda^{\left(\tau_{n}\right)} \cdot S=1-\left(\left(1-\lambda^{\left(\tau_{n}\right)} \cdot S_{-}\right) \tilde{a} 1_{\rrbracket \tau_{n}, T \rrbracket}\right) \cdot S
$$

and hence

$$
\begin{equation*}
1-\lambda^{\left(\tau_{n}\right)} \cdot S=\mathscr{E}\left(\left(-\tilde{a} 1_{\rrbracket \tau_{n}, T \rrbracket}\right) \cdot S\right) \tag{3.13}
\end{equation*}
$$

From

$$
1-\lambda^{\left(\tau_{n}\right)} \cdot S_{t}=\left(1-\lambda^{\left(\tau_{n}\right)} \cdot S_{\tau_{n+1}}\right)\left(1-\lambda^{\left(\tau_{n+1}\right)} \cdot S_{t}\right)
$$

and

$$
\mathscr{E}\left(\left(-\tilde{a} 1_{\rrbracket \tau_{n}, T \rrbracket}\right) \cdot S\right)_{t}=\mathscr{E}\left(\left(-\tilde{a} 1_{\rrbracket \tau_{n}, T \rrbracket}\right) \cdot S\right)_{\tau_{n+1}} \mathscr{E}\left(\left(-\tilde{a} 1_{\rrbracket \tau_{n+1}, T \rrbracket}\right) \cdot S\right)_{t}
$$

for $t \in \rrbracket \tau_{n+1}, \tau_{n+2} \rrbracket$ it follows recursively that (3.13) holds on [0, T]. Now let $\tau$ be arbitrary. On $\left\{\tau_{n} \leq \tau<\tau_{n+1}\right\}$ we have

$$
1-\lambda^{(\tau)} \cdot S=\frac{1-\lambda^{\left(\tau_{n}\right)} \cdot S}{1-\lambda^{\left(\tau_{n}\right)} \cdot S_{\tau}}=\frac{\mathscr{E}\left(\left(-\tilde{a} 1_{\rrbracket \tau_{n}, T \rrbracket}\right) \cdot S\right)}{\mathscr{E}\left(\left(-\tilde{a} 1_{\rrbracket \tau_{n}, T \rrbracket}\right) \cdot S\right)_{\tau}}=\mathscr{E}\left(\left(-\tilde{a} 1_{\rrbracket \tau, T \rrbracket}\right) \cdot S\right)
$$

as claimed.
DEFINITION 3.8. The (not necessarily unique) process $\tilde{a}$ from Lemma 3.7 is called adjustment process. Moreover, we call

$$
\hat{a}:=\left(1+\Delta A^{K}\right) \tilde{a}
$$

extended adjustment process.

The name adjustment process is taken from [44]:
Corollary 3.9. $E\left(\vartheta \cdot S_{T} \mathscr{E}(-\tilde{a} \cdot S)_{T}\right)=0$ for any $\vartheta \in \bar{\Theta}$, i.e., $\tilde{a}$ is an adjustment process in the sense of [44], Section 3 with $\bar{\Theta}$ substituted for $\Theta$.

Proof. This follows from Lemma 3.2(3).
Lemma 3.10. $\quad L, L_{-}$are $(0,1]$-valued.
Proof. Lemma 3.1(5) implies that $L_{t}=E\left(1-\lambda^{(t)} \cdot S_{T} \mid \mathscr{F}_{t}\right) \in(0,1]$ almost surely for fixed $t$, which yields by right-continuity that $L$ is [ 0,1$]$-valued outside some evanescent set.

Let $\tau:=\inf \left\{t \in[0, T]: L_{t}=0\right\} \wedge T$. Again by Lemma 3.1(5), we have

$$
0=L_{\tau \wedge T}=E\left(\left(1-\lambda^{(\tau \wedge T)} \cdot S_{T}\right)^{2} \mid \mathscr{F}_{\tau \wedge T}\right) \in(0,1]
$$

on $\left\{L_{t}=0\right.$ for some $\left.t \in[0, T]\right\}$, which implies that

$$
\begin{equation*}
P\left(L_{t}=0 \text { for some } t \in[0, T]\right)=0 . \tag{3.14}
\end{equation*}
$$

Finally let $\tau:=\inf \left\{t \in[0, T]: L_{t-}=0\right\} \wedge T$. Define an increasing sequence of stopping times $\left(\tau_{n}\right)_{n \in \mathbb{N}}$ via $\tau_{n}:=\inf \left\{t \in[0, T]: L_{t} \leq \frac{1}{n}\right\} \wedge T$. By (3.14) we have $\tau_{n} \uparrow \uparrow \tau$ on $\left\{L_{\tau-}=0\right\}$. Lemma 3.1(5) implies

$$
E\left(\left(1-\lambda^{\left(\tau_{n}\right)} \cdot S_{T}\right)^{2} 1_{\left\{\tau_{n}<T\right\}}\right)=E\left(L_{\tau_{n}} 1_{\left\{\tau_{n}<T\right\}}\right)
$$

By [39], Theorem V. 13 we have that

$$
1-\lambda^{\left(\tau_{n}\right)} \cdot S_{T}=\mathscr{E}\left(\left(-\tilde{a} 1_{\rrbracket \tau_{n}, T \rrbracket}\right) \cdot S\right)_{T} \rightarrow \mathscr{E}\left(\left(-\tilde{a} 1_{\llbracket \tau, T \rrbracket \cap\left\{L_{\tau-}=0\right\}}\right) \cdot S\right)_{T}
$$

in probability for $n \rightarrow \infty$. In view of Fatou's lemma and dominated convergence, we obtain

$$
0 \leq E\left(\left(\mathscr{E}\left(\left(-\tilde{a} 1_{\llbracket \tau, T \rrbracket}\right) \cdot S\right)_{T}\right)^{2} 1_{\left\{L_{\tau-}=0\right\}}\right) \leq E\left(L_{\tau-} 1_{\left\{L_{\tau-}=0\right\}}\right)=0
$$

Suppose that $\left\{L_{-}=0\right\}$ is not evanescent. Then there is some $n$ such that $P(D)>0$ for

$$
\begin{aligned}
D & :=\left\{L_{\tau-}=0 \text { and } 1-\lambda^{\left(\tau_{n}\right)} \cdot S_{\tau-}>0\right\} \\
& \supset\left\{L_{\tau-}=0 \text { and } \Delta(-\tilde{a} \cdot S)>-1 \text { on } \rrbracket \tau_{n}, \tau \llbracket\right\} .
\end{aligned}
$$

On $D \in \mathscr{F}_{\tau-}$ we have

$$
\frac{1-\lambda^{\left(\tau_{n}\right)} \cdot S_{T}}{1-\lambda^{\left(\tau_{n}\right)} \cdot S_{\tau-}}=\frac{\mathscr{E}\left(\left(-\tilde{a} 1_{\rrbracket \tau_{n}, T \rrbracket}\right) \cdot S\right)_{T}}{\mathscr{E}\left(\left(-\tilde{a} 1_{\rrbracket \tau_{n}, T \rrbracket}\right) \cdot S\right)_{\tau-}}=\mathscr{E}\left(\left(-\tilde{a} 1_{\llbracket \tau, T \rrbracket}\right) \cdot S\right)_{T}=0
$$

Consequently, the process $\lambda^{\left(\tau_{n}\right)} \cdot S$ cannot be a martingale under the $\sigma$-martingale measure from Assumption 2.1, which yields a contradiction to Corollary 2.5.

Since $L_{-}$does not vanish, the stochastic logarithm of $L$ is well defined:

Definition 3.11. We call

$$
K:=\mathscr{L}(L):=\frac{1}{L_{-}} \cdot L
$$

modified mean-variance tradeoff (MMVT) process.
The modified mean-variance tradeoff process is related to the mean-variance tradeoff (MVT) process of [42] (cf. Section 3.6).
3.3. Variance-optimal signed martingale measure. With the help of the modified mean-variance tradeoff process $K$ and the adjustment process $\tilde{a}$ we can define a signed measure $Q^{\star}$ which plays an important role in the context of quadratic hedging. This variance-optimal signed martingale measure appears more or less explicitly in many papers on the subject.

Definition 3.12. We call

$$
\begin{equation*}
N:=K-\tilde{a} \cdot S-[\tilde{a} \cdot S, K] \tag{3.15}
\end{equation*}
$$

variance-optimal logarithm process and the signed measure $Q^{\star}$ defined via

$$
\begin{equation*}
\frac{d Q^{\star}}{d P}:=\frac{L_{0}}{E\left(L_{0}\right)} \mathscr{E}(N)_{T}=\frac{\mathscr{E}(-\tilde{a} \cdot S)_{T}}{E\left(L_{0}\right)}=\frac{1-\lambda^{(0)} \cdot S_{T}}{E\left(1-\lambda^{(0)} \cdot S_{T}\right)} \tag{3.16}
\end{equation*}
$$

variance-optimal signed martingale measure (variance-optimal $S \sigma M M$ ).
The following result explains the terminology.
Proposition 3.13. 1. $Q^{\star}$ is a $S \sigma M M$ (cf. Definition 2.3) with density process

$$
Z^{Q^{\star}}:=\frac{L_{0}}{E\left(L_{0}\right)} \mathscr{E}(N)=\frac{L \mathscr{E}(-\tilde{a} \cdot S)}{E\left(L_{0}\right)}
$$

2. $Q^{\star}$ minimizes $Q \mapsto E\left(\left(\frac{d Q}{d P}\right)^{2}\right)$ over all $S \sigma M M$ 's $Q$. Hence it is the varianceoptimal signed $\bar{\Theta}$-martingale measure in the sense of [44], Section 1 , with $\Theta$ replaced by $\bar{\Theta}$ in the definition.

Proof. 1. Note that $L_{0} \mathscr{E}(N)=M^{(0)}$ is a martingale by Lemma 3.2. Lemma 3.2(3) implies that $Q^{\star}$ is a $\mathrm{S} \sigma \mathrm{MM}$.
2. For any other $\operatorname{S} \sigma$ MM $Q$ with $\frac{d Q}{d P} \in L^{2}(P)$ we have

$$
\begin{aligned}
& E\left(\left(\frac{d Q}{d P}\right)^{2}\right)-E\left(\left(\frac{d Q^{\star}}{d P}\right)^{2}\right) \\
& \quad \geq 2 E\left(\left(\frac{d Q}{d P}-\frac{d Q^{\star}}{d P}\right) \frac{d Q^{\star}}{d P}\right) \\
& \quad=2 E\left(\frac{d Q}{d P} \frac{1-\lambda^{(0)} \cdot S_{T}}{E\left(L_{0}\right)}\right)-2 E\left(\frac{d Q^{\star}}{d P} \frac{1-\lambda^{(0)} \cdot S_{T}}{E\left(L_{0}\right)}\right)=0
\end{aligned}
$$

by Corollary 2.5 .
If $Q \sim P$ is a probability measure with density process $Z=\mathscr{E}(M)$, then the density $\frac{d Q}{d P}$, the density process $Z$, and its stochastic logarithm $M$ uniquely determine one another. This is not true for the variance-optimal $\mathrm{S} \sigma \mathrm{MM} Q^{\star}$ because $\mathscr{E}(N)$ may vanish and hence $N$ cannot be fully recovered from $\mathscr{E}(N)$ or $\frac{d Q^{\star}}{d P}$. Therefore the following result does not follow immediately from the fact that $Q^{\star}$ is a $\mathrm{S} \sigma \mathrm{MM}$ whose density process is a multiple of $\mathscr{E}(N)$.

Lemma 3.14. The variance-optimal logarithm process $N$ and also $S+[S, N]$ are $\sigma$-martingales. Consequently, $S \mathscr{E}(N)$ is a $\sigma$-martingale as well.

Proof. Denote by $\left(\tau_{n}\right)_{n \in \mathbb{N}}$ the sequence of stopping times from the proof of Lemma 3.2. Since $\bigcup_{n \in \mathbb{N}} \rrbracket \tau_{n}, \tau_{n+1} \rrbracket=\Omega \times(0, T]$, it suffices to show that $1_{\rrbracket \tau_{n}, \tau_{n+1} \rrbracket} \cdot N$ and $1_{\rrbracket \tau_{n}, \tau_{n+1} \rrbracket} \bullet(S+[N, S])$ are $\sigma$-martingales for any $n \in \mathbb{N}$. Since

$$
\begin{aligned}
\mathscr{E}\left(N-N^{\tau_{n}}\right) & =\mathscr{E}\left(1_{\rrbracket \tau_{n}, T \rrbracket} \cdot(K-\tilde{a} \cdot S-[\tilde{a} \cdot S, K])\right) \\
& =\mathscr{E}\left(1_{\rrbracket \tau_{n}, T \rrbracket} \cdot K\right) \mathscr{E}\left(\left(-\tilde{a} 1_{\rrbracket \tau_{n}, T \rrbracket}\right) \cdot S\right) \\
& =\frac{L\left(1-\lambda^{\left(\tau_{n}\right)} \cdot S\right)}{L^{\tau_{n}}} \\
& =1+\frac{1_{\rrbracket \tau_{n}, T \rrbracket}}{L_{\tau_{n}}} \cdot M^{\left(\tau_{n}\right)}
\end{aligned}
$$

is a $\sigma$-martingale, we have that

$$
1_{\rrbracket \tau_{n}, \tau_{n+1} \rrbracket} \cdot N=\frac{1_{\rrbracket \tau_{n}, \tau_{n+1} \rrbracket}^{\mathscr{E}\left(N-N^{\tau_{n}}\right)_{-}} \cdot \mathscr{E}\left(N-N^{\tau_{n}}\right), ~\left(\frac{1}{}\right)}{}
$$

is a $\sigma$-martingale as well. Similarly,

$$
\begin{aligned}
1_{\rrbracket \tau_{n}, T \rrbracket} \cdot\left(\mathscr{E}\left(N-N^{\tau_{n}}\right) S\right) & =1_{\rrbracket \tau_{n}, T \rrbracket} \cdot \frac{L\left(1-\lambda^{\left(\tau_{n}\right)} \cdot S\right) S}{L^{\tau_{n}}} \\
& =\frac{1_{\rrbracket \tau_{n}, T \rrbracket}}{L_{\tau_{n}}} \cdot\left(1_{\rrbracket \tau_{n}, T \rrbracket} \cdot\left(M^{\left(\tau_{n}\right)} S\right)\right)
\end{aligned}
$$

is a $\sigma$-martingale by Lemma 3.2(2). Integration by parts yields

$$
\begin{array}{r}
1_{\rrbracket \tau_{n}, T \rrbracket} \cdot\left(\mathscr{E}\left(N-N^{\tau_{n}}\right) S\right)-S_{-} \cdot \mathscr{E}\left(N-N^{\tau_{n}}\right) \\
\quad=\left(\mathscr{E}\left(N-N^{\tau_{n}}\right)-1_{\rrbracket \tau_{n}, T \rrbracket}\right) \cdot(S+[N, S \rrbracket)
\end{array}
$$

which implies that

$$
\begin{aligned}
& 1_{\rrbracket} \tau_{n}, \tau_{n+1} \rrbracket \cdot(S+[N, S]) \\
& =\frac{1 \rrbracket \tau_{n}, \tau_{n+1} \rrbracket}{\mathscr{E}\left(N-N^{\tau_{n}}\right)_{-}} \cdot\left(\left(\mathscr{E}\left(N-N^{\tau_{n}}\right)_{-} 1_{\rrbracket \tau_{n}, T \rrbracket}\right) \cdot(S+[N, S])\right)
\end{aligned}
$$

is a $\sigma$-martingale as well. Finally,

$$
S \mathscr{E}(N)=S_{0}+\mathscr{E}(N)_{-} \cdot\left(S_{-} \cdot N+S+[S, N]\right)
$$

yields the last assertion.
3.4. Opportunity-neutral measure. In this section we define a measure $P^{\star}$ in terms of its density process

$$
Z^{P^{\star}}:=\frac{L}{E\left(L_{0}\right) \mathscr{E}\left(A^{K}\right)}
$$

For $Z^{P^{\star}}$ to be truly a density process, we need the following
Lemma 3.15. The process $Z^{P^{\star}}$ is a bounded positive martingale and satisfies

$$
Z^{P^{\star}}=\frac{L_{0}}{E\left(L_{0}\right)} \mathscr{E}\left(\frac{1}{1+\Delta A^{K}} \cdot M^{K}\right)
$$

Proof. Since $L$ is a submartingale by Corollary 3.4, we have $b^{L} \geq 0$ and hence $b^{K}=\frac{1}{L_{-}} b^{L} \geq 0$ outside some $P \otimes A$-null set. This implies that $A^{K}=b^{K} \cdot A$ and hence also $\mathscr{E}\left(A^{K}\right)$ are increasing processes. Thus we have $0<Z^{P^{\star}} \leq \frac{1}{E\left(L_{0}\right)}$. The equality of the two expressions for $Z^{P^{\star}}$ follows from Yor's formula. From the second representation we conclude that $Z^{P^{\star}}$ is a local martingale and hence a martingale because it is bounded.

DEFINITION 3.16. We call the probability measure $P^{\star} \sim P$ with density process $Z^{P^{\star}}$ opportunity-neutral probability measure.

The opportunity-neutral probability measure is typically not a martingale measure. In some instances it actually equals $P$ (cf. Section 3.6). For later use we determine the $P^{\star}$-characteristics of $S$.

Lemma 3.17. The components of $S$ are locally $P^{\star}$-square integrable semimartingales. Moreover,

$$
\begin{align*}
b^{S \star} & =\frac{\bar{b}}{1+\Delta A^{K}},  \tag{3.17}\\
\tilde{c}^{S \star} & =\frac{\bar{c}}{1+\Delta A^{K}}, \tag{3.18}
\end{align*}
$$

and

$$
\begin{align*}
\hat{c}^{S \star}\left(\hat{c}^{S \star}\right)^{-1} b^{S \star} & =b^{S \star},  \tag{3.20}\\
\tilde{c}^{S \star}\left(\tilde{c}^{S \star}\right)^{-1} b^{S \star} & =b^{S \star},  \tag{3.21}\\
\bar{c} \bar{c}^{-1} \bar{b} & =\bar{b} . \tag{3.22}
\end{align*}
$$

$P \otimes A$-almost everywhere, where

$$
\begin{align*}
\bar{b} & :=b^{S}+c^{S L} \frac{1}{L_{-}}+\int x \frac{y}{L_{-}} F^{S, L}(d(x, y))  \tag{3.23}\\
& =b^{S}+c^{S K}+\int x y F^{S, K}(d(x, y)) \tag{3.24}
\end{align*}
$$

and

$$
\begin{align*}
\bar{c} & :=c^{S}+\int x x^{\top}\left(1+\frac{y}{L_{-}}\right) F^{S, L}(d(x, y))  \tag{3.25}\\
& =c^{S}+\int x x^{\top}(1+y) F^{S, K}(d(x, y)) \tag{3.26}
\end{align*}
$$

Proof. The components of $S$ are locally $P^{\star}$-square-integrable semimartingales because $\frac{d P^{\star}}{d P}=Z_{T}^{P^{\star}}$ is bounded (cf. Lemma A.2). Let

$$
M:=\frac{1}{Z_{-}^{P^{\star}}} \cdot Z^{P^{\star}}=\frac{1}{1+\Delta A^{K}} \cdot M^{K}
$$

Observe that

$$
K^{c}-M^{c}=K^{c}-\frac{1}{1+\Delta A^{K}} \cdot K^{c}=\frac{\Delta A^{K}}{1+\Delta A^{K}} \cdot K^{c}
$$

Since

$$
\begin{aligned}
\left\langle\frac{\Delta A^{K}}{1+\Delta A^{K}} \cdot K^{c}, \frac{\Delta A^{K}}{1+\Delta A^{K}} \cdot K^{c}\right\rangle_{T} & =\left(\frac{\Delta A^{K}}{1+\Delta A^{K}}\right)^{2} \cdot\left\langle K^{c}, K^{c}\right\rangle_{T} \\
& =\sum_{t \leq T}\left(\frac{\Delta A_{t}^{K}}{1+\Delta A_{t}^{K}}\right)^{2} \Delta\left\langle K^{c}, K^{c}\right\rangle_{t} \\
& =0
\end{aligned}
$$

by continuity of $K^{c}$, we have $\frac{\Delta A^{K}}{1+\Delta A^{K}} \cdot K^{c}=0$ and hence $M^{c}=K^{c}$. Moreover, $M$ is a local martingale with $\Delta M=\frac{1}{1+\Delta A^{K}} \Delta K-\frac{\Delta A^{K}}{1+\Delta A^{K}}$. Together, it follows that $b^{S, M}=\left(b^{S}, 0\right)^{\top}, c^{S, M}=c^{S, K}$,

$$
F^{S, M}(G)=\int 1_{G}\left(x, \frac{y-\Delta A^{K}}{1+\Delta A^{K}}\right) F^{S, K}(d(x, y))
$$

for $G \in \mathscr{B}^{d+1}$ with $G \cap\left(\{0\}^{d} \times \mathbb{R}\right)=\varnothing$. By the Girsanov theorem as in Lemma A.9, $P^{\star}$-characteristics $\left(b^{S \star}, c^{S \star}, F^{S \star}, A\right)$ of $S$ are given by

$$
\begin{aligned}
b^{S \star} & =b^{S}+c^{S M}+\int x y F^{S, M}(d(x, y)) \\
& =b^{S}+c^{S K}+\int x \frac{y-\Delta A^{K}}{1+\Delta A^{K}} F^{S, K}(d(x, y)) \\
& =\frac{1}{1+\Delta A^{K}}\left(b^{S}+c^{S K}+\int x y F^{S, K}(d(x, y))\right),
\end{aligned}
$$

$c^{S \star}=c^{S}$ and

$$
\begin{aligned}
F^{S \star}(G) & =\int 1_{G}(x)(1+y) F^{S, M}(d(x, y)) \\
& =\frac{1}{1+\Delta A^{K}} \int 1_{G}(x)(1+y) F^{S, K}(d(x, y))
\end{aligned}
$$

for $G \in \mathscr{B}^{d}$ with $0 \notin G$. This yields (3.17), (3.18).
Using the same argument as in the proof of [14], Theorem 3.5, it follows that $b_{t}^{S \star} \in \hat{c}_{t}^{S \star} \mathbb{R}^{d}$ and also $b_{t}^{S \star} \in \tilde{c}_{t}^{S \star} \mathbb{R}^{d}(P \otimes A)$-almost everywhere on $\Omega \times[0, T]$. (Due to Assumption 2.1 local boundedness is not needed in our setup.) This implies (3.20), (3.21), and hence also (3.22) outside some $P \otimes A$-null set. Consequently,

$$
\begin{aligned}
(1+ & \left.\left(b^{S \star}\right)^{\top}\left(\hat{c}^{S \star}\right)^{-1} b^{S \star} \Delta A\right)\left(1-\left(b^{S \star}\right)^{\top}\left(\tilde{c}^{S \star}\right)^{-1} b^{S \star} \Delta A\right) \\
& =1+\left(b^{S \star}\right)^{\top}\left(\left(\hat{c}^{S \star}\right)^{-1}-\left(\tilde{c}^{S \star}\right)^{-1}-\left(\hat{c}^{S \star}\right)^{-1} b^{S \star}\left(b^{S \star}\right)^{\top}\left(\tilde{c}^{S \star}\right)^{-1} \Delta A\right) b^{S \star} \Delta A \\
& =1+\left(b^{S \star}\right)^{\top}\left(\hat{c}^{S \star}\right)^{-1}\left(\tilde{c}^{S \star}-\hat{c}^{S \star}-b^{S \star}\left(b^{S \star}\right)^{\top} \Delta A\right)\left(\tilde{c}^{S \star}\right)^{-1} b^{S \star} \Delta A \\
& =1 .
\end{aligned}
$$

REMARK 3.18. An inspection of the proofs of Lemmas 3.15 and 3.17 yields that $L$ need not be the opportunity process for (3.22) to hold. We only used the fact that $L=L_{0} \mathscr{E}(K)$ is a bounded semimartingale with $b^{L} \geq 0$ and $L, L_{-}>0$.
3.5. Characterization of $L$ and $\tilde{a}$. The opportunity process $L$ and the adjustment process $\tilde{a}$ play a crucial role in quadratic hedging. For example, they yield the density processes of the variance-optimal $\mathrm{S} \sigma \mathrm{MM} Q^{\star}$ and the opportunity-neutral measure $P^{\star}$, which in turn lead to formulas for the optimal hedge in Section 4. The characterizations of $L$ and $\tilde{a}$ in this section help to determine these processes in concrete models.

Lemma 3.19. We have

$$
\begin{align*}
b^{L} & =L_{-} \tilde{a}^{\top} \bar{b},  \tag{3.27}\\
\bar{b} & =\bar{c} \tilde{a},  \tag{3.28}\\
b^{K} & =\bar{b}^{\top} \bar{c}^{-1} \bar{b}=\left(b^{S \star}\right)^{\top}\left(\hat{c}^{S \star}\right)^{-1} b^{S \star} \tag{3.29}
\end{align*}
$$

outside some $P \otimes A$-null set, where $\bar{b}, \bar{c}$ are defined in (3.23) and (3.25).
Proof. We denote by $\tau_{n}$ the stopping times in the proof of Lemma 3.2. Fix $n \in \mathbb{N}$. Integration by parts and Lemma 3.2 yield that

$$
\left(\mathscr{E}\left(\left(-\tilde{a} 1_{\rrbracket \tau_{n}, T \rrbracket}\right) \cdot S\right)_{-} 1_{\rrbracket \tau_{n}, T \rrbracket}\right) \cdot\left(L-\left(L_{-} \tilde{a}\right) \cdot S-\tilde{a} \cdot[L, S]\right)=1_{\rrbracket \tau_{n}, T \rrbracket} \cdot M^{\left(\tau_{n}\right)}
$$

is a martingale. Consequently, its compensator

$$
\left(\mathscr{E}\left(\left(-\tilde{a} 1_{\rrbracket \tau_{n}, T \rrbracket}\right) \cdot S\right)_{-}\left(b^{L}-L_{-} \tilde{a}^{\top} b^{S}-\tilde{a}^{\top} \tilde{c}^{S L}\right) 1_{\rrbracket \tau_{n}, T \rrbracket}\right) \cdot A
$$

vanishes. Since $\mathscr{E}\left(\left(-\tilde{a} 1_{\rrbracket \tau_{n}, T \rrbracket}\right) \cdot S\right)_{-} \neq 0$ on $\rrbracket \tau_{n}, \tau_{n+1} \rrbracket$, this implies that

$$
b^{L}-\tilde{a}^{\top} L_{-} b^{S}-\tilde{a}^{\top} \tilde{c}^{S L}=0
$$

$P \otimes A$-almost everywhere on $\rrbracket \tau_{n}, \tau_{n+1} \rrbracket$. This yields (3.27).
Fix $n \in \mathbb{N}$. From Lemma 3.2(2) and integration by parts it follows that

$$
\begin{aligned}
0 \sim & 1_{\rrbracket \tau_{n}, T \rrbracket} \cdot\left(S M^{\left(\tau_{n}\right)}\right) \\
= & 1_{\rrbracket \tau_{n}, T \rrbracket} \cdot\left(S_{-} \cdot M^{\left(\tau_{n}\right)}+M_{-}^{\left(\tau_{n}\right)} \cdot S+\left[S, M^{\left(\tau_{n}\right)}\right]\right) \\
\sim & 1_{\rrbracket \tau_{n}, T \rrbracket} \cdot\left(\left(\mathscr{E}\left(\left(-a 1_{\rrbracket \tau_{n}, T \rrbracket}\right) \cdot S\right)_{-} L_{-}\right) \cdot S+\left[S, \mathscr{E}\left(\left(-a 1_{\rrbracket \tau_{n}, T \rrbracket}\right) \cdot S\right) L\right]\right) \\
= & \left(\mathscr{E}\left(\left(-a 1_{\rrbracket \tau_{n}, T \rrbracket}\right) \cdot S\right)_{-} 1_{\rrbracket \tau_{n}, T \rrbracket}\right) \\
& \cdot\left(L_{-} \cdot S+[S, L]-\tilde{a} \cdot\left(L_{-} \cdot[S, S]-[[L, S], S]\right)\right) \\
\sim & \left(\mathscr{E}\left(\left(-a 1_{\rrbracket \tau_{n}, T \rrbracket}\right) \cdot S\right)_{-} 1_{\rrbracket \tau_{n}, T \rrbracket}\right. \\
& \left.\times\left(L_{-} b^{S}+\tilde{c}^{S L}-\left(L_{-} \tilde{c}^{S}+\int x x^{\top} y F^{S, L}(d(x, y))\right) \tilde{a}\right)\right) \cdot A .
\end{aligned}
$$

Since $\mathscr{E}\left(\left(-\tilde{a} 1_{\rrbracket \tau_{n}, T \rrbracket}\right) \cdot S\right)_{-}$does not vanish on $\rrbracket \tau_{n}, \tau_{n+1} \rrbracket$, we have

$$
L_{-} b^{S}+\tilde{c}^{S L}-\left(L_{-} \tilde{c}^{S}+\int x x^{\top} y F^{S, L}(d(x, y))\right) \tilde{a}=0
$$

and hence (3.28) outside some $P \otimes A$-null set.
Finally, (3.27), (3.28), (3.22) yield

$$
L_{-} \bar{b}^{\top} \bar{c}^{-1} \bar{b}=L_{-} \bar{b}^{\top} \bar{c}^{-1} \bar{c} \tilde{a}=L_{-} \bar{b}^{\top} \tilde{a}=b^{L}
$$

which in turn implies the first equality in (3.29).
On the set $\{\Delta A=0\} \supset\left\{\Delta A^{K}=0\right\}$, the second equality follows from (3.17), (3.18). On $\left\{\Delta A^{K} \neq 0\right\}$ the same equations yield

$$
1=\left(1+\Delta A^{K}\right)-b^{K} \Delta A=\left(1+\Delta A^{K}\right)\left(1-\left(b^{S \star}\right)^{\top}\left(\tilde{c}^{S \star}\right)^{-1} b^{S \star} \Delta A\right)
$$

In view of (3.19) we have

$$
1+b^{K} \Delta A=1+\Delta A^{K}=1+\left(b^{S \star}\right)^{\top}\left(\hat{c}^{S \star}\right)^{-1} b^{S \star} \Delta A
$$

which in turn implies $b^{K}=\left(b^{S \star}\right)^{\top}\left(\hat{c}^{S \star}\right)^{-1} b^{S \star}$ on the set $\left\{\Delta A^{K} \neq 0\right\}$.
COROLLARY 3.20. The adjustment process and the extended adjustment process satisfy the equations

$$
\begin{equation*}
b^{S \star}=\tilde{c}^{S \star} \tilde{a}=\hat{c}^{S \star} \hat{a} \tag{3.30}
\end{equation*}
$$

or, put differently,

$$
A^{S \star}=\tilde{a} \cdot\langle S, S\rangle^{P^{\star}}=\hat{a} \cdot\left\langle M^{S \star}, M^{S \star}\right\rangle^{P^{\star}} .
$$

In the univariate case, this can be written more intuitively in terms of pathwise Radon-Nikodym derivatives:

$$
\tilde{a}_{t}=\frac{d A_{t}^{S \star}}{d\langle S, S\rangle_{t}^{P \star}}, \quad \hat{a}_{t}=\frac{d A_{t}^{S \star}}{d\left\langle M^{S \star}, M^{S \star}\right\rangle_{t}^{P \star}} .
$$

Proof. $\quad b^{S \star}=\tilde{c}^{S \star} \tilde{a}$ follows from (3.28), (3.17), (3.18). Together with (3.21), (3.19), (3.29) we have

$$
\begin{aligned}
\hat{c}^{S \star} \tilde{a} & =\left(\tilde{c}^{S \star}-b^{S \star}\left(b^{S \star}\right)^{\top} \Delta A\right) \tilde{a} \\
& =b^{S \star}\left(1-\left(b^{S \star}\right)^{\top}\left(\tilde{c}^{S \star}\right)^{-1} \tilde{c}^{S \star} \tilde{a} \Delta A\right) \\
& =b^{S \star}\left(1-\left(b^{S \star}\right)^{\top}\left(\tilde{c}^{S \star}\right)^{-1} b^{S \star} \Delta A\right) \\
& =\frac{b^{S \star}}{1+\left(b^{S \star}\right)^{\top}\left(\hat{c}^{S \star}\right)^{-1} b^{S \star} \Delta A} \\
& =\frac{b^{S \star}}{1+\Delta A^{K}},
\end{aligned}
$$

which yields $b^{S \star}=\hat{c}^{S \star} \hat{a}$.
Lemma 3.21. We have $\hat{a} \in L\left(M^{S \star}\right)$.
Proof. Equations (3.30), (3.17), (3.27) imply that

$$
\left(\hat{a}^{\top} \hat{c}^{S \star} \hat{a}\right) \cdot A_{T}=\left(\left(1+\Delta A^{K}\right) \tilde{a}^{\top} b^{S \star}\right) \cdot A_{T}=\left(\tilde{a}^{\top} \bar{b}\right) \cdot A_{T}=\frac{1}{L_{-}} \cdot A_{T}^{L}<\infty
$$

and hence $\hat{a} \in L_{\mathrm{loc}}^{2}\left(M^{S \star}\right) \subset L\left(M^{S \star}\right)$ relative to $P^{\star}$.
DEfinition 3.22. We call

$$
N^{\star}:=-\hat{a} \cdot M^{S \star}
$$

$P^{\star}$-minimal logarithm process.
The terminology is motivated by the fact that $\mathscr{E}\left(N^{\star}\right)$ is essentially the density process of the so-called minimal signed martingale measure relative to $P^{\star}$ instead of $P$ (in the sense of [44], (3.14)).

Lemma 3.23. We have

$$
\frac{L_{0}}{E\left(L_{0}\right)} \mathscr{E}(N)=Z^{P^{\star}} \mathscr{E}\left(N^{\star}\right)
$$

Consequently, $\mathscr{E}\left(N^{\star}\right)$ is the density process of $Q^{\star}$ relative to $P^{\star}$.

Proof. Integration by parts yields

$$
\frac{L_{0} \mathscr{E}(N)}{E\left(L_{0}\right) Z^{P^{\star}}}=\frac{\mathscr{E}(-\tilde{a} \cdot S) L \mathscr{E}\left(A^{K}\right)}{L}=\mathscr{E}\left(-\tilde{a} \cdot S+A^{K}-\left[\tilde{a} \cdot S, A^{K}\right]\right)
$$

The term in parentheses on the right-hand side equals

$$
\begin{align*}
x-\tilde{a} \cdot & M^{S \star}-\left(\tilde{a}^{\top} b^{S \star}\right) \cdot A+b^{K} \cdot A \\
& -\left(\tilde{a} \Delta A^{K}\right) \cdot M^{S \star}-\left(\tilde{a}^{\top} b^{S \star} \Delta A^{K}\right) \cdot A \tag{3.31}
\end{align*}
$$

(cf. [28], I.4.49b). Since

$$
b^{K}=\frac{1}{L_{-}} b^{L}=\tilde{a}^{\top} \bar{b}=\tilde{a}^{\top} b^{S \star}\left(1+\Delta A^{K}\right)
$$

by (3.27), (3.17), the expression in (3.31) equals $-\hat{a} \cdot M^{S \star}=N^{\star}$.
Roughly speaking, the next statement is another way of saying that $S$ is a $Q^{\star}-\sigma$-martingale.

Lemma 3.24. $N^{\star}$ and $S+\left[S, N^{\star}\right]$ are $P^{\star}$ - $\sigma$-martingales, which implies that $S \mathscr{E}\left(N^{\star}\right)$ is a $P^{\star}$ - $\sigma$-martingale as well.

Proof. $\quad N^{\star}$ is a $P^{\star}-\sigma$-martingale by definition. Moreover,

$$
\begin{aligned}
S+\left[S, N^{\star}\right] & =S-\hat{a} \cdot\left[S, M^{S \star}\right] \\
& =S-\hat{a} \cdot\left[M^{S \star}, M^{S \star}\right]-\hat{a} \cdot\left(\left(\Delta A^{S \star}\right) \cdot M^{S \star}\right) \\
& \sim^{\star}\left(b^{S \star}-\hat{c}^{S \star} \hat{a}\right) \cdot A=0
\end{aligned}
$$

by (3.30). The last statement follows as in Lemma 3.14.
Corollary 3.20 expresses the adjustment process in terms of the $P^{\star}$-characteristics of $S$. Of course this only helps if the opportunity-neutral measure is known in the first place. The following important result characterizes $L$ and $\tilde{a}$ directly in terms of $P$-characteristics.

THEOREM 3.25. The opportunity process is the unique semimartingale $L$ such that:

1. $L, L_{-}$are $(0,1]$-valued,
2. $L_{T}=1$,
3. The joint characteristics of $(S, L)$ solve the equation

$$
\begin{equation*}
b^{L}=L_{-} \bar{b}^{\top} \bar{c}^{-1} \bar{b} \tag{3.32}
\end{equation*}
$$

outside some $P \otimes A$-null set, where $\bar{b}, \bar{c}$ are defined as in (3.23), (3.25),
4.

$$
\begin{gather*}
a \mathscr{E}\left(\left(-a 1_{\rrbracket \tau, T \rrbracket}\right) \cdot S\right)_{-} 1_{\rrbracket \tau, T \rrbracket} \in \bar{\Theta}  \tag{3.33}\\
\mathscr{E}\left(\left(-a 1_{\rrbracket \tau, T \rrbracket}\right) \cdot S\right) L \text { is of class }(D) \tag{3.34}
\end{gather*}
$$

hold for $a:=\bar{c}^{-1} \bar{b}$ and any stopping time $\tau$.
In this case $a=\bar{c}^{-1} \bar{b}$ meets the requirement of an adjustment process $\tilde{a}$ in Lemma 3.7.

Proof. Suppose that $L$ is the opportunity process. Properties 1 and 2 are shown in Lemmas 3.2 and 3.10. Equation (3.29) and $b^{L}=L_{-} b^{K}$ yield (3.32). By (3.17), (3.18), (3.21), (3.29) we have

$$
\begin{aligned}
\left(a^{\top} \hat{c}^{S \star} a\right) \cdot A_{T} & \leq\left(a^{\top} \tilde{c}^{S \star} a\right) \cdot A_{T}=\left(\left(b^{S \star}\right)^{\top}\left(\tilde{c}^{S \star}\right)^{-1} b^{S \star}\right) \cdot A_{T} \\
& =\frac{1}{1+\Delta A^{K}} \cdot A_{T}^{K}<\infty,
\end{aligned}
$$

which implies $a \in L_{\mathrm{loc}}^{2}\left(M^{S \star}\right)$ relative to $P^{\star}$ by [28], III.4.3. Similarly, we have $a \in L\left(A^{S \star}\right)$ because $\left|a^{\top} b^{S \star}\right| \cdot A_{T} \leq \frac{1}{1+\Delta A^{K}} \cdot A_{T}^{K}<\infty$. Together, it follows that $a \in L(S)$.

More specifically, we have

$$
a \cdot A^{S \star}=\left(a^{\top} b^{S \star}\right) \cdot A=\frac{b^{L}}{\left(1+\Delta A^{K}\right) L_{-}} \cdot A
$$

and likewise for $\tilde{a}$ by (3.17), (3.27). Similarly, (3.27-3.29) yield

$$
\left\langle(a-\tilde{a}) \cdot M^{S \star},(a-\tilde{a}) \cdot M^{S \star}\right\rangle P^{P^{\star}} \leq\left((a-\tilde{a})^{\top} \tilde{c}^{S \star}(a-\tilde{a})\right) \cdot A=0
$$

which implies $(a-\tilde{a}) \cdot M^{S \star}=0$. Together, we have $a \cdot S=\tilde{a} \cdot S$. Hence one may choose $\tilde{a}=a$ in Lemma 3.7.

Finally, (3.33) follows from (3.12) and (3.34) from Lemma 3.2.
Conversely, let $L^{\prime}$ be a semimartingale satisfying properties $1-4$ with $\bar{b}^{\prime}, \bar{c}^{\prime}$ as in (3.23) and (3.25). Define $K^{\prime}:=\frac{1}{L_{-}^{\prime}} \cdot L^{\prime}$ and $N^{\prime}:=K^{\prime}-a \cdot S-\left[a \cdot S, K^{\prime}\right]$. We use the notation $L^{\prime}, \bar{b}^{\prime}, \bar{c}^{\prime}, K^{\prime}, N^{\prime}$ in this part of the proof because $L^{\prime}$ is yet to be shown to coincide with the true opportunity process. From

$$
\left[S, K^{\prime}\right]=\frac{1}{L_{-}^{\prime}} \cdot\left[M^{S}, M^{L^{\prime}}\right]+\left(\Delta A^{S}\right) \cdot M^{K^{\prime}}+\left(\Delta A^{K^{\prime}}\right) \cdot S
$$

and standard results (cf. [28], I.4.24, III.3.14) it follows that

$$
\left[S, K^{\prime}\right]=\left[S^{c}, K^{\prime c}\right]+\int_{[0, \cdot] \times \mathbb{R}^{d} \times \mathbb{R}} x y \mu^{\left(S, K^{\prime}\right)}(d(t, x, y))
$$

is an $\mathbb{R}^{d}$-valued special semimartingale with compensator $\left(c^{S K^{\prime}}+\int x y F^{S, K^{\prime}} \times\right.$ $(d(x, y))) \cdot A$. For $n \in \mathbb{N}$ define the predictable set $D_{n}:=\{|a| \leq n\}$. Since $1_{D_{n}}$ and $a 1_{D_{n}}$ are bounded, we have that

$$
1_{D_{n}} \cdot N^{\prime}=1_{D_{n}} \cdot K^{\prime}-\left(1_{D_{n}} a\right) \cdot S-\left(1_{D_{n}} a\right) \cdot\left[S, K^{\prime}\right]
$$

is a special semimartingale as well with compensator

$$
\begin{gathered}
\left(1_{D_{n}} b^{K^{\prime}}-1_{D_{n}} a^{\top}\left(b^{S}+c^{S K^{\prime}}+\int x y F^{S, K^{\prime}}(d(x, y))\right)\right) \cdot A \\
=\left(\left(\frac{b^{L^{\prime}}}{L_{-}^{\prime}}-\bar{b}^{\prime \top} \bar{c}^{\prime-1} \bar{b}^{\prime}\right) 1_{D_{n}}\right) \cdot A=0 .
\end{gathered}
$$

Consequently, $1_{D_{n}} \cdot N^{\prime}$ is actually a local martingale. Since $D_{n} \uparrow \Omega \times[0, T]$ up to an evanescent set, $N^{\prime}$ is a $\sigma$-martingale (cf. Remark A.5).

Similarly, we have that

$$
\begin{aligned}
1_{D_{n}} \cdot & \left(S^{i}+\left[S^{i}, N^{\prime}\right]\right) \\
= & 1_{D_{n}} \cdot S^{i}+1_{D_{n}} \cdot\left[S^{i}, K^{\prime}\right] \\
& \quad-\sum_{j=1}^{n}\left(1_{D_{n}} a^{j}\right) \cdot\left[S^{i}, S^{j}\right]-\sum_{j=1}^{n}\left(1_{D_{n}} a^{j}\right) \cdot\left[S^{i},\left[S^{j}, K^{\prime}\right]\right]
\end{aligned}
$$

is a special semimartingale with compensator

$$
\begin{aligned}
\left(1_{D_{n}}\right. & \left(b^{S}+c^{S K^{\prime}}+\int x y F^{S, K^{\prime}}(d(x, y))\right. \\
& \left.\left.-c^{S} a-\int x\left(x^{\top} a\right)(1+y) F^{S, K^{\prime}}(d(x, y))\right)^{i}\right) \cdot A \\
\quad= & \left(1_{D_{n}}\left(\bar{b}^{\prime}-\bar{c}^{\prime} a\right)^{i}\right) \cdot A
\end{aligned}
$$

for $i=1, \ldots, d$. Since $\bar{b}^{\prime}-\bar{c}^{\prime} a=\bar{b}^{\prime}-\bar{c}^{\prime} \bar{c}^{\prime-1} \bar{b}^{\prime}=0$ by Remark 3.18, it follows that the process $1_{D_{n}} \cdot\left(S^{i}+\left[S^{i}, N^{\prime}\right]\right)$ is a local martingale. This implies that $S+\left[S, N^{\prime}\right]$ is a $\sigma$-martingale as well.

Fix a stopping time $\tau$. Let $\vartheta:=a \mathscr{E}\left(-a 1_{\rrbracket \tau, T \rrbracket} \cdot S\right)_{-} 1_{\rrbracket \tau, T \rrbracket}$ and

$$
Z:=(1-\vartheta \cdot S) L^{\prime}=\mathscr{E}\left(\left(-a 1_{\rrbracket \tau, T \rrbracket}\right) \cdot S\right) L^{\prime}
$$

[In (3.33) and (3.34) it is implicitly assumed that $a \in L(S)$ for the integral to make sense. By similar arguments as in the first part of the proof one can show that this integrability condition is in fact implied by properties $1-3$ of Theorem 3.25.]

Since $N^{\prime}$ and $S+\left[S, N^{\prime}\right]$ are $\sigma$-martingales,

$$
\frac{Z}{Z^{\tau}}=\mathscr{E}\left(1_{\rrbracket \tau, T \rrbracket} \cdot K^{\prime}\right) \mathscr{E}\left(\left(-a 1_{\rrbracket \tau, T \rrbracket}\right) \cdot S\right)=\mathscr{E}\left(N^{\prime}-N^{\prime \tau}\right)
$$

and

$$
\begin{aligned}
& \frac{Z}{Z^{\tau}}\left(S-S^{\tau}\right) \\
& \quad=\mathscr{E}\left(N^{\prime}-N^{\prime \tau}\right)_{-} \cdot\left(\left(S-S^{\tau}\right)_{-} \cdot\left(N^{\prime}-N^{\prime \tau}\right)+1_{\rrbracket \tau, T \rrbracket} \cdot\left(S+\left[S, N^{\prime}\right]\right)\right)
\end{aligned}
$$

are $\sigma$-martingales as well.
We show that $\vartheta$ is efficient on $\rrbracket \tau, T \rrbracket$. Indeed, from (3.34) and Lemma A. 7 it follows that $Z-Z^{\tau}=\left(Z_{\tau} 1_{\rrbracket \tau, T \rrbracket}\right) \cdot \frac{Z}{Z^{\tau}}$ is a martingale. It is even a squareintegrable martingale because $Z_{T}-Z_{\tau} \in L^{2}(P)$. Let $\psi$ be a simple strategy with $\psi 1_{\llbracket 0, \tau \rrbracket}=0$. The same arguments as in step 1 of the proof of Lemma 2.4 yield that $(\psi \cdot S) Z=\left(\left(Z_{\tau} \psi\right) \cdot\left(S-S^{\tau}\right)\right) \frac{Z}{Z^{\tau}}$ is a martingale. Consequently,

$$
\begin{aligned}
& E\left(\left(1-(\vartheta+\psi) \cdot S_{T}\right)^{2}\right) \\
& \quad \geq E\left(\left(1-\vartheta \cdot S_{T}\right)^{2}\right)-2 E\left(\left(1-\vartheta \cdot S_{T}\right) L_{T}^{\prime}\left(\psi \cdot S_{T}\right)\right) \\
& \quad=E\left(\left(1-\vartheta \cdot S_{T}\right)^{2}\right)
\end{aligned}
$$

which implies the optimality of $\vartheta$. Since $Z-Z^{\tau}$ is a martingale, Lemma 3.2 yields that $L^{\prime}$ is the opportunity process.

Condition (3.33) looks somewhat unpleasant because of the involved definition of $\bar{\Theta}$. The following example shows that uniqueness in Theorem 3.25 does not generally hold without this condition. For related considerations see also [44] and [10].

Example 3.26. Let $T=1$ and $S$ be a standard Wiener process. By Theorem 3.25 the opportunity and adjustment processes are $L=1$ and $\tilde{a}=0$. Choose some doubling-type strategy $\psi \in L(S)$ with $1-\psi \cdot S \geq \frac{1}{2}$ and $1-\psi \cdot S_{T}=\frac{1}{2}$. Of course, $\psi$ cannot be admissible. We write $1-\psi \cdot S=\mathscr{E}(-\bar{a} \cdot S)$ with $\bar{a}:=\frac{\psi}{1-\psi \cdot S_{-}}$. Define

$$
\bar{L}:=\frac{1}{2 \mathscr{E}(-\bar{a} \cdot S)}=\frac{1}{2} \mathscr{E}\left(\bar{a} \cdot S+\bar{a}^{2} \cdot[S, S]\right)
$$

Straightforward calculations yield that $\bar{L}$ satisfies conditions 1-3 in Theorem 3.25. Moreover, $\bar{a}$ is the corresponding process in condition 4. Since $\mathscr{E}\left(\left(-\bar{a} 1_{\rrbracket \tau, T \rrbracket}\right) \cdot\right.$ $S) \bar{L}=\bar{L}^{\tau}$ is bounded, (3.34) is satisfied as well.

It is interesting to note that the "variance-optimal logarithm process" $\bar{N}$ corresponding to this wrong choice of $\bar{L}, \bar{a}$ satisfies $\mathscr{E}(\bar{N})=\frac{\bar{L}}{\bar{L}_{0}} \mathscr{E}(-\bar{a} \cdot S)=1$, that is, it coincides with the true variance-optimal logarithm process. In particular,

$$
\frac{d Q^{\star}}{d P}=\frac{1-\psi \cdot S_{T}}{E\left(1-\psi \cdot S_{T}\right)}
$$

which parallels the last expression in (3.16). Nevertheless, $\psi$ is not an efficient strategy on $\rrbracket 0, T \rrbracket$ because it is not admissible.

In concrete models, it may be easier to verify the following sufficient condition instead of (3.33), (3.34).

Lemma 3.27. Let L be a special semimartingale satisfying conditions $1-3$ in Theorem 3.25 with $\bar{b}, \bar{c}$ defined as in (3.23), (3.25). If $a:=\bar{c}^{-1} \bar{b}$ satisfies

$$
\sup \left\{E\left(\mathscr{E}\left(\left(-a 1_{\rrbracket \tau, T \rrbracket}\right) \cdot S\right)_{\sigma}^{2}\right): \sigma \text { stopping time }\right\}<\infty
$$

for any stopping time $\tau$, then condition 4 holds as well, that is, $L$ is the opportunity process.

Proof. Condition (3.34) is obvious because $L$ is bounded. Let $Q$ be an $\mathrm{S} \sigma \mathrm{MM}$ with density process $Z^{Q}$ and $\frac{d Q}{d P} \in L^{2}(P)$. Integration by parts yields that $(\vartheta \cdot S) Z^{Q}$ is a $\sigma$-martingale for

$$
\vartheta:=a \mathscr{E}\left(\left(-a 1_{\rrbracket \tau, T \rrbracket}\right) \cdot S\right)_{-} 1_{\rrbracket \tau, T \rrbracket}
$$

[cf. (2.2)]. Since $\sup _{t \in[0, T]}\left|Z_{t}^{Q}\right| \in L^{2}(P)$ by Doob's inequality and $1-\vartheta \cdot S$ is an $L^{2}$-semimartingale, we have that $(\vartheta \cdot S) Z^{Q}$ is of class (D) and hence a martingale (cf. Lemma A.7). Using Corollary 2.5 we obtain (3.33).
3.6. When does $P^{\star}=P$ hold? The opportunity-neutral measure plays a key role in quadratic hedging. Therefore we want to have a closer look at the question when $P^{\star}$ equals $P$. In line with [42], we call

$$
\widehat{K}:=\left(\left(b^{S}\right)^{\top}\left(\hat{c}^{S}\right)^{-1} b^{S}\right) \cdot A
$$

mean-variance tradeoff (MVT) process. Similarly, the MVT process relative to $P^{\star}$ is denoted by $\widehat{K}^{\star}$, that is,

$$
\widehat{K}^{\star}:=\left(\left(b^{S \star}\right)^{\top}\left(\hat{c}^{S \star}\right)^{-1} b^{S \star}\right) \cdot A .
$$

Observe that $\widehat{K}^{\star}=A^{K}$ by (3.29).

Proposition 3.28. The following statements are equivalent:

1. $P^{\star}=P$.
2. $K$ (or equivalently $L$ ) is a predictable process of finite variation and $L_{0}$ is deterministic.
3. $K=\widehat{K}$ and $L_{0}$ is deterministic.
4. $K=\widehat{K}^{\star}$ and $L_{0}$ is deterministic.
5. $\mathscr{E}(\widehat{K})_{T}$ is finite and deterministic.
6. $\mathscr{E}\left(\widehat{K}^{\star}\right)_{T}$ is deterministic.

In this case the opportunity process equals $L=\mathscr{E}(\widehat{K}) / \mathscr{E}(\widehat{K})_{T}$.

Proof. $\quad 1 \Rightarrow 2$ : Since $1=Z^{P \star}=L /\left(E\left(L_{0}\right) \mathscr{E}\left(A^{K}\right)\right)$, we have that $L$ and hence also $K=\mathscr{L}(L)$ are predictable processes of finite variation. $L_{0}$ is deterministic because $Z_{0}^{P \star}=1$.
$2 \Rightarrow 4$ : This is obvious because $K=A^{K}=\widehat{K}^{\star}$.
$4 \Rightarrow 1$ : This follows from

$$
Z^{P \star}=\frac{L}{E\left(L_{0}\right) \mathscr{E}\left(A^{K}\right)}=\frac{L_{0} \mathscr{E}(K)}{E\left(L_{0}\right) \mathscr{E}\left(\widehat{K}^{\star}\right)}=1
$$

$1 \Rightarrow 6$ : This follows from $Z_{T}^{P^{\star}}=1$ and $\widehat{K}^{\star}=A^{K}$.
$6 \Rightarrow 1$ : This holds because $Z_{T}^{P \star}=1 /\left(E\left(L_{0}\right) \mathscr{E}\left(\widehat{K}^{\star}\right)_{T}\right)$ is deterministic.
$1 \Rightarrow 3$ : In view of $(1 \Rightarrow 2)$, this follows from $K=A^{K}=\widehat{K}^{\star}=\widehat{K}$.
$3 \Rightarrow 5$ : This follows from $1=L_{T}=L_{0} \mathscr{E}(K)_{T}$.
$5 \Rightarrow 2$ : Let $L:=\mathscr{E}(\widehat{K}) / \mathscr{E}(\widehat{K})_{T}$. Since $\widehat{K}$ is an increasing predictable process, $L$ is a $(0,1]$-valued increasing predictable process. The predictability of $L$ implies $c^{S L}=0$ and $y=\Delta L_{t}\left(F_{t}^{S, L}(d(x, y)) A(d t)\right)$-almost everywhere. If $\bar{b}, \bar{c}$ are defined as in (3.23), (3.25), we have $\bar{b}=(1+\Delta \widehat{K}) b^{S}, \bar{c}=(1+\Delta \widehat{K}) \tilde{c}^{S}$ and hence

$$
L_{-} \bar{b}^{\top} \bar{c}^{-1} \bar{b}=L_{-}\left(1+\left(b^{S}\right)^{\top}\left(\hat{c}^{S}\right)^{-1} b^{S} \Delta A\right)\left(b^{S}\right)^{\top}\left(\tilde{c}^{S}\right)^{-1} b^{S}
$$

Observe that (3.19-3.21) can be derived literally for $P$ instead of $P^{\star}$. We obtain

$$
L_{-} \bar{b}^{\top} \bar{c}^{-1} \bar{b}=L_{-}\left(b^{S}\right)^{\top}\left(\hat{c}^{S}\right)^{-1} b^{S}=b^{L}
$$

which implies that $L$ satisfies conditions $1-3$ in Theorem 3.25. If we can show that $L$ is the true opportunity process, then $P^{\star}=P$ follows from Lemma 3.15.

Fix any stopping time $\tau$. For $a:=\bar{c}^{-1} \bar{b}=\left(\tilde{c}^{S}\right)^{-1} b^{S \star}$ and $X:=\left(-a 1_{\rrbracket \tau, T \rrbracket}\right) \cdot S$ we have

$$
\begin{aligned}
\left\langle M^{X}, M^{X}\right\rangle_{T} & =\left(a^{\top} \hat{c}^{S} a 1_{\rrbracket \tau, T \rrbracket}\right) \cdot A_{T} \\
& \leq\left(a^{\top} \tilde{c}^{S} a 1_{\rrbracket \tau, T \rrbracket}\right) \cdot A_{T} \\
& =\left(\left(b^{S}\right)^{\top}\left(\tilde{c}^{S}\right)^{-1} \tilde{c}^{S}\left(\tilde{c}^{S}\right)^{-1} b^{S} 1_{\rrbracket \tau, T \rrbracket}\right) \cdot A_{T} \\
& =\left(\frac{1_{\rrbracket \tau, T \rrbracket}}{1+\Delta \widehat{K}^{\top}} \bar{b}^{\top} \bar{c}^{-1} \bar{b}\right) \cdot A_{T} \\
& \leq\left(\left(b^{S}\right)^{\top}\left(\hat{c}^{S}\right)^{-1} b^{S}\right) \cdot A_{T} \\
& =\widehat{K}_{T} \leq \mathscr{E}(\widehat{K})_{T}
\end{aligned}
$$

Similarly, we have

$$
\operatorname{var}\left(A^{X}\right)_{T}=\left|a^{\top} b^{S} 1_{\rrbracket \tau, T \rrbracket}\right| \cdot A_{T}=\left(\frac{1_{\rrbracket \tau, T \rrbracket}}{1+\Delta \widehat{K}} \bar{b}^{\top} \bar{c}^{-1} \bar{b}\right) \cdot A_{T} \leq \mathscr{E}(\widehat{K})_{T}
$$

for the variation process of $A^{X}$. In view of Lemmas A. 3 and 3.27, $L$ is the opportunity process.

To relate the condition $P^{\star}=P$ to earlier literature, we define (myopic) portfolio weights

$$
\begin{align*}
& \tilde{\lambda}:=\left(\tilde{c}^{S}\right)^{-1} b^{S}, \\
& \hat{\lambda}:=(1+\Delta \widehat{K}) \tilde{\lambda} \tag{3.35}
\end{align*}
$$

in accordance with [42]. Repeating the arguments leading to (3.30) under $P$ rather than $P^{\star}$ yields $\hat{c}^{S} \hat{\lambda}=b^{S}$ [which implies that $\hat{\lambda}=\left(\hat{c}^{S}\right)^{-1} b^{S}$ if $\hat{c}^{S}$ is invertible]. By Theorem 1 of [43] we have $\hat{\lambda} \in L\left(M^{S}\right)$.

DEFINITION 3.29. If $\mathscr{E}\left(-\hat{\lambda} \cdot M^{S}\right)$ is of class (D) and hence a martingale, then it is the density process of some $\mathrm{S} \sigma \mathrm{MM} Q$. Only slightly extending [44], (3.14) we call $Q$ the minimal signed martingale measure (minimal $S \sigma M M$ ).

In view of Proposition 3.28, the following corollary can be interpreted as an extension of Proposition 5.1 in [30]. It also extends sufficient conditions for $Q^{\star}=$ $Q$ given in [44], Examples 1 and 2.

Corollary 3.30. Suppose $\mathscr{E}(-\tilde{a} \cdot S)_{T} \neq 0$ almost surely. Then there is equivalence between:

1. $P^{\star}=P$,
2. $\widehat{K}_{T}$ is finite, the minimal $S \sigma M M Q$ exists, $Q^{\star}=Q$, and $\tilde{a}$ can be chosen as $\tilde{\lambda}$.

The implication $1 \Rightarrow 2$ still holds without the assumption on $\mathscr{E}(-\tilde{a} \cdot S)$.
Proof. $\quad 1 \Rightarrow 2$ : This follows from Lemma 3.23, Theorem 3.25, and (3.17), (3.18).
$2 \Rightarrow 1$ : As in the proof of Lemma 3.17 it follows that $b^{S}=\hat{c}^{S}\left(\hat{c}^{S}\right)^{-1} b^{S}$ and hence $\hat{\lambda}^{\top} b^{S}=\left(b^{S}\right)^{\top}\left(\hat{c}^{S}\right)^{-1} b^{S}$. Hence, the density process of $Q$ equals

$$
\begin{aligned}
\mathscr{E}\left(-\hat{\lambda} \cdot M^{S}\right) & =\mathscr{E}\left(\left(\hat{\lambda}^{\top} b^{S}\right) \cdot A-\hat{\lambda} \cdot S\right) \\
& =\mathscr{E}\left(\left(\left(b^{S}\right)^{\top}\left(\hat{c}^{S}\right)^{-1} b^{S}\right) \cdot A-((1+\Delta \widehat{K}) \tilde{\lambda}) \cdot S\right) \\
& =\mathscr{E}(\widehat{K}-\tilde{\lambda} \cdot S-(\Delta \widehat{K}) \cdot(\tilde{\lambda} \cdot S)) \\
& =\mathscr{E}(\widehat{K}-\tilde{\lambda} \cdot S-[\widehat{K}, \tilde{\lambda} \cdot S]) \\
& =\mathscr{E}(\widehat{K}) \mathscr{E}(-\tilde{\lambda} \cdot S),
\end{aligned}
$$

where the fourth equality follows from [28], I.4.49b and the last from Yor's formula. This density process equals $\frac{L}{E\left(L_{0}\right)} \mathscr{E}(-\tilde{a} \cdot S)$ by $Q^{\star}=Q$ and Proposition 3.13. Since $\tilde{a}=\tilde{\lambda}$ and $\mathscr{E}(-\tilde{a} \cdot S)$ never vanishes (cf. [28], I.4.61), we have that $L=E\left(L_{0}\right) \mathscr{E}(\widehat{K})$ is predictable with $L_{0}=E\left(L_{0}\right)$. The assertion follows now from Proposition $3.28(2 \Rightarrow 1)$.

Finally, we consider the situation of deterministic mean-variance tradeoff, which is the focus of [42].

COROLLARY 3.31. If the MVT process $\widehat{K}$ is finite and deterministic, then $L:=\mathscr{E}(\widehat{K}) / \mathscr{E}(\widehat{K})_{T}$ is the opportunity process, $K:=\widehat{K}$ is the modified meanvariance tradeoff process, and $P^{\star}=P$.

Proof. This follows from Proposition $3.28(5 \Rightarrow 1,3)$ and from $1=L_{T}=$ $L_{0} \mathscr{E}(K)_{T}$.
3.7. Determination of the opportunity process. Unless we are in the fortunate situation of Corollary 3.31 or at least Proposition 3.28, the crucial step in concrete applications is to determine the opportunity process $L$. This is relatively easy in discrete time.

EXAMPLE 3.32. Suppose that we are actually considering a discrete-time model, that is, $A_{t}=[t]:=\max \{n \in \mathbb{N}: n \leq t\}$ and $\mathscr{F}_{t}=\mathscr{F}_{[t]}$ for $t \in[0, T]$ with $T \in \mathbb{N}$. In this case all processes in this paper are (or can be chosen) piecewise constant between integer times. For ease of notation suppose that $d=1$ (only one tradable asset). By [28], II.3.11 we have $b_{t}^{L}=E\left(\Delta L_{t} \mid \mathscr{F}_{t-1}\right), \bar{b}_{t}=$ $E\left(\Delta S_{t} L_{t} / L_{t-1} \mid \mathscr{F}_{t-1}\right)$, and $\bar{c}_{t}=E\left(\left(\Delta S_{t}\right)^{2} L_{t} / L_{t-1} \mid \mathscr{F}_{t-1}\right)$ for $t \in\{1,2, \ldots, T\}$. Consequently, (3.32) can be rewritten as

$$
\begin{equation*}
L_{t-1}=E\left(L_{t} \mid \mathscr{F}_{t-1}\right)-\frac{\left(E\left(\Delta S_{t} L_{t} \mid \mathscr{F}_{t-1}\right)\right)^{2}}{E\left(\left(\Delta S_{t}\right)^{2} L_{t} \mid \mathscr{F}_{t-1}\right)}, \tag{3.36}
\end{equation*}
$$

that is, the opportunity process is determined by a simple backward recursion starting in $L_{T}=1$. For the adjustment process we have

$$
\tilde{a}_{t}=\frac{\bar{b}_{t}}{\bar{c}_{t}}=\frac{E\left(\Delta S_{t} L_{t} \mid \mathscr{F}_{t-1}\right)}{E\left(\left(\Delta S_{t}\right)^{2} L_{t} \mid \mathscr{F}_{t-1}\right)} .
$$

The previous example indicates that the characteristic equation (3.32) may be interpreted as the continuous-time analogue of a backward recursion. True continuous-time models are typically Markovian in $S_{t}$ or at least ( $S_{t}, Y_{t}$ ) with some additional process $Y$ as, for example, stochastic volatility. If one makes the natural assumption $L_{t}=f\left(t, S_{t}, Y_{t}\right)$ with some $C^{2}$-function $f$, then (3.32) can be rewritten as an integro-differential equation for $f$ by means of Itô's formula. But as it is not obvious whether the smoothness assumption is justified, it may require substantial effort to make this statement precise. In [11] and ongoing research, $L$ is determined explicitly by an ansatz of the above type in specific stochastic volatility models.

Alternatively, the process $L$ can be interpreted as the solution to some backward stochastic differential equation (BSDE). To this end, we use the martingale representation theorem (cf. [28], III.4.24) to write the martingale part of $L$ as

$$
M^{L}=J \cdot S^{c}+W *\left(\mu^{S}-v^{S}\right)+U
$$

with some $J \in L_{\mathrm{loc}}^{2}\left(S^{c}\right), W \in G_{\mathrm{loc}}\left(\mu^{S}\right)$ and some local martingale $U \in \mathscr{H}_{\mathrm{loc}}^{2}$ such that $\left\langle U^{c}, S^{c}\right\rangle=0$ and $M_{\mu^{s}}^{P}(\Delta U \mid \widetilde{\mathscr{P}})=0$ in the sense of [28], III.3c. Using the notation

$$
\widehat{W}_{t}:=E\left(W\left(t, \Delta S_{t}\right) \mid \mathscr{F}_{t-}\right),
$$

the quadruple $(J, W, L, U)$ solves the BSDE

$$
\begin{align*}
L=J \cdot & S^{c}+W *\left(\mu^{S}-v^{S}\right)+U \\
+ & \left(\left(b^{S}+c^{S} \frac{J}{L_{-}}+\int \frac{W(x)-\widehat{W}}{L_{-}} x F^{S}(d x)\right)^{\top}\right. \\
& \times\left(c^{S}+\int x x^{\top}\left(1+\frac{W(x)-\widehat{W}}{L_{-}}\right) F^{S}(d x)\right)^{-1}  \tag{3.37}\\
& \left.\times\left(b^{S}+c^{S} \frac{J}{L_{-}}+\int \frac{W(x)-\widehat{W}}{L_{-}} x F^{S}(d x)\right) L_{-}\right) \cdot A
\end{align*}
$$

$$
L_{T}=1
$$

However, it is not obvious whether this representation is of any use.
One should note that (3.37) is not related to the BSDEs (3.6) and (4.10) in [44], which characterize the adjustment process and the optimal hedge. The latter are hard to use in practice because their terminal values involve the $L^{2}$-projection of 1, respectively, H on $K_{2}(0)$, which is generally unknown. If at all, one may rather observe a certain similarity between (3.36) and the recursive expression (2.1) in [44] for the adjustment process in discrete time. Mania and Tevzadze [34, 36] derive BSDE's for $1 / L$ in the case of a continuous asset price process $S$. These equations are quite different from both (3.37) and (3.32).
4. On the quadratic hedging problem. We now come back to the hedging problem from Definition 2.10. The processes and measures $\lambda^{(\tau)}, M^{(\tau)}, L, \tilde{a}, \hat{a}, K$, $N, Q^{\star}, Z^{P^{\star}}, P^{\star}, \bar{b}, \bar{c}$ are defined as in the previous section. Recall that $P$ is the default probability measure for expectations, martingales and so forth.
4.1. Mean value process and pure hedge coefficient. If $S$ is a martingale, the mean value process $V_{t}=E\left(H \mid \mathscr{F}_{t}\right)$ leads to the optimal hedge via (1.1). If $S$ fails to be a martingale, a similar role is played by the conditional expectation of $H$ relative to the variance-optimal $\mathrm{S} \sigma \mathrm{MM} Q^{\star}$. By the generalized Bayes formula we have

$$
\begin{equation*}
E_{Q^{\star}}\left(H \mid \mathscr{F}_{t}\right)=E\left(H \mathscr{E}\left(N-N^{t}\right)_{T} \mid \mathscr{F}_{t}\right) \tag{4.1}
\end{equation*}
$$

if $Q^{\star}$ is a true probability measure. In the general case we use the right-hand side of (4.1) as a substitute for the possibly undefined conditional expectation.

LEmmA 4.1. There is a unique semimartingale $V$ satisfying

$$
\left.\begin{array}{rl}
V_{t} & =E\left(\left.H \mathscr{E}\left(N-N^{t}\right)_{T}\right|_{\mathscr{F}} ^{t}\right.
\end{array}\right)
$$

for $t \in[0, T]$. Moreover, $\left(V_{S} M_{S}^{(t)}\right)_{s \in[t, T]}$ is a martingale for any $t \in[0, T]$.
Proof. In this proof $\varphi$ denotes an optimal hedging strategy for arbitrary initial endowment $v_{0} \in L^{2}\left(\Omega, \mathscr{F}_{0}, P\right)$ or, alternatively, $\left(v_{0}, \varphi\right)$ denotes an optimal endowment/strategy pair. Moreover, let $G:=v_{0}+\varphi \cdot S$ and define a squareintegrable martingale $Z$ by its terminal value $Z_{T}:=G_{T}-H$. Finally, we set $V:=G-\frac{Z}{L}$. The optimality of $\varphi$ implies that

$$
\begin{aligned}
0 & \leq E\left(\left(G_{T}+\varepsilon \vartheta \cdot S_{T}-H\right)^{2}\right)-E\left(\left(G_{T}-H\right)^{2}\right) \\
& =2 \varepsilon E\left(\left(\vartheta \cdot S_{T}\right) Z_{T}\right)+\varepsilon^{2} E\left(\left(\vartheta \cdot S_{T}\right)^{2}\right)
\end{aligned}
$$

for any simple strategy $\vartheta$ and any $\varepsilon \in \mathbb{R}$. Therefore

$$
\begin{equation*}
E\left(\left(\vartheta \cdot S_{T}\right) Z_{T}\right)=0 \tag{4.4}
\end{equation*}
$$

for any simple $\vartheta$, which implies that $S Z$ is a $\sigma$-martingale. By Remark 2.7 we have that $(\vartheta \cdot S) Z$ is a martingale for any $\vartheta \in \bar{\Theta}$. In particular, $(G-V) M^{(t)}=\left(1-\lambda^{(t)} \cdot\right.$ $S) Z$ is a martingale for any fixed $t \in[0, T]$. By Lemma 3.2(3), $\left.\left(G_{s} M_{s}^{(t)}\right)\right)_{s \in[t, T]}$ and hence also $\left.\left(V_{s} M_{s}^{(t)}\right)\right)_{s \in[t, T]}$ is a martingale. Using Lemma 3.7, we have

$$
\begin{aligned}
E\left(H \mathscr{E}\left(N-N^{t}\right)_{T} \mid \mathscr{F}_{t}\right) L_{t} & =E\left(V_{T} L_{T}\left(1-\lambda^{(t)} \cdot S_{T}\right) \mid \mathscr{F}_{t}\right) \\
& =V_{t} L_{t}\left(1-\lambda^{(t)} \cdot S_{t}\right) \\
& =V_{t} L_{t}
\end{aligned}
$$

which shows (4.2).
Along the same lines as Lemma 3.23 it follows that

$$
\begin{equation*}
\mathscr{E}\left(N-N^{t}\right)=\mathscr{E}\left(N^{\star}-\left(N^{\star}\right)^{t}\right) \frac{Z^{P^{\star}}}{\left(Z^{P^{\star}}\right)^{t}} \tag{4.5}
\end{equation*}
$$

Consequently,

$$
\begin{aligned}
E_{P^{\star}}\left(H \mathscr{E}\left(N^{\star}-\left(N^{\star}\right)^{t}\right)_{T} \mid \mathscr{F}_{t}\right) & =E\left(\left.H \mathscr{E}\left(N^{\star}-\left(N^{\star}\right)^{t}\right)_{T} \frac{Z_{T}^{P^{\star}}}{Z_{t}^{P^{\star}}} \right\rvert\, \mathscr{F}_{t}\right) \\
& =E\left(H \mathscr{E}\left(N-N^{t}\right)_{T} \mid \mathscr{F}_{t}\right)
\end{aligned}
$$

which yields (4.3). The uniqueness (up to indistinguishability) of $V$ is obvious.

DEFINITION 4.2. We call $V$ from Lemma 4.1 mean value process of the option.

The following technical statements mean essentially that $V$ is a $Q^{\star}-\sigma$-martingale.

Lemma 4.3. We have $1 . V+[V, N]$ and hence $V \mathscr{E}(N)$ are $\sigma$-martingales. 2. $V+\left[V, N^{\star}\right]$ and hence $V \mathscr{E}\left(N^{\star}\right)$ are $P^{\star}$ - $\sigma$-martingales.

Proof. 1. Fix $n \in \mathbb{N}$. If $\left(\tau_{n}\right)_{n \in \mathbb{N}}$ denotes the sequence of stopping times from the proof of Lemma 3.2, then

$$
\mathscr{E}\left(N-N^{\tau_{n}}\right)_{-}=\frac{L_{-}}{L_{-}^{\tau_{n}}}\left(1-\lambda^{\left(\tau_{n}\right)} \cdot S_{-}\right) \neq 0
$$

on $\rrbracket \tau_{n}, \tau_{n+1} \rrbracket$. For $t \in \rrbracket \tau_{n}, T \rrbracket$ we have

$$
\mathscr{E}\left(N-N^{t}\right) \mathscr{E}\left(N-N^{\tau_{n}}\right)_{t}=\mathscr{E}\left(N-N^{\tau_{n}}\right),
$$

which implies that $\left(L_{\tau_{n}} 1_{\rrbracket \tau_{n}, T \rrbracket}\right) \cdot\left(V \mathscr{E}\left(N-N^{\tau_{n}}\right)\right)$ is a martingale by (4.2). Integration by parts and Lemma 3.14 yield that

$$
\begin{aligned}
& 1_{\rrbracket \tau_{n}, \tau_{n+1} \rrbracket} \cdot(V+[V, N]) \\
& \quad=\frac{1_{\rrbracket \tau_{n}, \tau_{n+1} \rrbracket}^{\mathscr{E}}\left(N-N^{\tau_{n}}\right)_{-}}{} \cdot\left(V \mathscr{E}\left(N-N^{\tau_{n}}\right)\right)-\left(1_{\rrbracket \tau_{n}, \tau_{n+1} \rrbracket} V_{-}\right) \cdot N
\end{aligned}
$$

is a $\sigma$-martingale, which implies the first claim. The second follows as in Lemma 3.14.
2. By Lemma A. 8 we must show that $\left(V+\left[V, N^{\star}\right]\right) Z^{P^{\star}}$ is a $P-\sigma$-martingale. Integration by parts yields

$$
\left(V+\left[V, N^{\star}\right]\right) Z^{P^{\star}} \sim Z_{-}^{P^{\star}} \cdot\left(V+\left[V, N^{\star}\right]+\left[V+\left[V, N^{\star}\right], \frac{1}{1+\Delta A^{K}} \cdot M^{K}\right]\right)
$$

Hence we must show that the integrator is a $\sigma$-martingale. It equals

$$
\begin{aligned}
V+ & {\left[V, N^{\star}+\frac{1}{1+\Delta A^{K}} \cdot M^{K}+\left[N^{\star}, \frac{1}{1+\Delta A^{K}} \cdot M^{K}\right]\right] } \\
=V+[V, N]+[V, \tilde{a} \cdot S-K & +[\tilde{a} \cdot S, K]-\left(\tilde{a}\left(1+\Delta A^{K}\right)\right) \cdot M^{S \star} \\
& \left.+\frac{1}{1+\Delta A^{K}} \cdot M^{K}-\left[\tilde{a} \cdot M^{S \star}, M^{K}\right]\right] .
\end{aligned}
$$

Since $V+[V, N]$ is a $\sigma$-martingale by statement 1 , it suffices to show that the right-hand side of the long covariation term vanishes. To this end, observe that
using (3.27), (3.17) we get

$$
\begin{aligned}
\tilde{a} \cdot S & =\tilde{a} \cdot M^{S \star}+\left(\tilde{a}^{\top} b^{S \star}\right) \cdot A \\
& =\tilde{a} \cdot M^{S \star}+\frac{b^{K}}{1+\Delta A^{K}} \cdot A \\
& =\tilde{a} \cdot M^{S \star}+\frac{1}{1+\Delta A^{K}} \cdot A^{K}
\end{aligned}
$$

and hence

$$
\begin{aligned}
{[\tilde{a} \cdot S, K] } & =\left[\tilde{a} \cdot M^{S \star}, M^{K}\right]+\left[\tilde{a} \cdot M^{S \star}, A^{K}\right]+\left[\frac{1}{1+\Delta A^{K}} \cdot A^{K}, K\right] \\
& =\left[\tilde{a} \cdot M^{S \star}, M^{K}\right]+\left(\tilde{a} \Delta A^{K}\right) \cdot M^{S \star}+\frac{\Delta A^{K}}{1+\Delta A^{K}} \cdot K .
\end{aligned}
$$

This yields the first claim. The second follows again as in Lemma 3.14.

In general we do not know whether $V$ is locally square integrable, or even special, under $P$. Crucially, this integrability holds under $P^{\star}$, which is important for evaluation of the expected squared hedging error in Section 4.

LEmmA 4.4. We have $1 . V^{2} L,(v+\vartheta \cdot S)^{2} L$, and $(v+\vartheta \cdot S-V)^{2} L$ are submartingales for any admissible endowment/strategy pair $(v, \vartheta)$.
2. $V$ is a locally square-integrable semimartingale relative to $P^{\star}$.

Proof. 1. Let $G:=v+\vartheta \cdot S$ and fix $s \leq t$. From Lemmas 3.2(3), 4.1 and Hölder's inequality it follows that

$$
\begin{aligned}
\left(G_{s}-V_{s}\right)^{2} L_{s}^{2} & =\left(\left(G_{s}-V_{s}\right) M_{s}^{(s)}\right)^{2} \\
& =\left(E\left(\left(G_{t}-V_{t}\right) M_{t}^{(s)} \mid \mathscr{F}_{s}\right)\right)^{2} \\
& \leq E\left(\left(G_{t}-V_{t}\right)^{2} L_{t} \mid \mathscr{F}_{s}\right) E\left(\left(1-\lambda^{(s)} \cdot S_{t}\right) M_{t}^{(s)} \mid \mathscr{F}_{s}\right) \\
& =E\left(\left(G_{t}-V_{t}\right)^{2} L_{t} \mid \mathscr{F}_{s}\right) L_{s}
\end{aligned}
$$

Integrability follows by setting $t=T$. The claim for $V^{2} L$ and $G^{2} L$ follows analogously.
2. For any stopping time $\tau$ we have

$$
E_{P^{\star}}\left(V_{\tau}^{2}\right)=E\left(Z_{\tau}^{P \star} V_{\tau}^{2}\right) \leq \frac{E\left(L_{\tau} V_{\tau}^{2}\right)}{E\left(L_{0}\right)} \leq \frac{E\left(H^{2}\right)}{E\left(L_{0}\right)}
$$

by statement 1 , which implies the claim (cf. Lemma A.2).

LEmmA 4.5. Outside some $P \otimes A$-null set we have

$$
\begin{equation*}
b^{V \star}=\tilde{c}^{V S \star} \tilde{a}, \tag{4.6}
\end{equation*}
$$

$$
\begin{equation*}
\tilde{c}^{S \star}\left(\tilde{c}^{S \star}\right)^{-1} \tilde{c}^{S V \star}=\tilde{c}^{S V \star} \tag{4.7}
\end{equation*}
$$

Proof. By Lemma 4.3 and (3.27), (3.17) we have

$$
\begin{aligned}
0 & \sim^{\star} V+\left[V, N^{\star}\right] \\
& =V-\hat{a} \cdot\left[V, M^{S^{\star}}\right] \\
& =V-\hat{a} \cdot[V, S]+\hat{a} \cdot\left[V, A^{S^{\star}}\right] \\
& \sim^{\star} A^{V \star}-\hat{a} \cdot\langle V, S\rangle^{P^{\star}}+\left(\hat{a}^{\top} \Delta A^{S \star}\right) \cdot V \\
& \sim^{\star}\left(b^{V \star}-\tilde{c}^{V S \star} \hat{a}\right) \cdot A+\left(\left(1+\Delta A^{K}\right) \tilde{a}^{\top} b^{S \star} \Delta A\right) \cdot A^{V \star} \\
& =\left(b^{V \star}-\left(1+\Delta A^{K}\right) \tilde{c}^{V S^{\star}} \tilde{a}+\Delta A^{K} b^{V \star}\right) \cdot A \\
& =\left(\left(1+\Delta A^{K}\right)\left(b^{V \star}-\tilde{c}^{V S_{\star}} \tilde{a}\right)\right) \cdot A,
\end{aligned}
$$

which yields the first assertion.
Fix $(\omega, t) \in \Omega \times[0, T]$. Since

$$
\tilde{c}_{t}^{S, V \star}(\omega)=\left(\begin{array}{cc}
\tilde{c}_{t}^{S \star} & \tilde{c}_{t}^{S V \star} \\
\left(\tilde{c}_{t}^{S V \star}\right)^{\top} & \tilde{c}_{t}^{V \star}
\end{array}\right)(\omega)
$$

is a symmetric, nonnegative matrix, we have (4.7) by Albert [1], Theorem 9.1.6.

The next definition constitutes a first step toward optimal hedging.
Definition 4.6. We call the process

$$
\xi:=\left(\tilde{c}^{S \star}\right)^{-1} \tilde{c}^{S V \star}
$$

pure hedge coefficient.
The following representations of $\xi$ establish the link to (1.1).
PROPOSITION 4.7. The pure hedge coefficient $\xi$ satisfies

$$
\begin{equation*}
\xi \cdot\langle S, S\rangle^{P^{\star}}=\langle S, V\rangle^{P^{\star}} \tag{4.8}
\end{equation*}
$$

and

$$
\begin{equation*}
\xi \cdot\left\langle M^{S \star}, M^{S \star}\right\rangle^{P^{\star}}=\left\langle M^{S \star}, M^{V \star}\right\rangle^{P^{\star}} \tag{4.9}
\end{equation*}
$$

In the univariate case, (4.8) and (4.9) can be written more plainly as

$$
\begin{equation*}
\xi_{t}=\frac{d\langle S, V\rangle_{t}^{P^{\star}}}{d\langle S, S\rangle_{t}^{P^{\star}}}=\frac{d\left\langle M^{S \star}, M^{V \star}\right\rangle_{t}^{P^{\star}}}{d\left\langle M^{S \star}, M^{S \star}\right\rangle_{t}^{P^{\star}}} \tag{4.10}
\end{equation*}
$$

Proof. Lemma 4.5 yields

$$
\langle S, V\rangle^{P^{\star}}=\tilde{c}^{S V \star} \cdot A=\left(\tilde{c}^{S \star}\left(\tilde{c}^{S \star}\right)^{-1} \tilde{c}^{S V \star}\right) \cdot A=\xi \cdot\langle S, S\rangle^{P^{\star}} .
$$

By (3.30) and (4.6) we have

$$
\begin{align*}
\left\langle M^{S \star}, M^{S \star}\right\rangle^{P \star} & =\left(\tilde{c}^{S \star}-b^{S \star}\left(b^{S \star}\right)^{\top} \Delta A\right) \cdot A \\
& =\left(\left(1_{d}-\tilde{c}^{S \star} \tilde{a} \tilde{a}^{\top} \Delta A\right) \tilde{c}^{S \star}\right) \cdot A \tag{4.11}
\end{align*}
$$

and

$$
\begin{align*}
\left\langle M^{S \star}, M^{V \star}\right\rangle^{P^{\star}} & =\left(\tilde{c}_{t}^{S V \star}-b_{t}^{S \star} b_{t}^{V \star} \Delta A\right) \cdot A \\
& =\left(\left(1_{d}-\tilde{c}^{S \star} \tilde{a} \tilde{a}^{\top} \Delta A\right) \tilde{c}^{S V \star}\right) \cdot A, \tag{4.12}
\end{align*}
$$

where $1_{d}$ denotes the identity matrix. Equations (4.11), (4.12), (4.7) yield (4.9).

The pure hedge coefficient appears in the following decomposition:
Lemma 4.8. There exists a $P^{\star}$-local martingale $M$ with $M_{0}=0$ that is $P^{\star}$-orthogonal to $M^{S \star}$ (in the sense that $M^{S \star} M$ is a $P^{\star}$-local martingale) and such that

$$
\begin{equation*}
V=V_{0}+\xi \cdot S+M \tag{4.13}
\end{equation*}
$$

holds.
Proof. By (4.7), (4.6), (3.30) we have

$$
0=\left(\tilde{a}^{\top} \tilde{c}^{S V \star}-\tilde{a}^{\top} \tilde{c}^{S \star}\left(\tilde{c}^{S \star}\right)^{-1} \tilde{c}^{S V \star}\right) \cdot A=\left(b^{V \star}-\left(b^{S \star}\right)^{\top} \xi\right) \cdot A,
$$

which implies that $M:=V-V_{0}-\xi \cdot S$ is a $P^{\star}-\sigma$-martingale. By bilinearity and (4.7) the modified second $P^{\star}$-characteristics of $M$ in the sense of (1.2) equals

$$
\begin{aligned}
\tilde{c}^{M \star} & =\tilde{c}^{V \star}-2 \xi^{\top} \tilde{c}^{S V \star}+\xi^{\top} \tilde{c}^{S \star} \xi \\
& =\tilde{c}^{V \star}-2\left(\tilde{c}^{S V \star}\right)^{\top}\left(\tilde{c}^{S \star}\right)^{-1} \tilde{c}^{S V \star}+\left(\tilde{c}^{S V \star}\right)^{\top}\left(\tilde{c}^{S \star}\right)^{-1} \tilde{c}^{S \star}\left(\tilde{c}^{S \star}\right)^{-1} \tilde{c}^{S V \star} \\
& =\tilde{c}^{V \star}-\left(\tilde{c}^{S V \star}\right)^{\top}\left(\tilde{c}^{S \star}\right)^{-1} \tilde{c}^{S V \star} \leq \tilde{c}^{V \star} .
\end{aligned}
$$

Since $V$ is a locally square-integrable semimartingale relative to $P^{\star}$, it follows that $M$ is a locally square-integrable martingale relative to $P^{\star}$ (cf. [28], II.2.29). Since

$$
\left\langle M^{S \star}, M\right\rangle^{P^{\star}}=\langle S, V-\xi \cdot S\rangle^{P^{\star}}=\left(\tilde{c}^{S V \star}-\tilde{c}^{S \star}\left(\tilde{c}^{S \star}\right)^{-1} \tilde{c}^{S V \star}\right) \cdot A=0
$$

by (4.7), we have that $M^{S \star} M$ is a $P^{\star}$-local martingale.
Equation (4.13) can be interpreted as a process version of the $P^{\star}$-FöllmerSchweizer decomposition of $H$. The integrand in the latter yields the locally riskminimizing hedging strategy in the sense of [41] or [18] relative to $P^{\star}$.

### 4.2. Main results.

Lemma 4.9. For any $\mathscr{F}_{0}$-measurable random variable $v$, the feedback equation

$$
\begin{equation*}
\varphi_{t}=\xi_{t}-\left(v+\varphi \cdot S_{t-}-V_{t-}\right) \tilde{a}_{t} \tag{4.14}
\end{equation*}
$$

has a unique solution $\varphi(v):=\varphi \in L(S)$.

Proof. In the proof of Theorem 4.10 below it is shown that $\xi \in L(S)$. The stochastic differential equation

$$
\begin{equation*}
G=\left(\xi-\left(v-V_{-}\right) \tilde{a}\right) \cdot S-G_{-} \cdot(\tilde{a} \cdot S) \tag{4.15}
\end{equation*}
$$

has a unique solution $G$ by Jacod [27], (6.8). If we set $\varphi_{t}:=\xi_{t}-\left(v+G_{t-}-V_{t-}\right) \tilde{a}_{t}$, then $\varphi \in L(S)$ solves (4.14).

If, on the other hand, some $\tilde{\varphi} \in L(S)$ solves (4.14) as well, then $\widetilde{G}:=\tilde{\varphi} \cdot S$ is a solution to (4.15). This implies $\widetilde{G}=G$ and hence $\tilde{\varphi}=\varphi$.

We are now ready to state our first main result.

THEOREM 4.10. 1. The process $\varphi:=\varphi\left(v_{0}\right)$ given by the feedback expression (4.14) is an optimal hedging strategy for initial endowment $v_{0} \in L^{2}\left(\Omega, \mathscr{F}_{0}, P\right)$.
2. $\left(v_{0}, \varphi\right):=\left(V_{0}, \varphi\left(V_{0}\right)\right)$ is an optimal endowment/strategy pair.

Proof. 1. Denote by

$$
\left(\left(\begin{array}{l}
b^{S} \\
b^{V} \\
b^{K}
\end{array}\right),\left(\begin{array}{ccc}
c^{S} & c^{S V} & c^{S K} \\
c^{V S} & c^{V} & c^{V K} \\
c^{K S} & c^{K V} & c^{K}
\end{array}\right), F^{S, V, K}, A\right)
$$

$P$-differential characteristics of $(S, V, K)$ relative to the "truncation" function $h(x, z, y):=\left(x, z 1_{\{|z| \leq 1\}}, y\right)$ on $\mathbb{R}^{d} \times \mathbb{R} \times \mathbb{R}$. [Should $V$ be a $P$-special semimartingale, we could also choose the identity $h(x, z, y)=(x, z, y)$ as usual in this paper.] Along the same lines as in the proof of Lemma 3.17 it follows that

$$
\begin{align*}
\tilde{c}^{S V \star} & =\frac{1}{1+\Delta A^{K}}\left(c^{S V}+\int x z(1+y) F^{S, V, K}(d(x, z, y))\right),  \tag{4.16}\\
\tilde{c}^{V \star} & =\frac{1}{1+\Delta A^{K}}\left(c^{V}+\int z^{2}(1+y) F^{S, V, K}(d(x, z, y))\right) . \tag{4.17}
\end{align*}
$$

Let $\bar{\varphi}$ be an optimal hedging strategy for initial endowment $v_{0}$, denote by $G:=$ $v_{0}+\bar{\varphi} \cdot S$ its value process, and set $\bar{\xi}:=\bar{\varphi}+\left(G_{-}-V_{-}\right) \tilde{a}$. Moreover, let $\vartheta \in \bar{\Theta}$ and $\widetilde{G}:=\vartheta \cdot S$. In the proof of Lemma 4.1 it is shown that $Z \widetilde{G}$ is a martingale for
$Z:=(G-V) L$. Integration by parts yields $Z \widetilde{G}=L_{0} \mathscr{E}(K)(G-V) \widetilde{G}=L_{-} \cdot U$ with

$$
\begin{aligned}
U= & (G-V) \widetilde{G}+((G-V) \widetilde{G})_{-} \cdot K+[(G-V) \widetilde{G}, K] \\
= & (G-V)_{-} \cdot \widetilde{G}+\widetilde{G}_{-} \cdot(G-V)+[G-V, \widetilde{G}]+((G-V) \widetilde{G})_{-} \cdot K \\
& +(G-V)_{-} \cdot[\widetilde{G}, K]+\widetilde{G}_{-} \cdot[G-V, K]+[G-V,[\widetilde{G}, K]] \\
= & (G-V)_{-} \cdot\left(\widetilde{G}_{-} \cdot N+\vartheta \cdot(S+[S, N])\right) \\
& +\bar{\xi} \cdot\left(\widetilde{G}_{-} \cdot(S+[S, K])+[\widetilde{G}, S+[S, K]]\right) \\
& -\widetilde{G}_{-} \cdot(V+[V, K])-[\widetilde{G}, V+[V, K]] .
\end{aligned}
$$

The first term on the right-hand side is a $\sigma$-martingale by Lemma 3.14. By (3.15), the remaining two terms equal

$$
\begin{aligned}
G_{-} \cdot & (\bar{\xi} \cdot(S+[S, N])-(V+[V, N])) \\
& +\left(\tilde{a} G_{-}+\vartheta\right) \cdot(\bar{\xi} \cdot[S, S+[S, K]]-[V, S+[S, K]])
\end{aligned}
$$

The first line is a $\sigma$-martingale by Lemmas 3.14 and 4.3. By (3.26), (3.18) we have

$$
\begin{align*}
{[S, S+[S, K]] } & =\left[S^{c}, S^{c}\right]+\int x x^{\top}(1+y) \mu^{S, K}(d(x, y)) \\
& \sim\left(c^{S}+\int x x^{\top}(1+y) F^{S, K}(d(x, y))\right) \cdot A  \tag{4.18}\\
& =\bar{c} \cdot A \\
& =\left(\left(1+\Delta A^{K}\right) \tilde{c}^{S \star}\right) \cdot A .
\end{align*}
$$

Similarly, (4.16) yields

$$
\begin{align*}
{[S+[S, K], V] } & \sim\left(c^{S V}+\int x z(1+y) F^{S, V, K}(d(x, z, y))\right) \cdot A \\
& =\left(\left(1+\Delta A^{K}\right) \tilde{c}^{S V \star}\right) \cdot A \tag{4.19}
\end{align*}
$$

For later use, we observe that

$$
\begin{equation*}
[V+[V, K], V] \sim\left(\left(1+\Delta A^{K}\right) \tilde{c}^{V \star}\right) \cdot A 1 \tag{4.20}
\end{equation*}
$$

by (4.17). Altogether, we have that $\left(\left(\tilde{a} G_{-}+\vartheta\right)\left(1+\Delta A^{K}\right)\left(\tilde{c}^{S \star} \bar{\xi}-\tilde{c}^{S V \star}\right)\right) \cdot A$ is a $\sigma$-martingale. This being true for any $\vartheta$, we have

$$
\begin{equation*}
\tilde{c}^{S \star \bar{\xi}}-\tilde{c}^{S V \star}=0 \tag{4.21}
\end{equation*}
$$

$P \otimes A$-almost everywhere.

For $n \in \mathbb{N}$ define the predictable set $D_{n}:=\{|\xi| \vee|\bar{\xi}| \leq n\}$. Corollary 3.20 and (4.21), (4.7) yield

$$
\left((\bar{\xi}-\xi) 1_{D_{n}}\right) \cdot A^{S \star}=\left(\left(\bar{\xi}-\left(\tilde{c}^{S \star}\right)^{-1} \tilde{c}^{S V \star}\right)^{\top} \tilde{c}^{S \star} \tilde{a} 1_{D_{n}}\right) \cdot A=0
$$

as well as

$$
\begin{aligned}
& \left\langle\left((\bar{\xi}-\xi) 1_{D_{n}}\right) \cdot M^{S \star},\left((\bar{\xi}-\xi) 1_{D_{n}}\right) \cdot M^{S \star}\right)^{P \star} \\
& \quad=\left((\bar{\xi}-\xi)^{\top} \hat{c}^{S \star}(\bar{\xi}-\xi) 1_{D_{n}}\right) \cdot A \\
& \quad \leq\left((\bar{\xi}-\xi)^{\top} \tilde{c}^{S \star}(\bar{\xi}-\xi) 1_{D_{n}}\right) \cdot A=0 .
\end{aligned}
$$

Consequently, $\left((\bar{\xi}-\xi) 1_{D_{n}}\right) \cdot S=0$ for any $n$, which in turn implies $\bar{\xi}-\xi \in L(S)$ and $(\bar{\xi}-\xi) \cdot S=0$ by Lemma A.11. In particular, we have $\xi=\bar{\xi}-(\bar{\xi}-\xi) \in$ $L(S)$. The proof of Lemma 4.9 yields that $\varphi \cdot S=\bar{\varphi} \cdot S$ as well. In particular, $\varphi$ is admissible and an optimal hedging strategy for initial endowment $v_{0}$.
2. This follows essentially as statement 1 . We only have to determine the optimal initial endowment. Denote by ( $v_{0}, \bar{\varphi}$ ) an optimal endowment/strategy pair and let $Z$ be as in the first part of the proof. Parallel to (4.4), we conclude that $E\left(v Z_{T}\right)=0$ for any $v \in L^{2}\left(\Omega, \mathscr{F}_{0}, P\right)$, which implies $0=E\left(Z_{T} \mid \mathscr{F}_{0}\right)=Z_{0}=L_{0}\left(v_{0}-V_{0}\right)$. Consequently, $v_{0}=V_{0}$ as claimed.

As is well known, the gains process $\varphi \cdot S$ can be expressed more explicitly.
Corollary 4.11. The gains process of the optimal hedge in Theorem 4.10 equals

$$
\varphi \cdot S=\mathscr{E}(-\tilde{a} \cdot S)\left(\frac{\xi+\left(V_{-}-v_{0}\right) \tilde{a}}{\mathscr{E}(-\tilde{a} \cdot S)_{-}} \cdot\left(S+\frac{\tilde{a}}{1-\tilde{a}^{\top} \Delta S} \cdot[S, S]\right)\right)
$$

unless $\mathscr{E}(-\tilde{a} \cdot S)$ jumps to 0 .
Proof. By [27], (6.8) the stochastic differential equation $X=Y+X_{-} \cdot Z$ with two semimartingales $Y, Z$ such that $Y_{0}=0$ is uniquely solved by

$$
X=\mathscr{E}(Z)\left(\frac{1}{\mathscr{E}(Z)_{-}} \cdot Y-\frac{1}{\mathscr{E}(Z)} \cdot[Y, Z]\right)
$$

unless $\mathscr{E}(Z)$ jumps to 0 . Since

$$
\varphi \cdot S=\left(\xi-\left(v_{0}-V_{-}\right) \tilde{a}\right) \cdot S+(\varphi \cdot S)_{-} \cdot(-\tilde{a} \cdot S)
$$

the assertion follows.

Finally, we state formulas for the hedging error.

THEOREM 4.12. $\quad$ The expected squared hedging error of the optimal hedge in Theorem 4.10 equals

$$
\begin{align*}
E\left(\left(v_{0}\right.\right. & \left.\left.+\varphi \cdot S_{T}-H\right)^{2}\right) \\
= & E\left(\left(v_{0}-V_{0}\right)^{2} L_{0}+\left(\left(\tilde{c}^{V \star}-\left(\tilde{c}^{S V \star}\right)^{\top}\left(\tilde{c}^{S \star}\right)^{-1} \tilde{c}^{S V \star}\right) L\right) \cdot A_{T}\right) \\
= & E\left(\left(v_{0}-V_{0}\right)^{2} L_{0}+L \cdot\left(\langle V, V\rangle^{P^{\star}}-\xi \cdot\langle V, S\rangle^{P^{\star}}\right)_{T}\right) \\
= & E\left(\left(v_{0}-V_{0}\right)^{2} L_{0}+L \cdot\langle V-\xi \cdot S, V-\xi \cdot S\rangle_{T}^{P^{\star}}\right)  \tag{4.22}\\
= & E_{P^{\star}}\left(\left(v_{0}-V_{0}\right)^{2}\right. \\
& \left.\quad+\left(\left(\tilde{c}^{V \star}-\left(\tilde{c}^{S V \star}\right)^{\top}\left(\tilde{c}^{S \star}\right)^{-1} \tilde{c}^{S V \star}\right) \mathscr{E}\left(A^{K}\right)\right) \cdot A_{T}\right) E\left(L_{0}\right) \\
= & E_{P^{\star}}\left(\left(v_{0}-V_{0}\right)^{2}+\mathscr{E}\left(A^{K}\right) \cdot\langle V-\xi \cdot S, V-\xi \cdot S\rangle_{T}^{P^{\star}}\right) E\left(L_{0}\right) . \tag{4.23}
\end{align*}
$$

Proof. In view of Proposition 1.2, the second equality is obvious. The third and the last follow from

$$
\langle S, V-\xi \cdot S\rangle^{P^{\star}}=\left(\tilde{c}^{S V \star}-\tilde{c}^{S \star}\left(\tilde{c}^{S \star}\right)^{-1} \tilde{c}^{S V \star}\right) \cdot A=0
$$

Define $G:=v_{0}+\varphi \cdot S$ and $Z:=(G-V) L$ as in the proof of Theorem 4.10. Since $(G-V)^{2} L$ is a submartingale by Lemma 4.4, there exists a unique increasing predictable process $B$ with $B_{0}=0$ and such that $(G-V)^{2} L-B$ is a martingale. Since

$$
E\left(\left(v_{0}+\varphi \cdot S_{T}-H\right)^{2}\right)=E\left(\left(G_{T}-V_{T}\right)^{2} L_{T}\right)=E\left(\left(G_{0}-V_{0}\right)^{2} L_{0}\right)+E\left(B_{T}\right)
$$

the first equality in the theorem holds if

$$
\begin{equation*}
B=\left(\left(\tilde{c}^{V \star}-\xi^{\top} \tilde{c}^{S V \star}\right) L\right) \cdot A \tag{4.24}
\end{equation*}
$$

Since $G Z$ and $Z$ are martingales, we have

$$
\begin{align*}
-(G- & V)^{2} L \sim V Z \\
\sim & Z_{-} \cdot V+[V, Z]  \tag{4.25}\\
= & \left((G-V)_{-} L_{-}\right) \cdot V \\
& +\left[V,(G-V)_{-} \cdot L+L_{-} \cdot(G-V)+[G-V, L]\right]
\end{align*}
$$

In view of

$$
G-V=v_{0}+\xi \cdot S-\left((G-V)_{-} \tilde{a}\right) \cdot S-V
$$

(4.25) equals

$$
\begin{aligned}
&\left((G-V)_{-} L_{-}\right) \cdot(V+[V, K-\tilde{a} \cdot S-[\tilde{a} \cdot S, K]]) \\
& \quad+L_{-} \cdot {[V+[V, K], \xi \cdot S-V] }
\end{aligned}
$$

By Lemma 4.3 the first term is a $\sigma$-martingale and hence

$$
\begin{aligned}
(G-V)^{2} L & \sim-L_{-} \cdot(\xi \cdot[V+[V, K], S]-[V+[V, K], V]) \\
& =-L_{-} \cdot(\xi \cdot[V, S+[S, K]]-[V+[V, K], V]) \\
& \sim\left(L_{-}\left(1+\Delta A^{K}\right)\left(\tilde{c}^{V \star}-\xi^{\top} \tilde{c}^{S V \star}\right)\right) \cdot A \\
& =\left(\left(L-L_{-} \Delta M^{K}\right)\left(\tilde{c}^{V \star}-\xi^{\top} \tilde{c}^{S V \star}\right)\right) \cdot A
\end{aligned}
$$

by (4.19) and (4.20). Since $\Delta M^{K} \cdot U=\left[M^{K}, U\right]=\Delta U \cdot M^{K}$ is a $\sigma$-martingale for any predictable process $U$ of finite variation (cf. [28], I.4.49), we obtain

$$
(G-V)^{2} L \sim\left(L\left(\tilde{c}^{V \star}-\xi^{\top} \tilde{c}^{S V \star}\right)\right) \cdot A .
$$

Therefore the difference of both sides of (4.24) is a predictable $\sigma$-martingale of finite variation and hence 0 .

It remains to be shown that (4.23) equals (4.22). Integration by parts yields

$$
\begin{aligned}
Z^{P^{\star}}( & \left.E\left(L_{0}\right) \mathscr{E}\left(A^{K}\right) \cdot\langle V-\xi \cdot S, V-\xi \cdot S\rangle^{P^{\star}}\right) \\
= & \left(Z^{P^{\star}} E\left(L_{0}\right) \mathscr{E}\left(A^{K}\right)\right) \cdot\langle V-\xi \cdot S, V-\xi \cdot S\rangle^{P^{\star}} \\
& +\left(E\left(L_{0}\right) \mathscr{E}\left(A^{K}\right) \cdot\langle V-\xi \cdot S, V-\xi \cdot S\rangle^{P^{\star}}\right)_{-} \cdot Z^{P^{\star}} \\
= & L \cdot\langle V-\xi \cdot S, V-\xi \cdot S\rangle P^{\star}+M
\end{aligned}
$$

with some $P$-local martingale $M$. Hence

$$
\begin{aligned}
E_{P^{\star}} & \left(\mathscr{E}\left(A^{K}\right) \cdot\langle V-\xi \cdot S, V-\xi \cdot S\rangle_{T_{n}}^{P^{\star}}\right) E\left(L_{0}\right) \\
& =E\left(Z_{T_{n}}^{P^{\star}}\left(E\left(L_{0}\right) \mathscr{E}\left(A^{K}\right) \cdot\langle V-\xi \cdot S, V-\xi \cdot S\rangle_{T_{n}}^{P^{\star}}\right)\right) \\
& =E\left(L \cdot\langle V-\xi \cdot S, V-\xi \cdot S\rangle_{T_{n}}^{P^{\star}}\right)
\end{aligned}
$$

where $\left(T_{n}\right)_{n \in \mathbb{N}}$ denotes a localizing sequence for $M$. Monotone convergence yields that (4.23) equals (4.22).

If the results in this paper are to be applied to concrete models, it is not necessary to determine all the processes that have been introduced. Instead, one may proceed as follows: first one determines the opportunity process $L$ and the adjustment process $\tilde{a}$ using the characterization in Theorem 3.25. These processes yield the modified mean-variance tradeoff process $K$, the opportunity-neutral measure $P^{\star}$ and the variance-optimal logarithm process $N$. Finally, the mean-value process $V$ leads to the pure hedge coefficient $\xi$ and hence to the optimal hedge $\varphi$.
4.3. Connections to the literature. In this section we clarify the link of our results to the literature. If $S$ is a martingale, we are in the setup of Föllmer and Sondermann [19]. In our notation, they show that the optimal hedge $\varphi$ satisfies

$$
\begin{equation*}
\varphi_{t}=\frac{d\langle S, V\rangle_{t}}{d\langle S, S\rangle_{t}} \tag{4.26}
\end{equation*}
$$

where

$$
\begin{equation*}
V_{t}=E\left(H \mid \mathscr{F}_{t}\right) \tag{4.27}
\end{equation*}
$$

Applying our results to the martingale case, one immediately verifies that $L=1$, $\tilde{a}=0, K=0, N=0, Q^{\star}=P^{\star}=P$. Consequently, equation (4.2) for the meanvalue process of the option reduces to (4.27). Moreover, the optimal hedge $\varphi$ coincides with the pure hedge $\xi$, which satisfies $\xi \cdot\langle S, S\rangle=\langle S, V\rangle$ in accordance with (4.26).

Schweizer [42] goes beyond the martingale case. He shows that if the MVT process $\widehat{K}$ is deterministic, then the optimal hedging strategy for initial endowment $v_{0}$ contains a feedback element and is of the form

$$
\begin{equation*}
\varphi_{t}=\xi_{t}-\left(v_{0}+\varphi \cdot S_{t-}-V_{t-}\right) \tilde{\lambda}_{t} \tag{4.28}
\end{equation*}
$$

with $\tilde{\lambda}$ from (3.35). Here, the pure hedge coefficient $\xi$ is the integrand in the Föllmer-Schweizer decomposition of the claim, that is,

$$
H=V_{0}+\xi \cdot S_{T}+R_{T}
$$

where $V_{0}$ is a $\mathscr{F}_{0}$-measurable random variable and $R$ denotes a martingale which is orthogonal to $M^{S}$ (in the sense that $M^{S} R$ is a local martingale). In order to express the pure hedge coefficient similarly as in (4.26), recall that the minimal signed martingale measure $Q$ is given by

$$
\frac{d Q}{d P}:=\mathscr{E}\left(-\hat{\lambda} \cdot M^{S}\right)_{T}
$$

If we define $V$ as " $Q$-conditional expectation" of $H$ in the sense of

$$
\begin{equation*}
V_{t}:=E\left(H \mathscr{E}\left(\left(-\hat{\lambda} 1_{\rrbracket t, T \rrbracket}\right) \cdot M^{S}\right)_{T} \mid \mathscr{F}_{t}\right) \tag{4.29}
\end{equation*}
$$

then the pure hedge coefficient can be written as

$$
\begin{equation*}
\xi_{t}=\frac{d\langle S, V\rangle_{t}}{d\langle S, S\rangle_{t}}=\frac{d\left\langle M^{S}, M^{V}\right\rangle_{t}}{d\left\langle M^{S}, M^{S}\right\rangle_{t}} \tag{4.30}
\end{equation*}
$$

The hedging error satisfies the equation

$$
\begin{align*}
& E\left(\left(v_{0}+\varphi \cdot S_{T}-H\right)^{2}\right)  \tag{4.31}\\
& \quad=E\left(\left(v_{0}-V_{0}\right)^{2}+\mathscr{E}(\widehat{K}) \cdot\langle V-\xi \cdot S, V-\xi \cdot S\rangle_{T}\right) \frac{1}{\mathscr{E}(\widehat{K})_{T}}
\end{align*}
$$

In these formulas, all predictable covariation processes refer to the original probability measure $P$.

It is easy to see that (4.28)-(4.31) are special cases of our general results. To this end, recall that $L=\mathscr{E}(\widehat{K}) / \mathscr{E}(\widehat{K})_{T}, P^{\star}=P$, and $\tilde{a}=\tilde{\lambda}$ in the case of deterministic MVT (cf. Corollaries 3.31 and 3.30). Hence

$$
N^{\star}=-\left(\left(1+\Delta A^{K}\right) \tilde{a}\right) \cdot M^{S \star}=-((1+\Delta \widehat{K}) \tilde{\lambda}) \cdot M^{S}=-\hat{\lambda} \cdot M^{S}
$$

Consequently, (4.28), (4.29), (4.30), (4.31) correspond to (4.14), (4.3), (4.10), (4.23), respectively.

If the MVT process fails to be deterministic, the above formulas do not lead to the optimal hedge any more. Following Hipp [22], [44] observes that a key role in the general case is played by the variance optimal signed martingale measure $Q^{\star}$ and the adjustment process $\tilde{a}$. Schweizer characterizes both the adjustment process and the optimal hedging strategy in terms of backward stochastic differential equations. The use of these BSDEs in practice is complicated by their involved boundary conditions, which themselves depend on the unknown solution.

Rheinländer and Schweizer [40] show that the optimal hedging strategy $\varphi$ satisfies similar equations as in the case of deterministic MVT if $S$ is continuous. More specifically, it is of feedback form

$$
\varphi_{t}=\xi_{t}-\left(v_{0}+\varphi \cdot S_{t-}-V_{t-}\right) \tilde{a}_{t}
$$

where $V_{t}:=E_{Q^{\star}}\left(H \mid \mathscr{F}_{t}\right)$ is the martingale generated by $H$ relative to the varianceoptimal $\mathrm{S} \sigma \mathrm{MM} Q^{\star}$ and the pure hedge coefficient $\xi$ is the integrand in the Galtchouk-Kunita-Watanabe decomposition of $H$ relative to $Q^{\star}$ rather than $P$, that is,

$$
\xi_{t}=\frac{d\langle S, V\rangle_{t}^{Q^{\star}}}{d\langle S, S\rangle_{t}^{Q^{\star}}}
$$

This equation corresponds to our expression (4.10) because the predictable covariation does not depend on the probability measure for continuous processes.

An alternative approach in the continuous case is pursued by Gourieroux, Laurent and Pham [20] who use a new numeraire $\mathscr{E}(-\tilde{a} \cdot S)$ combined with a change of measure to transform the original semimartingale problem to a martingale problem à la Föllmer and Sondermann [19]. The task of computing $\tilde{a}$ has become a separate issue in the literature. It is tackled in a number of diffusion or jump-diffusion settings, for example, by Laurent and Pham [30], Biagini, Guasoni and Pratelli [6], Biagini and Guasoni [5], Hobson [24]. Our characterization of the adjustment process in Theorem 3.25 appears to be more suitable for direct computations than the methods available to date (cf. [11]).

The literature on discontinuous processes is more limited. Two partial results are reported by Arai [3] and Lim [33]. Arai extends the numeraire method of Gourieroux, Laurent and Pham [20] to discontinuous semimartingales assuming that $Q^{\star}$ is equivalent to $P$ and shows that $V$ in (4.3) is a $Q^{\star}$-martingale. However, Arai's results are hard to use for explicit computations since he does not provide a method for obtaining $\tilde{a}$.

Lim [33] uses BSDEs to compute the optimal hedge in a jump diffusion setting where asset price characteristics are adapted to a Brownian filtration. In addition he requires a certain martingale invariance property. He characterizes the optimal hedge explicitly at the cost of a somewhat restrictive model setup.

Finally, we want to explain another close link of our results to the formulas (4.28-4.31) of [42]. We already observed in Lemma 3.23 that the variance-optimal $\mathrm{S} \sigma \mathrm{MM} Q^{\star}$ is the minimal $\mathrm{S} \sigma \mathrm{MM}$ relative to $P^{\star}$. Moreover, $\tilde{a}$ and $\hat{a}$ coincide with the processes $\tilde{\lambda}$ and $\hat{\lambda}$ in [42] or Section 3.6 relative to $P^{\star}$ instead of $P$. Consequently, equations (4.14), (4.3), (4.10) are $P^{\star}$-versions of the formulas (4.28), (4.29), (4.30). The change of measure $P \rightarrow P^{\star}$ neutralizes the effect of stochastic mean-variance tradeoff which makes the results in [42] break down. With the hedging error one has to be slightly more careful. Since $\widehat{K}^{\star}=A^{K}$, we can view (4.23) essentially as a $P^{\star}$-version of (4.31). We only have to replace the deterministic second factor $1 / \mathscr{E}(\widehat{K})_{T}$ by

$$
E\left(\frac{1}{\mathscr{E}\left(\widehat{K}^{\star}\right)_{T}}\right)=E\left(\frac{1}{\mathscr{E}\left(A^{K}\right)_{T}}\right)=E_{P^{\star}}\left(\frac{E\left(L_{0}\right)}{L_{T}}\right)=E\left(L_{0}\right)
$$

APPENDIX

## A.1. Locally square-integrable semimartingales.

Definition A.1. For any special semimartingale $X$ we define

$$
\|X\|_{\mathscr{S}^{2}}:=\left\|X_{0}\right\|_{2}+\left\|\sqrt{\left[M^{X}, M^{X}\right]_{T}}\right\|_{2}+\left\|\operatorname{var}\left(A^{X}\right)_{T}\right\|_{2}
$$

where $\operatorname{var}\left(A^{X}\right)$ denotes the variation process of $A^{X}$ and $\|\cdot\|_{2}$ the $L^{2}$-norm. $X$ is said to belong to the set $\mathscr{S}^{2}$ of square-integrable semimartingales if $\|X\|_{\mathscr{S}^{2}}<\infty$. The elements of the corresponding localized class $\mathscr{S}_{\text {loc }}^{2}$ are called locally squareintegrable semimartingales.

Lemma A.2. For any semimartingale $X$, we have equivalence between:

1. $X \in \mathscr{S}_{\text {loc }}^{2}$.
2. $X_{0} \in L^{2}(P)$ and $X$ is a locally square-integrable semimartingale in the sense of [28], II.2.27, that is, it is a special semimartingale whose local martingale part is locally square-integrable.
3. $X$ is locally in $L^{2}$ in the sense of [15], that is, it belongs locally to the class of processes $Y$ with

$$
\sup \left\{E\left(Y_{\tau}^{2}\right): \tau \text { finite stopping time }\right\}<\infty .
$$

4. $X$ belongs locally to the class of processes $Y$ satisfying $E\left(Y_{\tau}^{2}\right)<\infty$ for any finite stopping time $\tau$.

Proof. We refer to the time set $\mathbb{R}_{+}$rather than $[0, T]$ in this proof.
$1 \Rightarrow 2$ : This follows from [28], II.2.28 and from the inequality

$$
E\left(\sup _{t \in \mathbb{R}_{+}}\left(Y_{t}-Y_{0}\right)^{2}\right) \leq 8\|Y\|_{\mathscr{S}^{2}}^{2}
$$

which holds for any semimartingale $Y$ (cf. [39], Theorem IV.5).
$2 \Rightarrow 3$ : This follows immediately from [28], II.2.28.
$3 \Rightarrow 4$ : This is trivial.
$4 \Rightarrow 1$ : Define a sequence of stopping times $\tau_{n}:=\inf \left\{t \in \mathbb{R}_{+}:\left|X_{t}\right|>n\right\} \wedge n$. Since $\sup _{t \in \mathbb{R}_{+}}\left|X_{t}^{\tau_{n}}\right| \leq n+\left|X_{\tau_{n}}\right|$ is integrable, $X$ is a special semimartingale (cf. [28], I.4.23). Choose a localizing sequence $\left(\sigma_{n}\right)_{n \in \mathbb{N}}$ for the locally bounded process $\operatorname{var}\left(A^{X}\right)$. Then

$$
\begin{aligned}
\sup _{t \in \mathbb{R}_{+}}\left|\left(M^{X}\right)_{t}^{\sigma_{n} \wedge \tau_{n}}\right|^{2} & \leq 3 \sup _{t \in \mathbb{R}_{+}}\left|X_{t}^{\tau_{n}}\right|^{2}+3 \sup _{t \in \mathbb{R}_{+}}\left|\left(A^{X}\right)_{t}^{\sigma_{n}}\right|^{2}+3\left|X_{0}\right|^{2} \\
& \leq 6 n^{2}+6\left|X_{\tau_{n}}\right|^{2}+3 \sup _{t \in \mathbb{R}_{+}}\left(\operatorname{var}\left(A^{X}\right)_{t}^{\sigma_{n}}\right)^{2}+3\left|X_{0}\right|^{2}
\end{aligned}
$$

is integrable for any $n \in \mathbb{N}$, which yields $X \in \mathscr{S}_{\text {loc }}^{2}$ (cf. [28], I.4.50c).
The following result on square integrability of exponential semimartingales is needed in the proof of Proposition 3.28. It extends a parallel statement for local martingales in [27], (8.27).

Lemma A.3. Let $X$ be a locally square-integrable semimartingale such that $\left\langle M^{X}, M^{X}\right\rangle$ and the variation process $\operatorname{var}\left(A^{X}\right)$ are bounded. Then

$$
E\left(\sup _{t} \mathscr{E}(X)_{t}^{2}\right)<\infty
$$

Proof. For ease of notation we prove the assertion for the time set $\mathbb{R}_{+}$rather than $[0, T]$. Denote by $m \in \mathbb{R}_{+}$an upper bound of $V:=\left\langle M^{X}, M^{X}\right\rangle+\operatorname{var}\left(A^{X}\right)$. We write $Z:=\mathscr{E}(X)$ and $Y_{t}^{*}:=\sup _{s \in[0, t]}\left|Y_{S}\right|$ for any process $Y$. For $n \in \mathbb{N}$ define stopping times

$$
\sigma_{n}:=\inf \left\{t \in \mathbb{R}_{+}:\left|Z_{t}\right| \geq n\right\} .
$$

Fix $n$ and set $\widetilde{Z}:=Z^{\sigma_{n}}$.
Step 1: We show that

$$
E\left(\left(\widetilde{Z}_{\tau_{-}}^{*}\right)^{2}\right) \leq 3+(12+3 m) E\left(\left(\widetilde{Z}_{-}^{2} \wedge n^{2}\right) \cdot V_{\tau-}\right)
$$

for any predictable stopping time $\tau$.
In view of $\widetilde{Z}=1+\left(\widetilde{Z}_{-} 1_{\llbracket 0, \sigma_{n} \rrbracket}\right) \cdot M^{X}+\left(\widetilde{Z}_{-} 1_{\llbracket 0, \sigma_{n} \rrbracket}\right) \cdot A^{X}$, we have

$$
E\left(\left(\widetilde{Z}_{\tau_{-}}^{*}\right)^{2}\right) \leq 3+3 E\left(\left(\left(\left(\widetilde{Z}_{-} 1_{\llbracket 0, \sigma_{n} \rrbracket}\right) \cdot M^{X}\right)_{\tau_{-}}^{*}\right)^{2}\right)+3 E\left(\left(\left(\left(\widetilde{Z}_{-} 1_{\llbracket 0, \sigma_{n} \rrbracket}\right) \cdot A^{X}\right)_{\tau_{-}}^{*}\right)^{2}\right)
$$

Since $\tau$ is predictable, Doob's inequality yields

$$
\begin{aligned}
E\left(\left(\left(\left(\widetilde{Z}_{-} 1_{\llbracket 0, \sigma_{n} \rrbracket}\right) \cdot M^{X}\right)_{\tau_{-}}^{*}\right)^{2}\right) & \leq 4 E\left(\left(\widetilde{Z}_{-}^{2} 1_{\llbracket 0, \sigma_{n} \rrbracket} \rrbracket\right) \cdot\left\langle M^{X}, M^{X}\right\rangle_{\tau_{-}}\right) \\
& \leq 4 E\left(\left(\widetilde{Z}_{-}^{2} \wedge n^{2}\right) \cdot V_{\tau_{-}}\right) .
\end{aligned}
$$

For the part of finite variation we have

$$
\begin{aligned}
\left(\left(\left(\widetilde{Z}_{-} 1_{\llbracket 0, \sigma_{n} \rrbracket}\right) \cdot A^{X}\right)_{\tau_{-}}^{*}\right)^{2} & \leq\left(\left(\left|\widetilde{Z}_{-}\right| \wedge n\right) \cdot \operatorname{var}\left(A^{X}\right)_{\tau_{-}}\right)^{2} \\
& \leq\left(\widetilde{Z}_{-}^{2} \wedge n^{2}\right) \cdot \operatorname{var}\left(A^{X}\right)_{\tau_{-}} \operatorname{var}\left(A^{X}\right)_{\infty}
\end{aligned}
$$

and hence

$$
E\left(\left(\left(\left(\widetilde{Z}_{-} 1_{\llbracket 0, \sigma_{n} \rrbracket}\right) \cdot A^{X}\right)_{\tau_{-}}^{*}\right)^{2}\right) \leq m E\left(\left(\widetilde{Z}_{-}^{2} \wedge n^{2}\right) \cdot V_{\tau_{-}}\right)
$$

Step 2: For $\vartheta \in \mathbb{R}_{+}$define the predictable stopping time $T_{\vartheta}:=\inf \left\{t \in \mathbb{R}_{+}: V_{t} \geq\right.$ $\vartheta\}$ (cf. [28], I.2.13). Step 1 yields that

$$
f(\vartheta):=E\left(\left(\widetilde{Z}_{T_{\vartheta}-}^{*} \wedge n\right)^{2}\right) \leq 3+(12+3 m) E\left(\left(\widetilde{Z}_{-}^{2} \wedge n^{2}\right) \cdot V_{T_{\vartheta}-}\right) .
$$

Since $\vartheta \mapsto T_{\vartheta}$ is the pathwise generalized inverse of $V$, we have

$$
\left(\widetilde{Z}_{-}^{2} \wedge n^{2}\right) \cdot V_{T_{\vartheta}-}=\int_{0}^{V_{T_{\vartheta}-}}\left(\widetilde{Z}_{T_{Q^{-}}-}^{2} \wedge n^{2}\right) d \varrho \leq \int_{0}^{\vartheta}\left(\widetilde{Z}_{T_{Q^{-}}}^{*} \wedge n\right)^{2} d \varrho
$$

and hence

$$
f(\vartheta) \leq 3+(12+3 m) \int_{0}^{\vartheta} f(\varrho) d \varrho
$$

for any $\vartheta \in \mathbb{R}_{+}$. By Gronwall's inequality this implies $f(\vartheta) \leq 3 e^{(12+3 m) \vartheta}$. Since $T_{m+1}=\infty$, we have

$$
E\left(n^{2} \wedge \sup _{t \leq \sigma_{n}} Z_{t}^{2}\right)=E\left(\left(\widetilde{Z}_{\infty-}^{*} \wedge n\right)^{2}\right) \leq 3 e^{(12+3 m)(m+1)}
$$

The assertion follows now from monotone convergence.
A.2. $\sigma$-martingales. The following facts on $\sigma$-martingales and integrability can be found, for example, in [29]. We summarize them here for the convenience of the reader.

Definition A.4. A semimartingale $X$ is called $\sigma$-martingale if there exists an increasing sequence $\left(D_{n}\right)_{n \in \mathbb{N}}$ of predictable sets such that $D_{n} \uparrow \Omega \times \mathbb{R}_{+}$up to an evanescent set and $1_{D_{n}} \cdot X$ is a uniformly integrable martingale for any $n \in \mathbb{N}$.

REMARK A.5. Uniformly integrable martingale can be replaced by local martingale in the previous definition.

Lemma A.6. Let $X$ be a semimartingale with differential characteristics $(b, c, F, A)$ relative to some truncation function $h$. Then $X$ is a $\sigma$-martingale if and only if $\int_{\{|x|>1\}}|x| F(d x)<\infty$ and

$$
b+\int(x-h(x)) F(d x)=0
$$

hold outside some $P \otimes A$-null set.

Lemma A.7. $\quad X$ is a uniformly integrable martingale if and only if it is a $\sigma$-martingale of class $(D)$.

Lemma A.8. Let $P^{\star} \sim P$ be a probability measure with density process $Z$. A real-valued semimartingale $X$ is a $P^{\star}$ - $\sigma$-martingale if and only if $X Z$ is a $P$ - $\sigma$-martingale.

Lemma A.9. Let $X$ be $a \mathbb{R}^{d}$-valued semimartingale and let $P^{\star} \sim P$ be a probability measure with density process $Z=Z_{0} \mathscr{E}(N)$. Denote by

$$
\left(b^{X, N}, c^{X, N}, F^{X, N}, A\right)=\left(\binom{b^{X}}{b^{N}},\left(\begin{array}{cc}
c^{X} & c^{X N} \\
c^{N X} & c^{N}
\end{array}\right), F^{X, N}, A\right)
$$

differential characteristics of the $\mathbb{R}^{d+1}$-valued seminartingale $(X, N)$ relative to some truncation function $h$. Then a version of the $P^{\star}$-differential characteristics of $(X, N)$ is given by $\left(b^{X, N \star}, c^{X, N \star}, F^{X, N^{\star}}, A\right)$, where

$$
\begin{aligned}
b^{X, N \star} & =b^{X, N}+c^{X N}+\int h(x, y) y F^{X, N}(d(x, y)), \\
c^{X, N \star} & =c^{X, N}, \\
\frac{d F^{X, N \star}}{d F^{X, N}}(x, y) & =1+y .
\end{aligned}
$$

Lemma A.10. If $X$ is a $\sigma$-martingale and $\vartheta \in L(X)$, then $\vartheta \cdot X$ is a $\sigma$-martingale as well.

Lemma A.11. Let $X$ be a $\mathbb{R}^{d}$-valued semimartingale and $\vartheta$ an $\mathbb{R}^{d}$-valued predictable process. Then $\vartheta \in L(X)$ if and only if there exists a semimartingale $Z$ with $Z_{0}=0$ and an increasing sequence $\left(D_{n}\right)_{n \in \mathbb{N}}$ of predictable sets such that $D_{n} \uparrow \Omega \times \mathbb{R}_{+}$up to an evanescent set, $\vartheta 1_{D_{n}}$ is bounded, and $1_{D_{n}} \cdot Z=\left(\vartheta 1_{D_{n}}\right) \cdot X$ for any $n \in \mathbb{N}$. In this case $Z=\vartheta \cdot X$.

Acknowledgments. Parts of this research were done while the second author was visiting Helsinki University of Technology. He wants to thank Esko Valkeila for his hospitality. Both authors thank the Isaac Newton Institute for the opportunity to work on the subject during the program on quantitative finance. We are grateful to Thorsten Rheinländer, Christophe Stricker, Shige Peng, Jean Jacod, and Martin Schweizer for valuable comments or discussions. Thanks are also due to an anonymous referee for his detailed suggestions which enhanced the presentation of the results.

## REFERENCES

[1] Albert, A. (1972). Regression and the Moore-Penrose Pseudoinverse. Academic Press, New York. MR0331659
[2] Arai, T. (2004). Minimal martingale measures for jump diffusion processes. J. Appl. Probab. 41 263-270. MR2036287
[3] Arai, T. (2005). An extension of mean-variance hedging to the discontinuous case. Finance Stoch. 9 129-139. MR2210931
[4] Benth, F., Di Nunno, G., Løkka, A., Øksendal, B. and Proske, F. (2003). Explicit representation of the minimal variance portfolio in markets driven by Lévy processes. Math. Finance 13 55-72. MR1968096
[5] Biagini, F. and Guasoni, P. (2002). Mean-variance hedging with random volatility jumps. Stochastic Anal. Appl. 20 471-494. MR1900301
[6] Biagini, F., Guasoni, P. and Pratelli, M. (2000). Mean-variance hedging for stochastic volatility models. Math. Finance 10 109-123. MR1802593
[7] Bobrovnytska, O. and Schweizer, M. (2004). Mean-variance hedging and stochastic control: Beyond the Brownian setting. IEEE Trans. Automat. Control 49 396-408. MR2062252
[8] ČERNÝ, A. (2004). Dynamic programming and mean-variance hedging in discrete time. Appl. Math. Finance 11 1-25.
[9] ČERNÝ, A. (2005). Optimal continuous-time hedging with leptokurtic returns. Math. Finance. To appear.
[10] ČERNÝ, A. and Kallsen, J. (2006). A counterexample concerning the variance-optimal martingale measure. Math. Finance. To appear. Available at http://ssrn.com/abstract=912952.
[11] ČERNÝ, A. and Kallsen, J. (2006). Mean-variance hedging and optimal investment in Heston's model with correlation. SSRN working paper. Available at http://ssrn.com/abstract= 909305.
[12] Choulli, T., Krawczyk, L. and Stricker, C. (1998). $\mathscr{E}$-martingales and their applications in mathematical finance. Ann. Probab. 26 853-876. MR1626523
[13] Delbaen, F., Monat, P., Schachermayer, W., Schweizer, M. and Stricker, C. (1997). Weighted norm inequalities and hedging in incomplete markets. Finance and Stochastics 1 181-227.
[14] Delbaen, F. and Schachermayer, W. (1995). The existence of absolutely continuous local martingale measures. Ann. Appl. Probab. 5 926-945. MR1384360
[15] Delbaen, F. and Schachermayer, W. (1996). Attainable claims with p'th moments. Ann. Inst. H. Poincaré Probab. Statist. 32 743-763. MR1422309
[16] DelbaEn, F. and Schachermayer, W. (1996). The variance-optimal martingale measure for continuous processes. Bernoulli 2 81-105. MR1394053
[17] DI Nunno, G. (2002). Stochastic integral representation, stochastic derivatives and minimal variance hedging. Stochastics Stochastics Rep. 73 181-198. MR1914983
[18] FöLlmer, H. and Schweizer, M. (1991). Hedging of contingent claims under incomplete information. In Applied Stochastic Analysis (M. H. A. Davis and R. J. Elliott, eds.) 389-414. Gordon and Breach, London. MR1108430
[19] Föllmer, H. and Sondermann, D. (1986). Hedging of nonredundant contingent claims. In Contributions to Mathematical Economics 205-223. North-Holland, Amsterdam. MR0902885
[20] Gourieroux, C., Laurent, J. and Pham, H. (1998). Mean-variance hedging and numéraire. Math. Finance 8 179-200. MR1635796
[21] Grandits, P. and Rheinländer, T. (2002). On the minimal entropy martingale measure. Ann. Probab. 30 1003-1038. MR1920099
[22] Hipp, C. (1993). Hedging general claims. In Proceedings of the 3rd AFIR Colloquium, Rome 2 603-613.
[23] Hipp, C. and Taksar, M. (2005). Hedging in incomplete markets and optimal control. To appear.
[24] Hobson, D. (2004). Stochastic volatility models, correlation, and the $q$-optimal martingale measure. Math. Finance 14 537-556. MR2092922
[25] Hou, C. and Karatzas, I. (2004). Least-squares approximation of random variables by stochastic integrals. In Stochastic Analysis and Related Topics in Kyoto 141-166. Math. Soc. Japan, Tokyo. MR2083708
[26] Hubalek, F., Krawczyk, L. and Kallsen, J. (2006). Variance-optimal hedging for processes with stationary independent increments. Ann. Appl. Probab. 16 853-885. MR2244435
[27] JACOD, J. (1979). Calcul Stochastique et Problèmes de Martingales, Springer, Berlin. MR0542115
[28] Jacod, J. and Shiryaev, A. (2003). Limit Theorems for Stochastic Processes, 2nd ed. Springer, Berlin. MR1943877
[29] Kallsen, J. (2004). $\sigma$-localization and $\sigma$-martingales. Theory Probab. Appl. 48 152-163. MR2013413
[30] Laurent, J. and Pham, H. (1999). Dynamic programming and mean-variance hedging. Finance Stoch. 3 83-110. MR1805322
[31] Leitner, J. (2001). Mean-variance efficiency and intertemporal price for risk. Technical Report 00/35, Center of Finance and Econometrics, Univ. Konstanz.
[32] Lim, A. (2004). Quadratic hedging and mean-variance portfolio selection with random parameters in an incomplete market. Math. Oper. Res. 29 132-161. MR2065719
[33] Lim, A. (2005). Mean-variance hedging when there are jumps. SIAM J. Control Optim. 44 1893-1922. MR2193511
[34] Mania, M. and TevZadze, R. (2000). A semimartingale Bellman equation and the varianceoptimal martingale measure. Georgian Math. J. 7 765-792. MR1811929
[35] Mania, M. and Tevzadze, R. (2003). Backward stochastic PDE and imperfect hedging. Int. J. Theor. Appl. Finance 6 663-692. MR2019723
[36] Mania, M. and Tevzadze, R. (2003). A semimartingale backward equation and the variance-optimal martingale measure under general information flow. SIAM J. Control Optim. 42 1703-1726. MR2046382
[37] Monat, P. and Stricker, C. (1995). Föllmer-Schweizer decomposition and mean-variance hedging for general claims. Ann. Probab. 23 605-628. MR1334163
[38] Pham, H. (2000). On quadratic hedging in continuous time. Math. Methods Oper. Res. 51 315-339. MR1761862
[39] Protter, P. (2004). Stochastic Integration and Differential Equations, 2nd ed. Springer, Berlin. MR2020294
[40] Rheinländer, T. and Schweizer, M. (1997). On $L^{2}$-projections on a space of stochastic integrals. Ann. Probab. 25 1810-1831. MR1487437
[41] Schweizer, M. (1991). Option hedging for semimartingales. Stochastic Process. Appl. 37 339-363. MR1102880
[42] SChWEIZER, M. (1994). Approximating random variables by stochastic integrals. Ann. Probab. 22 1536-1575. MR1303653
[43] Schweizer, M. (1995). On the minimal martingale measure and the Föllmer-Schweizer decomposition. Stochastic Anal. Appl. 13 573-599. MR1353193
[44] SCHWEIZER, M. (1996). Approximation pricing and the variance-optimal martingale measure. Ann. Probab. 24 206-236. MR1387633
[45] SCHWEIZER, M. (2001). A guided tour through quadratic hedging approaches. In Option Pricing, Interest Rates and Risk Management (E. Jouini, J. Cvitanic, and M. Musiela, eds.) 538-574. Cambridge Univ. Press. MR1848562
[46] Sekine, J. (2004). On the computation of $L^{2}$-hedging strategy with stochastic volatility. Unpublished manuscript.
[47] Stricker, C. (1990). Arbitrage et lois de martingale. Ann. Inst. H. Poincaré Probab. Statist. 26 451-460. MR1066088
[48] XIA, J. and YAN, J. (2006). Markowitz's portfolio optimization in an incomplete market. Math. Finance 16 203-216. MR2194902

Cass Business School
City University London
106 Bunhill Row
London ECIY 8TZ
United Kingdom
E-MAIL: cerny@martingales.info

HVB-Stiftungsinstitut Für
Finanzmathematik
Zentrum Mathematik
Technische Universität München
Boltzmannstrasse 3
85747 Garching bei MÜnchen
Germany
E-MAIL: kallsen@ma.tum.de


[^0]:    City Research Online:
    http://openaccess.city.ac.uk/
    publications@city.ac.uk

[^1]:    Received June 2005; revised October 2006.
    AMS 2000 subject classifications. 91B28, 60H05, 60G48, 93E20.
    Key words and phrases. Mean-variance hedging, opportunity process, opportunity-neutral measure, incomplete markets.

