# Yang-Mills as a constrained Gaussian 

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#### Abstract

Yang-Mills is reformulated in terms of the logarithmic derivative of the holonomies. The classical equations of motion are recovered, and the path integral is rewritten in two ways, both of which are of the form of a Gaussian satisfying a quadratic constraint.


## Introduction

It is an old and well-known idea that the standard formulation of Yang-Mills may not be the most appropriate one to study its infrared, nonperturbative features. Therefore, there were many attempts to reformulate the theory in terms of other variables which would make these features more transparent. The most studied choice has been the Wilson loop, i.e. the trace of the holonomy of the connection. A large amount of effort was expended and numerous wonderful results were obtained (see e.g. chapter 12 in [1] and the references therein). Nonetheless, the reformulation in terms of Wilson loops was never fully completed. In particular the Yang-Mills path integral was never rewritten in terms of these variables. This is due to the complexity of the constraints which the Wilson loops should satisfy [2].

Another set of variables, which, while known, has been rather understudied, is the logarithmic derivative of the holonomy. It seems to have been first introduced by Mandelstam in [3] and subsequently explored by several authors [4-12]. The definitive analysis of the kinematics of these variables was performed by Gross in [10], unifying and extending many of the previous results. We shall therefore refer below to these variables as Mandelstam-Gross variables.

The aim of this manuscript is to supply a reformulation of Yang-Mills in terms of these variables both at the classical and quantum level. In the next section we shall rewrite the classical action in terms of them and
then will recover the usual equation of motion ${ }^{1}$ In the section after that, we shall rewrite the Yang-Mills path integral in terms of these variables. This essentially amounts to computing the Jacobian of the transformation relating the Mandelstam-Gross variables to the connection.

Before we proceed, we give a summary of our notation choices and a very brief review of these variables. The reader at this stage should probably consult [10] as we will freely use the results from there.

We shall work in the Euclidean signature throughout. All the paths below are based at some arbitrary but fixed point. The parameter range for all the paths is $[0,1]$. Our gauge group will be $S U(N)$ throughout (though it is not difficult to generalize everything to an arbitrary simple Lie group). The Lie algebra of $S U(N)$ will be denoted by $s u(N)$. Now, given an $S U(N)$ connection $A$ on $\mathbb{R}^{D}$, one can define $\mathcal{P}(\gamma)$ along any loop $\gamma$. Following the notation of [10], the Mandelstam-Gross variables will be denoted by $\mathcal{B}_{\mu, s}$, where $\mu=1,2, \ldots, D$ and $s \in(0,1)$, and are defined to be

$$
\mathcal{B}_{\mu, s}(\gamma)=(\mathcal{P}(\gamma))^{-1} \frac{\delta \mathcal{P}(\gamma)}{\delta \gamma(s)}
$$

We would like to emphasize that $(\mu, s)$ should be thought of as a single index, thus making $\mathcal{B}$ an $s u(N)$-valued 1-form on the space of loops.

It can be shown that

$$
\mathcal{B}_{\mu, s}(\gamma)=\left(\mathcal{P}\left(\gamma_{s}\right)\right)^{-1} F_{\mu \nu}(\gamma(s)) \dot{\gamma}^{\nu}(s) \mathcal{P}\left(\gamma_{s}\right) \equiv \mathcal{F}_{\mu}\left(\gamma_{s}\right)
$$

where $F_{\mu \nu}$ stands for the curvature of the connection and $\gamma_{s}$ is the path $\gamma_{s}(t)=\gamma(s t)$.

Given $C_{\mu, s}(\gamma)$, an $s u(N)$-valued 1-form on the space of loops, we say that it is:
i- transverse if $C_{\mu, s}(\gamma) \dot{\gamma}^{\mu}(s)=0$ for every $s \in(0,1)$.
ii- nonanticipating if $\mathcal{C}_{\mu, s}\left(\gamma_{1}\right)=\mathcal{C}_{\mu, s}\left(\gamma_{2}\right)$ for all $\mu$ and all $s<s_{0}$ as long as $\gamma_{1}(s)=\gamma_{2}(s)$ for all $s<s_{0}$.

[^0]The crucial fact that was proven in [10] is that if an $s u(N)$-valued 1-form is transverse, nonanticipating and has curvature zero, i.e. if

$$
\mathcal{G}^{\mu, s ; \nu, t}[\mathcal{C}](\gamma)=\frac{\delta \mathcal{C}^{\nu, t}(\gamma)}{\delta \gamma^{\mu}(s)}-\frac{\delta \mathcal{C}^{\mu, s}(\gamma)}{\delta \gamma^{\nu}(t)}+\left[\mathcal{C}^{\mu, s}(\gamma), \mathcal{C}^{\nu, t}(\gamma)\right]=0
$$

for all $\mu, \nu$ and for every $s, t \in(0,1)$, then there is a connection $A$ whose corresponding $\mathcal{B}$ is the given 1 -form $\mathcal{C}_{\mu, s}(\gamma)$. Moreover, one can actually recover $A$ from $\mathcal{B}$ by a linear map $\mathcal{T}$ given by

$$
\begin{equation*}
A_{\mu}(x)=\mathcal{T}(\mathcal{B}) \equiv \int_{0}^{\frac{1}{2}} d s \mathcal{B}_{\mu, s}\left(\sigma_{x}\right) \partial^{\mu} \sigma_{x}(s) \tag{1}
\end{equation*}
$$

where $\sigma_{x}$ is the loop which is given by

$$
\sigma_{x}(s)= \begin{cases}2 s x & , \quad 0 \leq s \leq \frac{1}{2} \\ 2(1-s) x & , \quad \frac{1}{2} \leq s \leq 1\end{cases}
$$

It can be shown that the connection $\mathcal{T}(\mathcal{B})$ satisfies the radial gauge, and that the map $\mathcal{T}$ establishes a bijection between the set of all transverse, nonanticipating, curvature zero 1 -forms on the space of loops and the set of all connections satisfying the radial gauge with the composition $A \rightarrow$ $\mathcal{B}(A) \rightarrow \mathcal{T}(\mathcal{B}(A))$ being the identity on the space of such connections. It is assumed below that wherever gauge-fixing is needed, the radial gauge is used.

## The Classical Action and the Equations of Motion

We begin by giving the Yang-Mills action in terms of the Mandelstam-Gross variables [8]. It is given by the following formula

$$
\begin{equation*}
\frac{D}{\mathcal{V}} \int_{(0,1)} d s \int \mathcal{D} \gamma \frac{\operatorname{Tr}\left(\mathcal{B}_{\mu, s}(\gamma) \mathcal{B}^{\mu, s}(\gamma)\right)}{\dot{\gamma}_{\nu}(s) \dot{\gamma}^{\nu}(s)} \tag{2}
\end{equation*}
$$

where $D$ is the dimensions of the spacetime and $\mathcal{D} \gamma$ and $\mathcal{V}$ stand for the following formal expressions

$$
\mathcal{D} \gamma=\prod_{t \in(0,1)} d \gamma(t), \quad \text { and } \quad \mathcal{V}=\int \prod_{t \in(0,1) ; t \neq s} d \gamma(t)
$$

In other words, $\mathcal{D} \gamma$ stands for the (nonexistent) 'flat' measure on the space of loops $\Gamma$, and $\mathcal{V}$ is the (infinite) volume assigned by the measure to a
'slice' of the space, obtained from $\mathcal{D} \gamma$ by peeling off the integration over $\gamma(s)$.
Of course, at face value, these expressions are ill-defined. Nonetheless, we are going to continue using them below for two reasons. First, it simplifies the formal calculations considerably, mainly because it allows integration by parts. Second, it is in fact rather easy, if desired, to make sense of the formal expressions above by interpreting $\mathcal{D} \gamma$ to stand for an appropriate, well-defined measure. What measure would classify as appropriate? As we shall see below, it should satisfy the following properties:
i- $\mathcal{D} \gamma_{x}$ is supported on loops which are at least $C^{1}$.
ii- $\int \mathcal{D} \gamma_{x}<\infty$.
iii- For any $s \in(0,1)$, the integral $\int \mathcal{D} \gamma$ should factorize into an integral over the value of $\gamma$ and $\dot{\gamma}$ at $s$, an integral over the part of $\gamma$ from 0 to $s$ and a final integral over the part from $s$ to 1.
iv- The measure corresponding to the integral over $\dot{\gamma}(s)$ in the previous item is even.

Contrary to what it might seem, it is rather straightforward to construct measures which satisfy the conditions above. In fact the Brownian bridge measure almost fits the bill. The only issue is that the Brownian paths are not sufficiently regular. An obvious substitute would be a Gaussian measure whose covariance is sufficiently smoothing. The simplest choice for a covariance that works is to take the bilinear form $\frac{1}{\left(-\nabla^{2}+1\right)^{2}}$. An additional attractive feature is that one can have an entire family of measures $\mathcal{D} \gamma_{\epsilon}$ parametrized by $\epsilon>0$, whose covariances are given by $\frac{1}{\epsilon\left(-\nabla^{2}+1\right)^{2}}$, such that as $\epsilon \rightarrow 0$, they can be thought to approach the 'flat' one. This would justify the integration by parts manipulations in what follows.

Now, let us demonstrate that (2) is indeed the Yang-Mills action in disguise. Note that

$$
\begin{aligned}
& \frac{D}{\mathcal{V}} \int_{(0,1)} d s \int \mathcal{D} \gamma \frac{\operatorname{Tr}\left(\mathcal{B}_{\mu, s}(\gamma) \mathcal{B}^{\mu, s}(\gamma)\right)}{\dot{\gamma}_{\nu}(s) \dot{\gamma}^{\nu}(s)}=\frac{D}{\mathcal{V}} \int_{(0,1)} d s \int \mathcal{D} \gamma \frac{\operatorname{Tr}\left(\mathcal{F}_{\mu}\left(\gamma_{s}\right) \mathcal{F}^{\mu}\left(\gamma_{s}\right)\right)}{\dot{\gamma}_{\nu}(s) \dot{\gamma}^{\nu}(s)} \\
& =\frac{D}{\mathcal{V}} \int_{(0,1)} d s \int d \gamma(s) \int_{t \in(0,1) ; t \neq s} d \gamma(t) \frac{\operatorname{Tr}\left(F_{\mu \rho}(\gamma(s)) F^{\mu \rho^{\prime}}(\gamma(s))\right) \dot{\gamma}^{\rho}(s) \dot{\gamma}_{\rho^{\prime}}(s)}{\dot{\gamma}_{\nu}(s) \dot{\gamma}^{\nu}(s)} \\
& =\int_{(0,1)} d s \int d x \operatorname{Tr}\left(F_{\mu \rho}(x) F^{\mu \rho^{\prime}}(x)\right) \delta_{\rho^{\prime}}^{\rho}=\int d x \operatorname{Tr}\left(F_{\mu \rho}(x) F^{\mu \rho}(x)\right) .
\end{aligned}
$$

The above calculation carries through if we use a well-defined measure satisfying (i)-(iv). The first two items will make the expression for the action well-defined. The factorization described in (iii) would be the analogue of the factorization in the third line above, except that we will have an additional integration over $\dot{\gamma}(s)$ (here it is subsumed in the $\prod_{t \in(0,1) ; t \neq s} d \gamma(t)$ ). Finally, the last item would reproduce the Kronecker delta in the fourth line.

We want to show now that we recover the usual Yang-Mills equations of motion by setting to zero the variation of the action (2) over the set of transverse, nonanticipating $\mathcal{B}_{\mu, s}$ 's satisfying the constraint

$$
\begin{equation*}
\mathcal{G}^{\mu, s ; \nu, t}(\gamma)=\frac{\delta \mathcal{B}^{\nu, t}(\gamma)}{\delta \gamma^{\mu}(s)}-\frac{\delta \mathcal{B}^{\mu, s}(\gamma)}{\delta \gamma^{\nu}(t)}+\left[\mathcal{B}^{\mu, s}(\gamma), \mathcal{B}^{\nu, t}(\gamma)\right]=0 \tag{3}
\end{equation*}
$$

At this point, one can use the result (see Theorem 3.13 and Corollary 3.5 in [10]) that (under suitable smoothness assumptions) $\mathcal{B}^{\mu, s}$ 's provide a faithful coordinatization of the restricted gauge orbit space. Since the action is a smooth function(al) on this space, and since the critical points of a function(a) are unaffected by changes of variables applied to its domain, we have what we want. Despite the fact that there is nothing wrong with the argument above, it would certainly be more satisfying if one can supply a more explicit computational proof.

Therefore, in order to perform the variation of (2) subject to the constraint (3), we introduce the Lie algebra valued Lagrange multiplier $\Xi_{\mu, s ; \nu, t}(\gamma)$. Since $\mathcal{G}^{\mu, s ; \nu, t}$ is antisymmetric under $\mu, s \longleftrightarrow \nu, t$, we can, without loss of generality, assume that $\Xi_{\mu, s ; \nu, t}$ is antisymmetric as well. We thus add the following term

$$
\begin{equation*}
\int \mathcal{D} \gamma \int_{(0,1)^{2}} d s d t \operatorname{Tr}\left(\Xi_{\mu, s ; \nu, t}(\gamma) \mathcal{G}^{\mu, s ; \nu, t}(\gamma)\right) \tag{4}
\end{equation*}
$$

to our action and vary $\mathcal{B}^{\mu, s}$ and $\Xi_{\mu, s ; \nu, t}$. Of course, variation of the $\Xi$ simply gives back the constraint. We therefore only need to vary the $\mathcal{B}$ 's. Performing the variation is straightforward as the full action is quadratic in the $\mathcal{B}$ 's with the result being the following equation

$$
\begin{equation*}
\frac{\mathcal{B}^{\mu, s}(\gamma)}{\mathcal{V} \dot{\gamma}_{\rho}(s) \dot{\gamma}^{\rho}(s)}=\int_{(0,1)} d t\left(\frac{\delta \Xi_{\mu, s ; \nu, t}(\gamma)}{\delta \gamma^{\nu}(t)}+\left[\mathcal{B}^{\nu, t}(\gamma), \Xi_{\mu, s ; \nu, t}(\gamma)\right]\right) \tag{5}
\end{equation*}
$$

Before we proceed, let us identify our target, i.e. let us give the YangMills equations of motion in terms of the Mandelstam-Gross variables. This
is well-known [13] and is given by the expression

$$
\frac{\delta \mathcal{F}^{\mu}\left(\gamma^{\prime}\right)}{\delta \gamma^{\prime \mu}(1)}=0
$$

where this equation should hold for any path $\gamma^{\prime}$.
Now, note that

$$
\begin{aligned}
\int_{(0,1)} & d s\left(\frac{\delta}{\delta \gamma^{\mu}(s)}\left(\frac{\mathcal{B}^{\mu, s}(\gamma)}{\mathcal{V} \dot{\gamma}_{\rho}(s) \dot{\gamma}^{\rho}(s)}\right)\right) \\
& =\int_{(0,1)} d s\left(\frac{\delta}{\delta \gamma^{\mu}(s)}\left(\frac{\mathcal{F}^{\mu}\left(\gamma_{s}\right)}{\mathcal{V} \dot{\gamma}_{\rho}(s) \dot{\gamma}^{\rho}(s)}\right)\right) \\
& =\int_{(0,1)} d s\left(\frac{\frac{\delta \mathcal{F}^{\mu}\left(\gamma_{s}\right)}{\delta \gamma^{\mu}(s)}}{\mathcal{V} \dot{\gamma}_{\rho}(s) \dot{\gamma}^{\rho}(s)}-\mathcal{F}^{\mu}\left(\gamma_{s}\right) \frac{2 \dot{\delta}(0) \delta_{\mu}^{\rho} \dot{\gamma}_{\rho}(s)}{\mathcal{V}\left(\dot{\gamma}_{\rho}(s) \dot{\gamma}^{\rho}(s)\right)^{2}}\right) \\
& =\int_{(0,1)} d s\left(\frac{1}{\mathcal{V} \dot{\gamma}_{\rho}(s) \dot{\gamma}^{\rho}(s)} \frac{\delta \mathcal{F}^{\mu}\left(\gamma_{s}\right)}{\delta \gamma_{s}^{\mu}(1)}\right)
\end{aligned}
$$

where the second term on the third line vanishes (no matter how we define $\dot{\delta}(0))$ in view of the transversality constraint that the $\mathcal{B}$ 's (and hence the $\mathcal{F}^{\prime}$ 's) satisfy. Clearly, we have that if $\frac{\delta \mathcal{F}^{\mu}\left(\gamma^{\prime}\right)}{\delta \gamma^{\prime \mu}(1)}=0$ for every path $\gamma^{\prime}$, then

$$
\int_{(0,1)} d s\left(\frac{\delta}{\delta \gamma^{\mu}(s)}\left(\frac{\mathcal{B}^{\mu, s}(\gamma)}{\mathcal{V} \dot{\gamma}_{\rho}(s) \dot{\gamma}^{\rho}(s)}\right)\right)=0
$$

for every loop $\gamma$. We claim that the converse implication holds as well. To see this, let $\gamma$ be a loop and let $s_{0} \in(0,1)$. We know that

$$
\begin{aligned}
0 & =\int_{(0,1)} d s\left(\frac{1}{\mathcal{V} \dot{\gamma}_{\rho}(s) \dot{\gamma}^{\rho}(s)} \frac{\delta \mathcal{F}^{\mu}\left(\gamma_{s}\right)}{\delta \gamma_{s}^{\mu}(1)}\right) \\
& =\int_{0}^{s_{0}} d s\left(\frac{1}{\mathcal{V} \dot{\gamma}_{\rho}(s) \dot{\gamma}^{\rho}(s)} \frac{\delta \mathcal{F}^{\mu}\left(\gamma_{s}\right)}{\delta \gamma_{s}^{\mu}(1)}\right)+\int_{s_{0}}^{1} d s\left(\frac{1}{\mathcal{V} \dot{\gamma}_{\rho}(s) \dot{\gamma}^{\rho}(s)} \frac{\delta \mathcal{F}^{\mu}\left(\gamma_{s}\right)}{\delta \gamma_{s}^{\mu}(1)}\right)
\end{aligned}
$$

If we replace $\gamma$ with another loop $\tilde{\gamma}$ which coincides with $\gamma$ for $t \leq s_{0}$, then the first term above does not change. The second term can be made as small as one pleases by having the $\dot{\tilde{\gamma}}(t)$ very large while $\tilde{\gamma}(t)$ stays bounded
for $t \geq s_{0} \|^{2}$ Therefore, we can conclude that for any $s_{0} \in(0,1)$

$$
\int_{0}^{s_{0}} d s\left(\frac{1}{\mathcal{V} \dot{\gamma}_{\rho}(s) \dot{\gamma}^{\rho}(s)} \frac{\delta \mathcal{F}^{\mu}\left(\gamma_{s}\right)}{\delta \gamma_{s}^{\mu}(1)}\right)=0
$$

From this it follows that the integrand itself vanishes and we have our claim.

The upshot of the discussion above is that, to verify that the Yang-Mills equations of motion hold, it is sufficient to apply $\int_{(0,1)} d s \frac{\delta}{\delta \gamma^{\mu}(s)}$ to the right hand side of (5) and then verify that the result vanishes.

Therefore consider

$$
\begin{gather*}
\int_{(0,1)^{2}} d s d t\left(\frac{\delta^{2} \Xi_{\mu, s ; \nu, t}(\gamma)}{\delta \gamma^{\mu}(s) \delta \gamma^{\nu}(t)}+\left[\frac{\delta \mathcal{B}^{\nu, t}(\gamma)}{\delta \gamma^{\mu}(s)}, \Xi_{\mu, s ; \nu, t}(\gamma)\right]\right.  \tag{6}\\
\left.+\left[\mathcal{B}^{\nu, t}(\gamma), \frac{\delta \Xi_{\mu, s ; \nu, t}(\gamma)}{\delta \gamma^{\mu}(s)}\right]\right)
\end{gather*}
$$

The first term vanishes due to the antisymmetry of $\Xi$. The second term can be rewritten in the following way

$$
\begin{array}{rl}
\int_{(0,1)^{2}} & d s d t\left(\left[\frac{\delta \mathcal{B}^{\nu, t}(\gamma)}{\delta \gamma^{\mu}(s)}, \Xi_{\mu, s ; \nu, t}(\gamma)\right]\right)  \tag{7}\\
& =\int_{(0,1)^{2}} d s d t\left(\frac{1}{2}\left[\frac{\delta \mathcal{B}^{\nu, t}(\gamma)}{\delta \gamma^{\mu}(s)}-\frac{\delta \mathcal{B}^{\mu, s}(\gamma)}{\delta \gamma^{\nu}(t)}, \Xi_{\mu, s ; \nu, t}(\gamma)\right]\right) \\
& =-\int_{(0,1)^{2}} d s d t\left(\frac{1}{2}\left[\left[\mathcal{B}^{\mu, s}(\gamma), \mathcal{B}^{\nu, t}(\gamma)\right], \Xi_{\mu, s ; \nu, t}(\gamma)\right]\right),
\end{array}
$$

where the antisymmetry of $\Xi$ justifies the first equality, and (3) justifies the second.

[^1]Looking at the third term in (6), and using (5) we get
(8) $\quad \int_{(0,1)^{2}} d s d t\left(\left[\mathcal{B}^{\nu, t}(\gamma), \frac{\delta \Xi_{\mu, s ; \nu, t}(\gamma)}{\delta \gamma^{\mu}(s)}\right]\right)$

$$
\begin{aligned}
& =\int_{(0,1)^{2}} d s d t\left(\left[\mathcal{B}^{\nu, t}(\gamma),-\frac{\mathcal{B}^{\nu, t}(\gamma)}{\mathcal{V} \dot{\gamma}_{\rho}(t) \dot{\gamma}^{\rho}(t)}+\left[\mathcal{B}^{\mu, s}(\gamma), \Xi_{\nu, t ; \mu, s}(\gamma)\right]\right)\right. \\
& =\int_{(0,1)^{2}} d s d t\left(\left[\mathcal{B}^{\nu, t}(\gamma),\left[\mathcal{B}^{\mu, s}(\gamma), \Xi_{\nu, t ; \mu, s}(\gamma)\right]\right)\right.
\end{aligned}
$$

Now, adding together (7) and (8) and dropping the $\gamma$ 's to reduce clutter we get

$$
\begin{aligned}
& \int_{(0,1)^{2}} d s d t\left\{-\frac{1}{2}\left[\left[\mathcal{B}^{\mu, s}, \mathcal{B}^{\nu, t}\right], \Xi_{\mu, s ; \nu, t}\right]+\left[\mathcal{B}^{\nu, t},\left[\mathcal{B}^{\mu, s}, \Xi_{\nu, t ; \mu, s}\right]\right\}\right. \\
& =\int_{(0,1)^{2}} d s d t\left\{-\frac{1}{2}\left(\mathcal{B}^{\mu, s} \mathcal{B}^{\nu, t} \Xi_{\mu, s ; \nu, t}-\mathcal{B}^{\nu, t} \mathcal{B}^{\mu, s} \Xi_{\mu, s ; \nu, t}-\Xi_{\mu, s ; \nu, t} \mathcal{B}^{\mu, s} \mathcal{B}^{\nu, t}\right.\right. \\
& \left.\left.\quad+\Xi_{\mu, s ; \nu, t} \mathcal{B}^{\nu, t} \mathcal{B}^{\mu, s}\right)\right)+\left(\mathcal{B}^{\nu, t} \mathcal{B}^{\mu, s} \Xi_{\nu, t ; \mu, s}-\mathcal{B}^{\nu, t} \Xi_{\nu, t ; \mu, s} \mathcal{B}^{\mu, s}-\mathcal{B}^{\mu, s} \Xi_{\nu, t ; \mu, s} \mathcal{B}^{\nu, t}\right. \\
& \left.\left.\quad+\Xi_{\nu, t ; \mu, s} \mathcal{B}^{\mu, s} \mathcal{B}^{\nu, t}\right)\right\} \\
& =\int_{(0,1)^{2}} d s d t\left(\left(\mathcal{B}^{\nu, t} \mathcal{B}^{\mu, s} \Xi_{\mu, s ; \nu, t}+\Xi_{\mu, s ; \nu, t} \mathcal{B}^{\mu, s} \mathcal{B}^{\nu, t}\right)\right. \\
& \left.\quad+\left(\mathcal{B}^{\nu, t} \mathcal{B}^{\mu, s} \Xi_{\nu, t ; \mu, s}+\Xi_{\nu, t ; \mu, s} \mathcal{B}^{\mu, s} \mathcal{B}^{\nu, t}\right)\right) \\
& =0
\end{aligned}
$$

We have thus recovered the Yang-Mills equations of motion.

## The Path Integral

Our goal now is to extend the reformulation of Yang-Mills in terms of the $\mathcal{B}$ 's to the quantum setting. We shall do this by rewriting the Yang-Mills path integral in terms of the new variables. As we've mentioned above, this amounts to computing the Jacobian of the change of variables $A \rightarrow$ $\mathcal{B}(A)$. This shall be done in three steps. First, we'll derive the relevant finite dimensional formulae. This will be the longest part of the calculation. Then we'll calculate the Jacobian for the closely related case of the principal chiral model. We will finish by dealing with the Yang-Mills case.

## Finite Dimensions

Suppose that we have the following three maps

$$
\begin{aligned}
F: \mathbb{R}^{n+m} \rightarrow \mathbb{R} & \text { "the integrand" } \\
g: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n+m} & \text { "the parametrization" } \\
h: \mathbb{R}^{n+m} \rightarrow \mathbb{R}^{l} \quad & \text { "the constraint" }
\end{aligned}
$$

We are going to suppose that all the maps above are as smooth as is required. Moreover, we will assume that $g$ is an embedding, i.e. that it is an immersion and that $g\left(\mathbb{R}^{n}\right)$ is an embedded $n$-dimensional submanifold of $\mathbb{R}^{n+m}$. We shall denote this submanifold by $S$. Finally, we will assume that $g\left(\mathbb{R}^{n}\right)=S=h^{-1}(0)$.

We are interested in rewriting

$$
\int_{\mathbb{R}^{n}} F(g(x)) d x
$$

in terms of an integral over $\mathbb{R}^{n+m}$ with an insertion of a delta function containing $h$. As a first step, we rewrite the integral above in terms of the Hausdorff measure on the submanifold $S$. Then, using the area formula ${ }^{3}$ we have the following equality

$$
\begin{equation*}
\int_{\mathbb{R}^{n}} F(g(x)) d x=\int_{S} F(y) \frac{1}{J g(y)} d \mathcal{H}(y) \tag{9}
\end{equation*}
$$

where $d \mathcal{H}$ stands for the Hausdorff measure on $S$. There are two equivalent definitions of the Jacobian factor $J g(y) \|^{4}$ The more familiar one is

$$
\begin{equation*}
J g(y)=\sqrt{\operatorname{det}\left(\left.\left.\left(D g^{\dagger}\right)\right|_{g^{-1}(y)} D g\right|_{g^{-1}(y)}\right)} \tag{10}
\end{equation*}
$$

where $D g$ is the derivative of $g$. The second definition [15] of $J g(y)$ is
(11) $J g(y)=$ volume of the parallelepiped spanned by $D g\left(v_{1}\right), \ldots, D g\left(v_{n}\right)$,

[^2]where $D g$ is evaluated at $g^{-1}(y)$ and $v_{1}, \ldots, v_{n}$ are an orthonormal basis in $\mathbb{R}^{n}$. The word 'volume' is understood to be the $n$-dimensional Hausdorff measure.

Formula (9) can be easily generalized to the case when the domain of the "parametrization" $g$ instead of being $\mathbb{R}^{n}$ is an $n$-dimensional Riemannian manifold $M$. In this case, (9) is modified to become

$$
\begin{equation*}
\int_{M} F(g(x)) d \operatorname{vol}(x)=\int_{S} F(y) \frac{1}{J g(y)} d \mathcal{H}(y) \tag{12}
\end{equation*}
$$

where $d \mathrm{vol}$ stands for the volume form induced from the Riemannian metric on $M . J g$ is again given by (11) where $v_{1}, \ldots, v_{n}$ are now an orthonormal basis of the tangent plane to a point of $M \square^{5}$

Now, let $K: \mathbb{R}^{n+m} \rightarrow \mathbb{R}$ be a function and assume for now that $h$ is a submersion (and thus $l=m$ ). Using the coarea formuld ${ }^{6}$, we have that

$$
\begin{equation*}
\int_{\mathbb{R}^{n+m}} K(z) \frac{e^{-\frac{h^{2}(z)}{\epsilon}}}{(\sqrt{\pi \epsilon})^{m}} J h(z) d z=\int_{\mathbb{R}^{m}} \frac{e^{-\frac{c^{2}}{\epsilon}}}{(\sqrt{\pi \epsilon})^{m}}\left(\int_{h^{-1}(c)} K(t) d \mathcal{H}(t)\right) d c \tag{13}
\end{equation*}
$$

Again, analogously to $J g(y)$, there are two definitions for $J h(z)$. The first one is the expression

$$
\begin{equation*}
J h(z)=\sqrt{\operatorname{det}\left(\left.\left.D h\right|_{z}(D h)^{\dagger}\right|_{z}\right)} \tag{14}
\end{equation*}
$$

The second one, which will play a much bigger role below, is
(15) $J h(z)=$ volume of the parallelepiped spanned by $D h\left(v_{1}\right), \ldots, D h\left(v_{l}\right)$,
where $v_{1}, \ldots, v_{l}$ is an othonormal basis of the orthogonal complement of the kernel of of $D h$.

[^3]If we now send $\epsilon \rightarrow 0$, let $K(z)=F(z) \frac{1}{J g(z)}$, and recall that $S=h^{-1}(0)$, we get that

$$
\int_{S} F(y) \frac{1}{J g(y)} d \mathcal{H}(y)=\int_{\mathbb{R}^{n+m}} F(z) \delta(h(z)) \frac{J h(z)}{J g(z)} d z
$$

We would like to generalize this expression to the case when $h$ is not a submersion, but is of constant rank. Therefore, suppose that the rank of $h$ is $m \leq l$. Assume first that, in fact, $D h$ (and thus $h$ ) maps $\mathbb{R}^{n+m}$ into the $m$ dimensional subspace of $\mathbb{R}^{l}$ consisting of all those $l$-tuples whose last $l-m$ entries are zero. Formula (13) still applies, but there are two important differences. First, only definition (15) of $J h$ works as long as 'volume' is understood to mean the $m$-dimensional Hausdorff measure in $\mathbb{R}^{l}$. The second is that the limit $\epsilon \rightarrow 0$ does not produce $\delta(h)$ on the left hand side. If we denote the components of $h$ by $\left(h_{1}, h_{2}, \ldots, h_{l}\right)$, then very ${ }^{7}$ formally,

$$
\begin{equation*}
\delta(h)=\delta\left(h_{1}\right) \delta\left(h_{2}\right) \ldots \delta\left(h_{l}\right)=\lim _{\epsilon \rightarrow 0} \frac{1}{(\sqrt{\pi \epsilon})^{l}} e^{-\frac{h_{1}^{2}}{\epsilon}} \ldots e^{-\frac{h_{1}^{2}}{\epsilon}}=\lim _{\epsilon \rightarrow 0} \frac{e^{-\frac{h^{2}(z)}{\epsilon}}}{(\sqrt{\pi \epsilon})^{l}} \tag{16}
\end{equation*}
$$

which differs by a factor $\sqrt{\pi \epsilon}^{l-m}$ from the expression appearing in 13 .
Note however that our end goal is to compute certain path integrals, and in such computations it is not the value of the integral that matters, but its logarithmic derivatives with respect to the sources. Therefore, introducing the notation $\simeq$ which stands for "proportional" $\|^{8}$ and keeping in mind that $\delta(h)$ should be understood as a limit (or perhaps through its Fourier transform) we can write that

$$
\begin{equation*}
\int_{\mathbb{R}^{n+m}} K(z) \delta(h(z)) J h(z) d z \simeq \int_{S} K(y) d \mathcal{H}(y) \tag{17}
\end{equation*}
$$

[^4]Above, we've assumed that the range of $D h(z)$ is a certain fixed $\mathbb{R}^{m}$ for all $z$. There is no difficulty however in generalizing the arguments above to the general case. Simply replace $h(z)$ with $\tilde{h}(z)=O(z) h(z)$ where $O(z)$ it he rotation which takes $D h(z)$ to the $m$-dimensional subspace consisting of all those $l$-tuples whose last $l-m$ entries vanish. All the arguments above then apply, and equation (17) applies with $h$ replaced by $\tilde{h}$. However, as is evident from (15) and 16$), J h(z)=J \tilde{h}(z)$ and $\delta(h)=\delta(\tilde{h})$. Therefore, we can see that 17) holds for any $h$ of constant rank.

Now, combining 12 and 17 , and setting $K(z)=F(z) \frac{1}{J g(z)}$, we get that

$$
\begin{equation*}
\int_{M} F(g(x)) d \operatorname{vol}(x) \simeq \int_{\mathbb{R}^{n+m}} F(z) \delta(h(z)) \frac{J h(z)}{J g(z)} d z \tag{18}
\end{equation*}
$$

In addition to the general formula above, we will need a certain special case. This is the situation when the domain of the parametrization $g$ happens to be a subspace of $\mathbb{R}^{n+m}$ and $g$ is itself is a graph. More precisely, for $z \in \mathbb{R}^{n+m}$, write it as $z=\left(z_{\|}, z_{\perp}\right)$, where $z_{\|} \in \mathbb{R}^{n}$ and $z_{\perp} \in \mathbb{R}^{m}$. Assume that $g$ is a function of $z_{\|}$and has the form

$$
g\left(z_{\|}\right)=\left(z_{\|}, g_{\perp}\left(z_{\|}\right)\right)
$$

The derivative of $g$ is easily computed to be

$$
D g=\left(\frac{I}{D g_{\perp}}\right)
$$

where the right hand side stands for an $n+m$ by $n$ matrix with the top $n$ by $n$ chunk being the identity while the bottom $m$ by $n$ piece is the derivative of $g_{\perp}$. Hence, it follows from (10) that

$$
\begin{align*}
J g & =\sqrt{\operatorname{det}\left((D g)^{T}(D g)\right)}  \tag{19}\\
& =\sqrt{\operatorname{det}\left(I+\left(D g_{\perp}\right)^{T}\left(D g_{\perp}\right)\right)} \\
& =\sqrt{\operatorname{det}\left(I+\left(D g_{\perp}\right)\left(D g_{\perp}\right)^{T}\right)}
\end{align*}
$$

where Sylvester's determinant theorem was used to get to the last line.

Assume now that $h$ is a submersion, and write

$$
D h=\left(\begin{array}{l|l}
D_{\|} h & \mid D_{\perp} h
\end{array}\right),
$$

where $D_{\|} h$ and $D_{\perp} h$ stands for the submatrices of $D h$ containing the partial derivatives with respect to the components of $z_{\|}$and $z_{\perp}$ respectively. It follows from the implicit function theorem that

$$
D g_{\perp}=-\left(D_{\perp} h\right)^{-1} D_{\|} h
$$

Consequently, we get that

$$
\begin{align*}
\operatorname{det} & \left((D h)(D h)^{T}\right) \operatorname{det}\left(D_{\perp} h\right)^{-2}  \tag{20}\\
& =\operatorname{det}\left(\left(D_{\perp} h\right)^{-1}(D h)(D h)^{T}\left(\left(D_{\perp} h\right)^{-1}\right)^{T}\right) \\
& =\operatorname{det}\left(\left(-D g_{\perp} \mid \quad I\right)\left(\frac{-\left(D g_{\perp}\right)^{T}}{I}\right)\right) \\
& =\operatorname{det}\left(I+\left(D g_{\perp}\right)\left(D g_{\perp}\right)^{T}\right)
\end{align*}
$$

Recalling (14), and putting together (19) and (20), we get that in the case we are describing

$$
\frac{J h}{J g}= \pm \operatorname{det}\left(D_{\perp} h\right)
$$

Plugging this into (18), we get that

$$
\begin{equation*}
\int_{R^{n}} F\left(g\left(z_{\|}\right)\right) d z_{\|} \simeq \int_{\mathbb{R}^{n+m}} F(z) \delta(h(z)) J h_{\perp} d z \tag{21}
\end{equation*}
$$

where $J h_{\perp}$ has the following two equivalent definitions

$$
\begin{align*}
J h_{\perp} & =\operatorname{det}\left(D_{\perp} h\right)  \tag{22}\\
& =\text { volume of the parallelepiped spanned by } D h\left(v_{1}\right), \ldots, D h\left(v_{m}\right)
\end{align*}
$$

where $v_{1}, \ldots, v_{m}$ is an orthonormal basis of the space of $z_{\perp}$ 's.
So far, we know that (21) holds only if $h$ is a submersion. However, if $h$ is merely of constant rank, then proceeding as before by replacing $h$ with $\tilde{h}=O h$, where $O$ is an appropriately chosen rotation, and repeating the same arguments as before, we see that the relation continues to hold in
that case as well. Of course, in this case, only the second definition of $J h_{\perp}$ above makes sense where 'volume' now means, as before, the $m$-dimensional Hausdorff measure.

## Principal Chiral Model

Let us move on now and consider the two dimensional $S U(N)$ principal chiral model. The fundamental fields in this case are group-valued maps $\phi: \mathbb{R}^{2} \rightarrow S U(N)$ and the action is $\int_{\mathbb{R}^{2}} \operatorname{Tr}\left(\left(\phi^{-1} \partial_{\mu} \phi\right)\left(\phi^{-1} \partial^{\mu} \phi\right)\right)$. It is natural to change variables to $A_{\mu}=\phi^{-1} \partial_{\mu} \phi$. Of course, this is exactly the form that a flat connection should have. Therefore, we shall attempt to rewrite the path integral of the model as an integral of a quadratic action, simply the Proca mass term, over the space of flat connections.

We would like to use (18) above assuming that it generalizes formally to the infinite dimensional setting. Our left hand side now is

$$
\begin{equation*}
\int \mathcal{D} \phi e^{-\int_{\mathbb{R}^{2}} \operatorname{Tr}\left(\left(\phi^{-1} \partial_{\mu} \phi\right)\left(\phi^{-1} \partial^{\mu} \phi\right)\right)} \tag{23}
\end{equation*}
$$

Therefore, we see that the analogue of the manifold $M$ is the space of fundamental fields. We claim that this is a Riemannian manifold and the formal measure $\mathcal{D} \phi$ is indeed the one induced from the Riemannian structure. To see this, note that a tangent vector can be represented by the derivative at $t=0$ of a curve of the form

$$
\begin{equation*}
t \rightarrow(x \rightarrow \psi(x, t) \phi(x)) \tag{24}
\end{equation*}
$$

where $\psi(x, 0)=\mathbb{I}$. The inner product of two such vectors, corresponding to $\psi_{1}$ and $\psi_{2}$ is given by $\int_{\mathbb{R}^{2}} \operatorname{Tr}\left(\dot{\psi}_{1}(x, 0) \dot{\psi}_{2}(x, 0)\right) d x$. In other words, this is just the usual Riemannian metric on $S U(N)$ 'summed' over all points of spacetime. Since the measure $\mathcal{D} \phi$ is similarly a product of Haar measures over all points of spacetime, and since the Haar measure is induced by the Riemannian structure on the group, we have our claim.

Continuing the comparison of (23) with the left hand side of (18) we see that the map $x \rightarrow g(x)$ corresponds to $\phi \rightarrow A_{\mu}=\phi^{-1} \partial_{\mu} \phi$ and that the analogue of $F(\cdot)$ is $e^{-\int_{\mathbb{R}^{2}} \operatorname{Tr}\left(A_{\mu} A^{\mu}\right)}$.

Let us now discuss the analogue of the right hand side of (18) which corresponds to 23 . Clearly, the analogue of $\mathbb{R}^{n+m}$ is the space of all $s u(N)-$ valued 1-forms. The analogue of the standard Euclidean metric on this space is

$$
\begin{equation*}
A_{\mu} \cdot \tilde{A}^{\mu}=\int_{\mathbb{R}^{2}} \operatorname{Tr}\left(A_{\mu} \tilde{A}^{\mu}\right) \tag{25}
\end{equation*}
$$

The analogue of the space $\mathbb{R}^{l}$ is, as we shall see below, the space of $s u(N)$-valued 2-forms. The analogue of the standard Euclidean metric on this space is

$$
\begin{equation*}
C_{\mu \nu} \cdot \tilde{C}^{\mu \nu}=\int_{\mathbb{R}^{2}} \operatorname{Tr}\left(C_{\mu \nu} \tilde{C}^{\mu \nu}\right) \tag{26}
\end{equation*}
$$

Finally, the analogue of $d z$ would be $\mathcal{D} A$ which is formally equal to a product over Lebesgue measures at each point of spacetime. It remains to identify the analogues of $h, J g$ and $J h$.

We shall begin with $J g$. Let us thus compute the action of the differential of the map $\phi \rightarrow A_{\mu}$ on a tangent vector. Therefore, take a curve of the form (24). Note that

$$
\begin{aligned}
\left.\frac{d}{d t}\left(\phi^{-1} \psi^{-1} \partial_{\mu}(\psi \phi)\right)\right|_{t=0} & =\left.\phi^{-1}\left(\dot{\psi^{-1}}\right) \partial_{\mu}(\psi \phi)\right|_{t=0}+\left.\phi^{-1} \psi^{-1} \partial_{\mu}(\dot{\psi} \phi)\right|_{t=0} \\
& =\phi^{-1} \partial_{\mu}\left(\left.\dot{\psi}\right|_{t=0}\right) \phi
\end{aligned}
$$

We therefore see that the differential of our map is the composition of $\partial_{\mu}$ followed by the adjoint action $9^{9}$ Now, since the metric 25 is invariant under the adjoint action, and since $\partial_{\mu}$ is independent of $\phi$, we see, recalling (11), that the analogue of $J g$ would be a constant. We can therefore discard it.

Since $A_{\mu}=\phi^{-1} \partial_{\mu} \phi$, it follows, as mentioned at the beginning of this subsection, that $A_{\mu}$ is a flat connection. We thus take for $h$ the flatness

[^5]constraint
$$
G_{\mu \nu}(A)=\partial_{\mu} A_{\nu}-\partial_{\nu} A_{\mu}+\left[A_{\mu}, A_{\nu}\right]=0 .
$$

It remains to find $J h$. The derivative of $G$ is easily computed to be

$$
\left.G^{\prime}\right|_{A} v=\partial_{[\nu} v_{\mu]}+\left[A_{[\mu}, v_{\nu]}\right],
$$

where as is traditional, the square brackets on the indices stand for antisymmetrization. Now, note that

$$
\begin{aligned}
\phi^{-1}\left(\partial_{\mu}\left(\phi v_{\nu} \phi^{-1}\right)\right) \phi & =\phi^{-1}\left(\phi\left(\partial_{\mu} v_{\nu}\right) \phi^{-1}+\left(\partial_{\mu} \phi\right) v_{\nu} \phi^{-1}+\phi v_{\nu}\left(\partial_{\mu} \phi^{-1}\right)\right) \phi \\
& =\partial_{\mu} v_{\nu}+\phi^{-1}\left(\partial_{\mu} \phi\right) v_{\nu}-v_{\nu} \phi^{-1}\left(\partial_{\mu} \phi\right)
\end{aligned}
$$

and thus

$$
\left.G^{\prime}\right|_{A} v=\phi^{-1}\left(\partial_{[\mu}\left(\phi v_{\nu]} \phi^{-1}\right)\right) \phi
$$

We thus see that the derivative of $G$ consists of the composition of the adjoint action on the $s u(N)$ one-form, followed by the exterior derivative $d$, followed by an adjoint action on the $s u(N)$ two-form. Clearly, the adjoint action does not affect the metrics (25) and (26). Recalling the definition of Jh given in 15), we see that its analogue in this case coincides with the Jacobian of the exterior derivative. But this is again a constant independent of the field $A$ and thus can be discarded as well.

Putting everything together, we have thus arrived at

$$
\left.\int e^{-\int_{\mathbb{R}^{2}} \operatorname{Tr}\left(\left(\phi^{-1} \partial_{\mu} \phi\right)\left(\phi^{-1} \partial^{\mu} \phi\right)\right.}\right) \mathcal{D} \phi \simeq \int e^{-\int_{\mathbb{R}^{2}} \operatorname{Tr}\left(A_{\mu} A^{\mu}\right)} \delta\left(G_{\mu \nu}(A)\right) \mathcal{D} A_{\mu} .
$$

To reduce clutter, we've only displayed the relation between the two partition functions since the inclusion of sources is trivial.

## Yang-Mills

We are ready to rewrite the Yang-Mills path integral

$$
\int \mathcal{D} A e^{-S_{Y M}(A)},
$$

where, we remind the reader, we are using the radial gauge for our gauge-fixing and thus $\mathcal{D} A$ stands for integration over the set of connections
in this gauge. Note that the Faddeev-Popov determinant in this case is ignorable.

Now, decompose the space of $\mathcal{B}$ 's into the kernel of the map $\mathcal{T}$, denoting it by $\mathcal{B}_{\perp}$, and its orthogona ${ }^{10}$ complement, denoted by $\mathcal{B}_{\|}$. Note that for a connection $A$ in the radial gauge, we have that $\mathcal{T}(\mathcal{B}(A))=A$. It follows from this that the space of connections in the radial gauge can be identified with $\mathcal{B}_{\|}$, and thus, the map $A \rightarrow \mathcal{B}(A)$ is in fact of the form $\mathcal{B}_{\|} \rightarrow\left(\mathcal{B}_{\|}, \mathcal{B}_{\perp}\left(\mathcal{B}_{\|}\right)\right)$. We can therefore use 21) and rewrite the Yang-Mills path integral as

$$
\int \mathcal{D B} e^{-S_{Y M}(\mathcal{B})} \delta(\mathcal{G}(\mathcal{B})) J \mathcal{G}_{\perp}
$$

We have already computed $D \mathcal{G}$ in the subsection dealing with the principal chiral model and have found it to be equal to the composition of an adjoint action with an exterior derivative followed by another adjoint action ${ }^{11}$ Now, note that it follows from (1) that

$$
\mathcal{T}\left(\phi(\sigma) \mathcal{B}(\sigma) \phi^{-1}(\sigma)\right)=\phi\left(\sigma_{x}\right) \mathcal{T}(\mathcal{B})(x)\left(\phi\left(\sigma_{x}\right)\right)^{-1}
$$

where $\phi(\sigma)$ is an $S U(N)$ valued function of the loop $\sigma$. From this equation, it follows at once that the adjoint action takes $\mathcal{B}_{\perp}$ and $\mathcal{B}_{\|}$to themselves. Therefore, an orthonogmal basis of $\mathcal{B}_{\| \mid}$is taken to an orthonormal basis. Combining this fact with the definition (22), we get that $J \mathcal{G}_{\perp}$ is again a constant and thus can be discarded.

Finally, putting together the result of this subsection with that of the previous one ${ }^{12}$, we see that we can write

$$
\int \mathcal{D} A e^{-S_{Y M}(A)} \simeq \int \mathcal{D} \mathcal{B} e^{-S_{Y M}(\mathcal{B})} \delta(\mathcal{G}(\mathcal{B})) \simeq \int \mathcal{D} \phi e^{-S_{Y M}(\phi)}
$$

where $S_{Y M}(A)$ is the usual Yang-Mills action in terms of a connection, $S_{Y M}(\mathcal{B})$ is given by $(2)$ and $S_{Y M}(\phi)$ is the analogue of the principal chiral

[^6]model action given by
$$
S_{Y M}(\phi)=\frac{D}{\mathcal{V}} \int_{(0,1)} d s \int \mathcal{D} \gamma \frac{\operatorname{Tr}\left(\phi^{-1}(\gamma) \frac{\delta \phi(\gamma)}{\delta \gamma^{\mu}(s)} \phi^{-1}(\gamma) \frac{\delta \phi(\gamma)}{\delta \gamma_{\mu}(s)}\right)}{\dot{\gamma}_{\nu}(s) \dot{\gamma}^{\nu}(s)}
$$

Also, $\mathcal{D B}$ stands for integration over the set of all transverse, nonanticipating 1-forms over the space of loops and $\mathcal{D} \phi$ stands for integration over the space of all $S U(N)$ valued maps on the space of loops.

Again, we've only given the equality between the the partition functions since inclusion of sources is straightforward.

We have thus achieved our goal of rewriting the Yang-Mills path integral, and have done so in two ways. Moreover, both of them are of the form of a Gaussian integral satisfying a quadratic constraint.

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[^0]:    ${ }^{1}$ The only attempt to do so seems to be in Appendix A of 8. Our proof is different and we believe is more transparent.

[^1]:    ${ }^{2}$ An objection might be raised here that there is a hidden factor of the $\dot{\tilde{\gamma}}$ hidden inside $\mathcal{F}^{\mu}$. This is indeed the case, but the denominator in $\int_{s_{0}}^{1} d s\left(\frac{1}{\mathcal{V} \dot{\gamma}_{\rho}(s) \dot{\gamma}^{\rho}(s)} \frac{\delta \mathcal{F}^{\mu}\left(\gamma_{s}\right)}{\delta \gamma_{s}^{\mu}(1)}\right)$ has two such factors and thus dominates.

[^2]:    ${ }^{3}$ See e.g. Theorem 3.9 in 14 .
    ${ }^{4}$ There are slightly different conventions for the normalization of the Hausdorff measure in the literature, and formula (9) above holds for only one of them. The exact choice of convention however is not going to matter since below, we will only care about the proportionality of two expressions and not their equality.

[^3]:    ${ }^{5}$ The analogue of 10 also holds but we won't need it.
    ${ }^{6}$ See e.g. Theorem 3.11 in [14].

[^4]:    ${ }^{7} \delta(h)$ is in fact badly defined, since it contains factors of the form $\delta(0)$. This is less of a problem than it seems for two reasons. First, as discussed below, constant factors (even if formally infinite) can often be safely ignored in path integral calculations. Second, the delta functions are formal expressions which are usually implemented through a limiting procedure (or a Fourier transform), and the intermediate expressions are perfectly well-defined with the only problem arising, like here, at the limit (see the discussion about the 'soft' imposition of constraints in [16).
    ${ }^{8}$ Of course, any two things are proportional. The nontrivial statement in our case is that the same constant makes two integrals equal for all $F$ 's.

[^5]:    ${ }^{9}$ Incidentally, this shows that the map $\phi \rightarrow A_{\mu}$ is an immersion, provided one imposes the Dirichlet boundary conditions (i.e. $\phi=$ identity) on the fields before removing the infrared regulator. The need for a careful consideration of the boundary conditions in the regularization is discussed at the beginning of chapter 2 of [13.

[^6]:    ${ }^{10}$ Orthogonality is defined with respect to the metric 25 suitably generalized (we integrate over the space of loops instead of $\mathbb{R}^{2}$ ).
    ${ }^{11}$ Of course, we've only done so for connections on $\mathbb{R}^{2}$, but the formal computations are the same in any dimension and thus, we assume, hold for functionals as well.
    ${ }^{12}$ More precisely, with its generalization to an infinite dimensional setting as opposed to $\mathbb{R}^{2}$.

